## Control Theory and Design:

An RH2 and RH Viewpoint

## by Patrizio Colaneri, Jose C. Geromel, Arturo Locatelli



- ISBN: 0121791904
- Pub. Date: June 1997
- Publisher: Elsevier Science \& Technology Books


## Preface \& Acknowledgments

Robust control theory has been the object of much of the research activity developed in the last fifteen years within the context of linear systems control. At the stage, the results of these efforts constitute a fairly well established part of the scientific community background, so that the relevant techniques can reasonably be exploited for practical purposes. Indeed, despite their complex derivation, these results are of simple implementation and capable of accounting for a number of interesting real life applications. Therefore the demand of including these topics in control engineering courses is both timely and suitable and motivated the birth of this book which covers the basic facts of robust control theory, as well as more recent achievements, such as robust stability and robust performance in presence of parameter uncertainties. The book has been primarily conceived for graduate students as well as for people first entering this research field. However, the particular care which has been dedicated to didactic instances renders the book suited also to undergraduate students who are already acquainted with basic system and control. Indeed, the required mathematical background is supplied where necessary.

Part of the here collected material has been structured according to the textbook Controllo in $\mathrm{RH}_{2}-\mathrm{RH}_{\infty}$ (in Italian) by the authors. They are deeply indebted to the publisher Pitagora for having kindly permitted it. The first five chapters introduces the basic results on $R H_{2}$ and $R H_{\infty}$ theory whereas the last two chapters are devoted to present more recent results on robust control theory in a general and self-contained setting. The authors gratefully acknowledge the financial support of Centro di Teoria dei Sistemi of the Italian National Research Council - CNR, the Brazilian National Research Council - CNPq (under grant 301373/80) and the Research Council of the State of São Paulo, Brazil - FAPESP (under grant 90/3607-0).

This book is a result of a joint, fruitful and equal scientific cooperation. For this reason, the authors' names appear in the front page in alphabetical order.

Milan, Italy
José C. Geromel
Campinas, Brazil
Arturo Locatelli
Milan, Italy

## Table of Contents

Preface \& Acknowledgments
1 Introduction ..... 1
2 Preliminaries ..... 3
3 Feedback Systems Stability ..... 69
4 RH[subscript 2] Control ..... 87
5 RH[actual symbol not reproducible] Control ..... 121
Nonclassical Problems in RH[subscript 2] and RH[actual ..... 195 symbol not reproducible]
7 Uncertain Systems Control Design ..... 263
A Some Facts on Polynomials ..... 301
B Singular Values of Matrices ..... 303
C Riccati Equation ..... 315
D Structural Properties ..... 325
E The Standard 2-Block Scheme ..... 327
F Loop Shifting ..... 337
G Worst Case Analysis ..... 343
H Convex Functions and Sets ..... 349
I Convex Programming Numerical Tools ..... 359
Bibliography ..... 371
Index ..... 375

## Chapter 1

## Introduction

Frequency domain techniques have longly being proved to be particularly fruitful and simple in the design of (linear time invariant) SISO ${ }^{1}$ control systems. Less appealing have appeared for many years the attempts of generalizing such nice techniques to the MIMO ${ }^{2}$ context. This partially motivated the great deal of interest which has been devoted to time domain design methodologies starting in the early 60 's. Indeed, this stream of research originated a huge number of results both of remarkable conceptual relevance and practical impact, the most celebrated of which is probably the LQG ${ }^{3}$ design. Widely acknowledged are the merits of such an approach: among them the relatively small computational burden involved in the actual definition of the controller and the possibility of affecting the dynamical behavior of the control system through a guided sequence of experiments aimed at the proper choice of the parameters of both the performance index (weighting matrices) and uncertainty description (noises intensities). Equally well known are the limits of the LQG design methodology, the most significant of which is the possible performance decay caused by operative conditions even slightly differing from the (nominal) ones referred to in the design stage. Specifically, the lack of robustness of the classical LQG design originates from the fact that it does not account for the uncertain knowledge or unexpected perturbations of the plant, actuators and sensors parameters.

The need of simultaneously complying with design requirements naturally specified in the frequency domain and guaranteeing robustness of the control system in the face of uncertainties and/or parameters deviations, focused much of the research activity on the attempt of overcoming the traditional and myopic dichotomy between time and frequency domain approaches. At the stage, after about two decades of intense efforts on these lines, the control system designer can rely on a set of well established results which give proper answers to the significant instances of performance and stability robustness. The value of the results achieved so far partially stems in the construction of a unique formal theoretical picture which naturally includes both the classical LQG design ( $\mathrm{RH} \mathrm{H}_{2}$ design) , revisited at the light of a transfer function-like approach, and the new challenging developments of the so called robust design ( $R H_{\infty}$ design), which encompasses most of the above mentioned robustness instances.

The design methodologies which are presented in the book are based on the minimization of a performance index, simply consisting of the norm of a suitable transfer

[^0]function. A distinctive feature of these techniques is the fact that they do not come up with a unique solution to the design problem; rather, they provide a whole set of (admissible) solutions which satisfy a constraint on the maximum deterioration of the performance index. The attitude of focusing on the class of admissible controllers instead of determining just one of them can be traced back to a fundamental result which concerns the parametrization of the class of controllers stabilizing a given plant. Chapter 3 is actually dedicated to such a result and deals also with other questions on feedback systems stability. In subsequent Chapters 4 and 5 the main results of $R H_{2}$ and $R H_{\infty}$ design are presented, respectively. In addition, a few distinguishing aspects of the underlying theory are emphasized as well, together with particular, yet significant, cases of the general problem. Chapter 5 contains also a preliminary discussion on the robustness requirements which motivate the formulation of the so called standard $R H_{\infty}$ control problem. Chapter 6 and 7 go beyond the previous ones in the sense that the design problems to be dealt with are setting in a more general framework. One of the most interesting examples of this situation is the so called mixed $R H_{2} / R H_{\infty}$ problem which is expressed in terms of both $R H_{2}$ and $R H_{\infty}$ norms of two transfer functions competing with each other to get the best tradeoff between performance and robustness. Other problems that fall into this framework are those related to regional pole placement, time-domain specification and structural constraints. All of them share basically the same difficulty to be faced numerically. Indeed, they can not be solved by the methodology given in the previous Chapters but by means of mathematical programming methods. More specifically, all can (after a proper change of variables) be converted into convex problems. This feature is important in both practical and theoretical points of view since numerical efficiency allows the treatment of real-word problems of generally large dimension while global optimality is always assured. Chapter 7 is devoted to the controllers design for systems subject to structured convex bounded uncertainties which models in an adequate and precise way many classes of parametric uncertainties with practical appealing. The associated optimal control problems are formulated and solved jointly with respect to the controller transfer function and the feasible uncertainty in order to guarantee minimum loss in the performance index. One of such situation of great importance for its own is the design problem involving actuators failure. Robust stability and performance are addressed for two classes of nonlinear perturbations, leading to what are called Persidiskii and Lur'e design. In general terms, the same technique involving the reduction of the related optimal control design problems to convex programming problems is again used. The main point to be remarked is that the two classes of nonlinear perturbations considered impose additional linear and hence convex constraints, to the matrices variables to be determined.

Treating these arguments requires a fairly deep understanding of some facts from mathematics not so frequently included in the curricula of students in Engineering. Covering the relevant mathematical background is the scope of Chapter 2, where the functional (Hardy) spaces which permeate all over the book are characterized. Some miscellaneous facts on matrix algebra, system and control theory and convex optimization are collected in Appendix A through I.

## Chapter 2

## Preliminaries

### 2.1 Introduction

The scope of this chapter is twofold: on one hand it is aimed at presenting the extension of the concepts of poles and zeros, well known for single-input single-output (SISO) systems, to the multivariable case; on the other, it is devoted to the introduction of the basic notions relative to some functional spaces whose elements are matrices of rational functions (spaces $R L_{2}, R L_{\infty}, R H_{2}, R H_{\infty}$ ). The reason of this choice stems from the need of presenting a number of results concerning significant control problems for linear, continuous-time, finite dimensional and time-invariant systems.

The derivation of the related results takes substantial advantage on the nature of the analysis and design methodology adopted; such a methodology was actually developed so as to take into account state-space and frequency based techniques at the same time.

For this reason, it should not be surprising the need of carefully extending to multiinput multi-output (MIMO) systems the notions of zeros and poles, which proved so fruitful in the context of SISO systems. In Section 2.5, where this attempt is made, it will be put into sharp relief few fascinating and in some sense unexpected relations between poles, zeros, eigenvalues, time responses and ranks of polynomial matrices.

Analogously, it should be taken for granted the opportunity of going in depth on the characterization of transfer matrices (transfer functions for MIMO systems) in their natural embedding, namely, in the complex plane. The systems considered hereafter obviously have rational transfer functions. This leads to the need of providing, in Section 2.8 the basic ideas on suitable functional spaces and linear operators so as to throw some light on the connections between facts which naturally lie in time-domain with others more suited with the frequency-domain setting.

Although the presentation of these two issues is intentionally limited to few basic aspects, nevertheless it requires some knowledge on matrices of polynomials, matrices of rational functions, singular values and linear operators. To the acquisition of such notions are dedicated Sections 2.3-2.7.

### 2.2 Notation and terminology

The continuous-time linear time-invariant dynamic systems, object of the present text, are described, depending on circumstances, by a state space representation

$$
\begin{aligned}
\dot{x} & =A x+B u \\
y & =C x+D u
\end{aligned}
$$

or by their transfer function

$$
G(s)=C(s I-A)^{-1} B+D
$$

The signals which refer to a system are indifferently intended to be in time-domain or in frequency-domain all the times the context does not lead to possible misunderstandings. Sometimes, it is necessary to explicitly stress that the derivation is in frequency-domain. In this case, the subscript "L" indicates the Laplace transform of the considered signal, whereas the subscript "L0" denotes the Laplace transform when the system state at the initial time is zero (typically, this situation occurs when one thinks in terms of transfer functions). For instance, with reference to the above system, one may write

$$
\begin{aligned}
y_{L 0} & =G(s) u_{L} \\
y_{L} & =y_{L 0}+C(s I-A)^{-1} x(0)
\end{aligned}
$$

Occasionally, the transfer function $G(s)$ of a system $\Sigma$ is explicitly related to one of its realizations by writing

$$
G(s)=\Sigma(A, B, C, D)
$$

or

$$
G(s)=\left[\begin{array}{l|l}
A & B \\
\hline C & D
\end{array}\right]
$$

The former notation basically has a compactness value, whereas the latter is mainly useful when one wants to display possible partitions in the input and/or output matrices. For example, the system

$$
\begin{aligned}
\dot{x} & =A x+B_{1} w+B_{2} u \\
z & =C_{1} x+D_{12} u \\
y & =C_{2} x+D_{21} w
\end{aligned}
$$

is related to its transfer function $G(s)$ by writing

$$
G(s)=\left[\begin{array}{c|cc}
A & B_{1} & B_{2} \\
\hline C_{1} & 0 & D_{12} \\
C_{2} & D_{21} & 0
\end{array}\right]
$$

When a purely algebraic (i.e. nondynamic) system is considered, these notations become

$$
G(s)=\Sigma(\emptyset, \emptyset, \emptyset, D)=\left[\begin{array}{l|l}
\emptyset & \emptyset \\
\hline \emptyset & D
\end{array}\right]
$$

Referring to the class of systems considered here, the transfer functions are in fact rational matrices of complex variable, namely, matrices whose generic element is a rational function, i.e., a ratio of polynomials with real coefficients. The transfer function is said to be proper when each element is a proper rational function, i.e., a ratio of polynomials with the degree of the numerator not greater than the degree of the denominator. When this inequality holds in a strict sense for each element of the matrix, the transfer function is said to be strictly proper. Briefly, $G(s)$ is proper if

$$
\lim _{s \rightarrow \infty} G(s)=K<\infty
$$

where the notation $K<\infty$ means that each element of matrix $K$ is finite. Analogously, $G(s)$ is strictly proper if

$$
\lim _{s \rightarrow \infty} G(s)=0
$$

A rational matrix $G(s)$ is said to be analytic in $\operatorname{Re}(s) \geq 0$ (resp. $\leq 0$ ) if all the elements of the matrix are bounded functions in the closed right (resp. left) half plane.

In connection with a system characterized by the transfer function

$$
G(s)=\left[\begin{array}{c|c}
A & B  \tag{2.1}\\
\hline C & D
\end{array}\right]
$$

the so-called adjoint system has transfer function

$$
G^{\sim}(s):=G^{\prime}(-s)=\left[\begin{array}{c|c}
-A^{\prime} & -C^{\prime} \\
\hline B^{\prime} & D^{\prime}
\end{array}\right]
$$

whereas the transfer function of the so-called transpose system is

$$
G^{\prime}(s):=\left[\begin{array}{c|c}
A^{\prime} & C^{\prime} \\
\hline B^{\prime} & D^{\prime}
\end{array}\right]
$$

System (2.1) is said to be input-output stable if its transfer function $G(s)$ is analytic in $\operatorname{Re}(s) \geq 0(G(s)$ is stable, by short). It is said to be internally stable if matrix $A$ is stable, i.e., if all its eigenvalues have negative real parts.

Now observe that a system is input-output stable if and only if all elements of $G(s)$, whenever expressed as ratio of polynomials without common roots, have their poles in the open left half plane only. If the realization of system (2.1) is minimal, the system is input-output stable if and only if it is internally stable.

Finally, the conjugate transpose of the generic (complex) matrix $A$ is denoted by $A^{\sim}$ and, if it is square, $\lambda_{i}(A)$ is its $i$-th eigenvalue, while

$$
r_{s}(A):=\max _{i}\left|\lambda_{i}(A)\right|
$$

denotes its spectral radius.

### 2.3 Polynomial matrices

A polynomial matrix is a matrix whose elements are polynomials in a unique unknown. Throughout the book, such an unknown is denoted by the letter $s$. All the polynomial
coefficients are real. Hence, the element $n_{i j}(s)$ in position $(i, j)$ in the polynomial matrix $N(s)$ takes the form

$$
n_{i j}(s)=\alpha_{\nu} s^{\nu}+\alpha_{\nu-1} s^{\nu-1}+\cdots+\alpha_{1} s+\alpha_{0}, \alpha_{k} \in R, \quad \forall k
$$

The degree of a polynomial $p(s)$ is denoted by $\operatorname{deg}[p(s)]$. If the leading coefficient $\alpha_{\nu}$ is equal to one, the polynomial is said to be monic.

The rank of a polynomial matrix $N(s)$, denoted by $\operatorname{rank}[N(s)]$, is defined by analogy from the definition of the rank of a numeric matrix, i.e., it is the dimension of the largest square matrix which can be extracted from $N(s)$ with determinant not identically zero.

A square polynomial matrix is said to be unimodular if it has full rank (it is invertible) and its determinant is constant.

Example 2.1 The polynomial matrices

$$
N_{1}(s)=\left[\begin{array}{cc}
1 & s+1 \\
0 & 3
\end{array}\right], \quad N_{2}(s)=\left[\begin{array}{cc}
s+1 & s-2 \\
s+2 & s-1
\end{array}\right]
$$

are unimodular since $\operatorname{det}\left[N_{1}(s)\right]=\operatorname{det}\left[N_{2}(s)\right]=3$.
A very peculiar property of a unimodular matrix is that its inverse is still a polynomial (and obviously unimodular) matrix. Not differently from what is usually done for polynomials, the polynomial matrices can be given the concepts of divisor and greatest common divisor as well.

Definition 2.1 (Right divisor) Let $N(s)$ be a polynomial matrix. A square polynomial matrix $R(s)$ is said to be a right divisor of $N(s)$ if it is such that

$$
N(s)=\bar{N}(s) R(s)
$$

with $\bar{N}(s)$ a suitable polynomial matrix.
An analogous definition can be formulated for the left divisor.
Definition 2.2 (Greatest common right divisor) Let $N(s)$ and $D(s)$ be polynomial matrices with the same number of columns. A square polynomial matrix $R(s)$ is said to be a Greatest Common Right Divisor (GCRD) of $(N(s), D(s))$ if it is such that
i) $R(s)$ is a right divisor of $D(s)$ and $N(s)$, i.e.

$$
\begin{aligned}
& N(s)=\bar{N}(s) R(s) \\
& D(s)=\bar{D}(s) R(s)
\end{aligned}
$$

with $\bar{N}(s)$ and $\bar{D}(s)$ suitable polynomial matrices
ii) For each polynomial matrix $\hat{R}(s)$ such that

$$
\begin{aligned}
& N(s)=\hat{N}(s) \hat{R}(s) \\
& D(s)=\hat{D}(s) \hat{R}(s)
\end{aligned}
$$

with $\hat{N}(s)$ and $\hat{D}(s)$ polynomial matrices, it turns out that $R(s)=W(s) \hat{R}(s)$ where $W(s)$ is again a suitable polynomial matrix.

A similar definition can be formulated for the Greatest Common Left Divisor (GCLD).
It is easy to see, by exploiting the properties of unimodular matrices, that, given two polynomial matrices $N(s)$ and $D(s)$, there exist infinite GCRD's (and obviously GCLD's). A way to compute a GCRD (resp. GCLD) of two assigned polynomial matrices $N(s)$ and $D(s)$ relies on their manipulation through a unimodular matrix which represents a sequence of suitable elementary operations on their rows (resp. columns). The elementary operations on the rows (resp. columns) of a polynomial matrix $N(s)$ are

1) Interchange of the i-th row (resp. i-th column) with the j-th row (resp. j-th column)
2) Multiplication of the i-th row (resp. i-th column) by a nonzero scalar
3) Addition of a polynomial multiple of the $i$-th row (resp. i-th column) to the $j$-th row (resp. j-th column).

It is readily seen that each elementary operation can be performed premultiplying (resp. postmultiplying) $N(s)$ by a suitable polynomial and unimodular matrix $T(s)$. Moreover, matrix $T(s) N(s)$ (resp. $N(s) T(s)$ ) turns out to have the same rank as $N(s)$.
Remark 2.1 Notice that, given two polynomials $r_{0}(s)$ and $r_{1}(s)$ with $\operatorname{deg}\left[r_{0}(s)\right] \geq \operatorname{deg}\left[r_{1}(s)\right]$, it is always possible to define two sequences of polynomials $\left\{r_{i}(s), i=2,3, \cdots, p+2\right\}$ and $\left\{q_{i}(s), i=1,2, \cdots, p+1\right\}$, with $0 \leq p \leq \operatorname{deg}\left[r_{1}(s)\right]$, such that

$$
\begin{aligned}
& r_{i}(s)=q_{i+1}(s) r_{i+1}(s)+r_{i+2}(s), \quad i=0,1, \cdots, p \\
& \operatorname{deg}\left[r_{i+2}(s)\right]<\operatorname{deg}\left[r_{i+1}(s)\right], \quad i=0,1, \cdots, p \\
& r_{p+2}(s)=0
\end{aligned}
$$

Letting

$$
\begin{aligned}
& T_{i}(s):=\left[\begin{array}{cc}
1 & -q_{i}(s) \\
0 & 1
\end{array}\right], \quad n_{i}(s):=\left[\begin{array}{c}
r_{i-1}(s) \\
r_{i}(s)
\end{array}\right], i=1,3,5, \cdots \\
& T_{i}(s):=\left[\begin{array}{cc}
1 & 0 \\
-q_{i}(s) & 1
\end{array}\right], \quad n_{i}(s):=\left[\begin{array}{c}
r_{i}(s) \\
r_{i-1}(s)
\end{array}\right], i=2,4,6, \cdots \\
& T(s):=\prod_{i=0}^{p} T_{p+1-i}(s)
\end{aligned}
$$

and noticing that $T(s)$ is unimodular (product of unimodular matrices), it turns out that

$$
\begin{aligned}
& T(s) n_{1}(s)=\left[\begin{array}{c}
r_{p+1}(s) \\
0
\end{array}\right], p=1,3,5, \cdots \\
& T(s) n_{1}(s)=\left[\begin{array}{c}
0 \\
r_{p+1}(s)
\end{array}\right], p=2,4,6, \cdots
\end{aligned}
$$

For instance, take $r_{0}(s)=s^{4}+2 s^{2}-s+2, r_{1}(s)=s^{3}+s-2$. It follows that $q_{1}(s)=s$, $q_{2}(s)=s-1, r_{2}(s)=s^{2}+s+2$ and $r_{3}(s)=0$.

By repeatedly exploiting the facts shown in Remark 2.1, it is easy to verify that, given a polynomial matrix $N(s)$ with the number of rows not smaller than the number of columns, there exists a suitable polynomial and unimodular matrix $T(s)$ such that

$$
T(s) N(s)=\left[\begin{array}{c}
R(s) \\
0
\end{array}\right]
$$

where $R(s)$ is a square polynomial matrix.
Algorithm 2.1 (GCRD of two polynomial matrices) Let $N(s)$ and $D(s)$ be two polynomial matrices with the same number, say $m$, of columns and with $n_{n}$ and $n_{d}$ rows, respectively.

1) Assume that $m \leq n_{d}+n_{n}$, otherwise go to point 4). Let $P(s):=\left[D^{\prime}(s) N^{\prime}(s)\right]^{\prime}$ and determine a polynomial and unimodular matrix $T(s)$ such that

$$
\left.T(s) P(s)=\left[\begin{array}{c}
R(s) \\
0
\end{array}\right]\right\} m \mathrm{rows}
$$

Notice that $T(s)$ can be partitioned as follows

$$
\left.T(s):=\left[\begin{array}{ll}
\underbrace{T_{d 1}(s)}_{n_{d} \text { columns }} & T_{n 1}(s) \\
T_{d 2}(s) & T_{n 2}(s)
\end{array}\right]\right\} m \text { rows }
$$

2) Letting $S(s):=T^{-1}(s)$ and writing

$$
\left.S(s):=\left[\begin{array}{ll}
\underbrace{S_{d 1}(s)}_{n_{d} \text { columns }} & S_{d 2}(s) \\
S_{n 1}(s) & S_{n 2}(s)
\end{array}\right]\right\} m \text { rows }
$$

it turns out

$$
\begin{aligned}
& D(s)=S_{d 1}(s) R(s) \\
& N(s)=S_{n 1}(s) R(s)
\end{aligned}
$$

so that $R(s)$ is a right divisor of both $D(s)$ and $N(s)$.
3) It also holds that

$$
\begin{equation*}
R(s)=T_{d 1}(s) D(s)+T_{n 1}(s) N(s) \tag{2.2}
\end{equation*}
$$

Hence, suppose that $\hat{R}(s)$ is any other right divisor of both $D(s)$ and $N(s)$. Therefore, for some polynomial matrices $\bar{D}(s)$ and $\bar{N}(s)$ it follows that $D(s)=$ $\bar{D}(s) \hat{R}(s)$ and $N(s)=\bar{N}(s) \hat{R}(s)$. The substitution of these two expressions in eq. (2.2) leads to $R(s)=\left[T_{d 1}(s) \bar{D}(s)+T_{n 1}(s) \bar{N}(s)\right] \hat{R}(s)$ so that $R(s)$ is a GCRD of $(N(s), D(s))$.
4) If $m>n_{d}+n_{n}$, take two matrices $\hat{D}(s)$ and $\hat{N}(s)$ both with $m$ columns and $n_{d}$ and $n_{n}$ rows, respectively

$$
\begin{aligned}
& \hat{D}(s):=\left[\begin{array}{lll}
I & 0 & 0
\end{array}\right] \\
& \hat{N}(s):=\left[\begin{array}{lll}
I & 0 & 0
\end{array}\right]
\end{aligned}
$$

and let

$$
\left.R(s):=\left[\begin{array}{c}
D(s)  \tag{2.3}\\
N(s) \\
0
\end{array}\right]\right\} m-n_{d}-n_{n} \text { rows }
$$

Thus, $D(s)=\hat{D}(s) R(s)$ and $N(s)=\hat{N}(s) R(s)$. Hence, $R(s)$ is a right divisor of both $D(s)$ and $N(s)$. Assume now that $\hat{R}(s)$ is any other right divisor, i.e. there exist two polynomial matrices $\bar{D}(s)$ and $\bar{N}(s)$ such that $D(s)=\bar{D}(s) \hat{R}(s)$ and $N(s)=\bar{N}(s) \hat{R}(s)$. By substituting these two last expressions in eq. (2.3) one obtains

$$
R(s):=\left[\begin{array}{c}
\bar{D}(s) \\
\bar{N}(s) \\
0
\end{array}\right] \hat{R}(s)
$$

so leading to the conclusion that $R(s)$ is a GCRD of $(N(s), D(s))$.

Example 2.2 Consider the matrices

$$
\begin{aligned}
& D(s)=\left[\begin{array}{cc}
s^{2}-s & s^{2} \\
2 s^{2}+9 s+5 & 2 s^{2}+5 s+5
\end{array}\right] \\
& N(s)=\left[\begin{array}{ll}
s^{2}+1 & s^{2}+2 s+1
\end{array}\right]
\end{aligned}
$$

Now take

$$
\begin{array}{lll}
T_{1}(s)=\left[\begin{array}{ccc}
1 & 0 & 0 \\
-2 & 1 & 0 \\
-1 & 0 & 1
\end{array}\right] & & T_{2}(s)=\left[\begin{array}{ccc}
1 & 0 & -s \\
0 & 1 & -11 \\
0 & 0 & 1
\end{array}\right] \\
T_{3}(s)=\left[\begin{array}{ccc}
1 & -s / 3 & 0 \\
0 & 1 & 0 \\
0 & s / 6 & 1
\end{array}\right] & , T_{4}(s)=\left[\begin{array}{lll}
1 & 0 & 0 \\
0 & 1 & 6 \\
0 & 0 & 1
\end{array}\right] \\
T_{5}(s)=\left[\begin{array}{ccc}
0 & 0 & 1 \\
0 & -1 & 0 \\
1 & 0 & 0
\end{array}\right] & , T_{6}(s)=\left[\begin{array}{ccc}
1 & 0 & 0 \\
0 & 1 & -51 / 14 \\
0 & 0 & 1
\end{array}\right] \\
T_{7}(s)=\left[\begin{array}{ccc}
1 & 0 & 0 \\
0 & 14 / 103 & 0 \\
0 & -196 s / 309 & 1
\end{array}\right] & T_{8}(s)=\left[\begin{array}{ccc}
1 & 0 & 0 \\
0 & 1 & 0 \\
0 & -1 & 1
\end{array}\right]
\end{array}
$$

Then

$$
\begin{gathered}
T(s)=\prod_{i=0}^{7} T_{8-i}(s)= \\
=\left[\begin{array}{ccc}
(3 s-2) / 2 & s / 6 & (6-11 s) / 6 \\
-(24 s+93) / 103 & (3 s-14) / 103 & (18 s+70) / 103 \\
\left(112 s^{2}+252 s+196\right) / 103 & \left(-14 s^{2}+28 s+14\right) / 103 & -\left(84 s^{2}+70 s+70\right) / 103
\end{array}\right]
\end{gathered}
$$

so that

$$
T(s) P(s):=T(s)\left[\begin{array}{c}
D(s) \\
N(s)
\end{array}\right]=\left[\begin{array}{cc}
1 & \left(-17 s^{2}+6 s+6\right) / 6 \\
0 & s \\
0 & 0
\end{array}\right]
$$

and

$$
R(s)=\left[\begin{array}{cc}
1 & \left(-17 s^{2}+6 s+6\right) / 6 \\
0 & s
\end{array}\right]
$$

Finally, notice that

$$
T^{-1}(s):=S(s)=\left[\begin{array}{ll}
S_{d 1(s)} & S_{d 2}(s) \\
S_{n 1}(s) & S_{n 2}(s)
\end{array}\right]
$$

with

$$
\begin{aligned}
& S_{d 1}(s)=\left[\begin{array}{cc}
s(s-1) & \left(17 s^{3}-23 s^{2}+6 s+6\right) / 6 \\
2 s^{2}+9 s+5 & \left(34 s^{3}+141 s^{2}+31 s-54\right) / 6
\end{array}\right] \\
& S_{n 1}(s)=\left[\begin{array}{ll}
s^{2}+1 & \left(17 s^{3}-6 s^{2}+17 s+6\right) / 6
\end{array}\right]
\end{aligned}
$$

It is then easy to verify that $D(s)=S_{d 1}(s) R(s)$ and that $N(s)=S_{n 1}(s) R(s)$.
The familiar concept of coprimeness, easily introduced for polynomials, can be properly extended to polynomial matrices as follows.

Definition 2.3 (Right coprimeness) Two polynomial matrices $N(s)$ and $D(s)$ having the same number of columns, are said to be right coprime if the two equations

$$
\begin{aligned}
& N(s)=\bar{N}(s) T(s) \\
& D(s)=\bar{D}(s) T(s)
\end{aligned}
$$

where $\bar{N}(s)$ and $\bar{D}(s)$ are suitable polynomial matrices, are verified by a unimodular polynomial matrix $T(s)$ only.

Example 2.3 The matrices

$$
N(s)=\left[\begin{array}{c}
s^{2}-1 \\
s^{3}+3 s^{2}-s-3
\end{array}\right], \quad D(s)=\left[\begin{array}{c}
2 s^{2}+6 s+4 \\
s^{3}+3 s^{2}+2 s
\end{array}\right]
$$

are not right coprime. Actually, it turns out that

$$
\begin{aligned}
& N(s)=\bar{N}(s) R(s), \quad \bar{N}(s)=\left[\begin{array}{c}
s-1 \\
s^{2}+2 s-3
\end{array}\right], \quad R(s)=s+1 \\
& D(s)=\bar{D}(s) R(s), \quad \bar{D}(s)=\left[\begin{array}{c}
2 s+4 \\
s^{2}+2 s
\end{array}\right]
\end{aligned}
$$

and $R(s)$ is not unimodular $(\operatorname{det}[R(s)]=s+1)$.
Of course, an analogous definition can be stated for the left coprimeness. Definitions 2.1-2.3 also yield that two matrices are right (resp. left) coprime if all their common right (resp. left) divisors are actually unimodular. In particular, each GCRD (resp. GCLD) of two right (resp. left) coprime matrices must be unimodular. In view of Algorithm 2.1, this entails that a possible way to verify whether or not two matrices are right (resp. left) coprime, is computing and evaluating the determinant of a greatest common divisor. As a matter of fact, if a GCRD (resp. GCLD) is unimodular, then all other greatest common divisors are unimodular as well. More precisely, if $R_{1}(s)$ and $R_{2}(s)$ are two GCRD's and $R_{1}(s)$ is unimodular, it results $R_{1}(s)=$ $W(s) R_{2}(s)$, with $W(s)$ polynomial. Since $\operatorname{det}\left[R_{1}(s)\right] \neq 0$ it follows that $\operatorname{det}\left[R_{2}(s)\right]$ $\neq 0$ as well.

Again from Algorithm 2.1 (step 1) it can be concluded that two polynomial matrices $D(s)$ and $N(s)$ with the same number of columns, say $m$, are right coprime if the rank of $P(s):=\left[D^{\prime}(s) N^{\prime}(s)\right]^{\prime}$ is $m$ for any $s$. As a matter of fact, coprimeness is equivalent to $R(s)$ being unimodular so that $\operatorname{rank}[T(s) P(s)]$ must be constant and equal to $m$. Since $T(s)$ is unimodular, $\operatorname{rank}[P(s)]$ must be constant and equal to $m$ as well.

Lemma 2.1 Let $N(s)$ and $D(s)$ be two polynomial matrices with the same number of columns and let $R(s)$ be a GCRD of $(N(s), D(s))$. Then,
i) $T(s) R(s)$ is a $G C R D$ of $(N(s), D(s))$ for any polynomial and unimodular matrix $T(s)$
ii) If $\bar{R}(s)$ is an arbitrary $G C R D$ of $(N(s), D(s))$, then there exists a polynomial and unimodular matrix $T(s)$ such that

$$
\bar{R}(s)=T(s) R(s)
$$

Proof Point $i$ ) Being $R(s)$ a GCRD of $(N(s), D(s))$ it follows that

$$
\begin{array}{lll}
N(s)=\widetilde{N}(s) R(s), & D(s)=\widetilde{D}(s) R(s) \\
N(s)=\bar{N}(s) \bar{R}(s), & D(s)=\widetilde{D}(s) \bar{R}(s) \\
R(s)=W(s) \bar{R}(s) &
\end{array}
$$

where $\widetilde{N}(s), \widetilde{D}(s), \bar{N}(s), \bar{D}(s)$ and $W(s)$ are suitable polynomial matrices. Taken an arbitrary polynomial and unimodular matrix $T(s)$, let $\hat{R}(s):=T(s) R(s)$. It follows that $N(s)=\hat{N}(s) \hat{R}(s), \hat{N}(s)=\widetilde{N}(s) T^{-1}(s), D(s)=\hat{D}(s) \hat{R}(s), \hat{D}(s)=\widetilde{D}(s) T^{-1}(s)$. Furthermore, it is $\hat{R}(s)=\hat{W}(s) \bar{R}(s), \hat{W}(s)=T(s) W(s)$. Hence, $\hat{R}(s)$ is a GCRD of $(N(s), D(s))$ as well.

Point ii) If $R(s)$ and $\bar{R}(s)$ are two GCRD's of $(N(s), D(s))$, then, for two suitable polynomial matrices $W(s)$ and $\bar{W}(s)$, it results $R(s)=W(s) \bar{R}(s)$ and $\bar{R}(s)=$ $\bar{W}(s) R(s)$. From these relations it follows that $\operatorname{rank}[\bar{R}(s)] \leq \operatorname{rank}[R(s)] \leq \operatorname{rank}[\bar{R}(s)]$. Therefore, $\operatorname{rank}[R(s)]=\operatorname{rank}[\bar{R}(s)]$. Let now $U(s)$ and $\bar{U}(s)$ be two polynomial and unimodular matrices such that

$$
U(s) R(s)=\left[\begin{array}{c}
H(s) \\
0
\end{array}\right], \quad \bar{U}(s) \bar{R}(s)=\left[\begin{array}{c}
\bar{H}(s) \\
0
\end{array}\right]
$$

where the two submatrices $H(s)$ and $\bar{H}(s)$ have the same number of rows equal to $\operatorname{rank}[R(s)]$. Consequently,

$$
\begin{align*}
{\left[\begin{array}{c}
H(s) \\
0
\end{array}\right] } & =U(s) R(s) \\
& =U(s) W(s) \bar{R}(s) \\
& =U(s) W(s) \bar{U}^{-1}(s)\left[\begin{array}{c}
\bar{H}(s) \\
0
\end{array}\right] \\
& =\Gamma(s)\left[\begin{array}{c}
\bar{H}(s) \\
0
\end{array}\right] \\
& =\left[\begin{array}{ll}
\Gamma_{11}(s) & \Gamma_{12}(s) \\
\Gamma_{21}(s) & \Gamma_{22}(s)
\end{array}\right]\left[\begin{array}{c}
\bar{H}(s) \\
0
\end{array}\right] \tag{2.4}
\end{align*}
$$

and

$$
\begin{aligned}
{\left[\begin{array}{c}
\bar{H}(s) \\
0
\end{array}\right] } & =\bar{U}(s) \bar{R}(s) \\
& =\bar{U}(s) \bar{W}(s) R(s) \\
& =\bar{U}(s) \bar{W}(s) U^{-1}(s)\left[\begin{array}{c}
H(s) \\
0
\end{array}\right]
\end{aligned}
$$

$$
\begin{align*}
& =\bar{\Gamma}(s)\left[\begin{array}{c}
H(s) \\
0
\end{array}\right] \\
& =\left[\begin{array}{ll}
\bar{\Gamma}_{11}(s) & \bar{\Gamma}_{12}(s) \\
\bar{\Gamma}_{21}(s) & \bar{\Gamma}_{22}(s)
\end{array}\right]\left[\begin{array}{c}
H(s) \\
0
\end{array}\right] \tag{2.5}
\end{align*}
$$

From eq.(2.4), (2.5) it follows that

$$
\begin{align*}
& 0=\Gamma_{21}(s) \bar{H}(s)  \tag{2.6}\\
& 0=\bar{\Gamma}_{21}(s) H(s) \tag{2.7}
\end{align*}
$$

Being $R(s)$ and $\bar{R}(s)$ square, matrices $H(s)$ and $\bar{H}(s)$ have ranks equal to the number of their rows, which is obviously not greater than the number of their columns. Therefore, from eqs. (2.6) and (2.7) it follows that $\Gamma_{21}(s)=\bar{\Gamma}_{21}(s)$. Equations (2.4), (2.5) outline that matrices $\Gamma_{12}(s), \Gamma_{22}(s), \bar{\Gamma}_{12}(s)$ and $\bar{\Gamma}_{21}(s)$ are in fact arbitrary. Hence, one can set $\Gamma_{12}(s)=\bar{\Gamma}_{12}(s)=0$ and $\Gamma_{22}(s)=\bar{\Gamma}_{22}(s)=I$. Based on these considerations, one can henceforth assume that $\Gamma(s)$ and $\bar{\Gamma}(s)$ have the form

$$
\Gamma(s)=\left[\begin{array}{cc}
\Gamma_{11} & 0 \\
0 & I
\end{array}\right], \quad \bar{\Gamma}(s)=\left[\begin{array}{cc}
\bar{\Gamma}_{11} & 0 \\
0 & I
\end{array}\right]
$$

so that, from eqs. (2.4) and (2.5) it follows

$$
\begin{aligned}
{\left[\begin{array}{c}
H(s) \\
0
\end{array}\right] } & =\left[\begin{array}{cc}
\Gamma_{11} & 0 \\
0 & I
\end{array}\right]\left[\begin{array}{c}
\bar{H}(s) \\
0
\end{array}\right] \\
& =\left[\begin{array}{cc}
\Gamma_{11} & 0 \\
0 & I
\end{array}\right]\left[\begin{array}{cc}
\bar{\Gamma}_{11} & 0 \\
0 & I
\end{array}\right]\left[\begin{array}{c}
H(s) \\
0
\end{array}\right]
\end{aligned}
$$

In particular, it is $H(s)=\Gamma_{11}(s) \bar{\Gamma}_{11}(s) H(s)$, so that, recalling the properties of $H(s)$, it results $I=\Gamma_{11}(s) \bar{\Gamma}_{11}(s)$. Hence, both $\Gamma_{11}(s)$ and $\bar{\Gamma}_{11}(s)$ are unimodular, since their inverses are still polynomial matrices. The same holds for $\Gamma(s)$ and $\bar{\Gamma}(s)$ as well. Finally,

$$
\begin{aligned}
\bar{R}(s) & =\bar{U}^{-1}(s)\left[\begin{array}{c}
\bar{H}(s) \\
0
\end{array}\right] \\
& =\bar{U}^{-1}(s) \bar{\Gamma}(s)\left[\begin{array}{c}
H(s) \\
0
\end{array}\right] \\
& =\bar{U}^{-1}(s) \bar{\Gamma}(s) U(s) R(s) \\
& :=T(s) R(s)
\end{aligned}
$$

and $T(s)$ is actually unimodular since it is the product of unimodular matrices.
Remark 2.2 In view of the results now proved and the given definitions, it is apparent that when the matrices are in fact scalars it results: (i) a right divisor is also a left divisor and vice-versa; (ii) two GCRD's differ only for a multiplicative scalar, since all unimodular polynomials $p(s)$ take the form $p(s):=\alpha \in R, \alpha \neq 0$; (iii) two polynomials are coprime if and only if they do not have common roots.

Right coprime polynomial matrices enjoy the property stated in the following lemma, which provides the generalization of a well known result relative to integers and polynomials (see also Theorem A.1).

Lemma 2.2 Let $N(s)$ and $D(s)$ be two polynomial matrices with the same number of columns. Then, they are right coprime if and only if there exist two polynomial matrices $X(s)$ and $Y(s)$ such that

$$
\begin{equation*}
X(s) N(s)+Y(s) D(s)=I \tag{2.8}
\end{equation*}
$$

Proof Based on the results illustrated in Algorithm 2.1, it is always possible to write a generic GCDR $R(s)$ of $(N(s), D(s))$ as $R(s)=\hat{X}(s) N(s)+\hat{Y}(s) D(s)$. Moreover, if $N(s)$ and $D(s)$ are coprime, $R(s)$ must be unimodular so that

$$
\begin{aligned}
I & =R^{-1}(s) R(s) \\
& =R^{-1}(s)[\hat{X}(s) N(s)+\hat{Y}(s) D(s)] \\
& =X(s) N(s)+Y(s) D(s)
\end{aligned}
$$

where $X(s):=R^{-1}(s) \hat{X}(s), Y(s):=R^{-1}(s) \hat{Y}(s)$.
Conversely, suppose that there exist two matrices $X(s)$ and $Y(s)$ satisfying eq. (2.8) and let $R(s)$ be a GCRD of $(N(s), D(s))$, i.e.

$$
N(s)=\hat{N}(s) R(s), \quad D(s)=\hat{D}(s) R(s)
$$

It is then possible to write $I=[X(s) \hat{N}(s)+Y(s) \hat{D}(s)] R(s)$ yielding

$$
\begin{equation*}
R^{-1}(s)=X(s) \hat{N}(s)+Y(s) \hat{D}(s) \tag{2.9}
\end{equation*}
$$

The right side of equation (2.9) is a polynomial matrix. This entails that $R(s)$ is a unimodular matrix so that $N(s)$ and $D(s)$ are right coprime.

Example 2.4 Consider the two polynomial matrices

$$
N(s)=\left[\begin{array}{cc}
2 s^{2}+1 & 2 s \\
s & 1
\end{array}\right], \quad D(s)=\left[\begin{array}{cc}
2 s^{3}+s & 2 s^{2}
\end{array}\right]
$$

They are right coprime. As a matter of fact

$$
\begin{aligned}
& X(s):=\left[\begin{array}{cc}
-2 s^{3}-2 s^{2}+1 & -2 s \\
2 s^{4}+2 s^{3}+s^{2} & 2 s^{2}+1
\end{array}\right] \\
& Y(s):=\left[\begin{array}{c}
2 s^{2}+2 s \\
-2 s^{3}-2 s^{2}-s-1
\end{array}\right]
\end{aligned}
$$

are such that $X(s) N(s)+Y(s) D(s)=I$. Moreover, taken

$$
T(s):=\left[\begin{array}{ccc}
-s^{3}-s^{2}+1 & -s & s^{2}+s \\
s^{2}+s & 1 & -s-1 \\
-s & 0 & 1
\end{array}\right]
$$

one can easily realize that it is unimodular and that

$$
T(s)\left[\begin{array}{l}
N(s) \\
D(s)
\end{array}\right]=\left[\begin{array}{c}
R(s) \\
0
\end{array}\right]
$$

with

$$
R(s)=\left[\begin{array}{cc}
s^{2}+1 & s \\
s & 1
\end{array}\right]
$$

which is unimodular as well.

Of course, a version of the result provided in Lemma 2.2 can be stated for left coprimeness as well. An important and significant canonical form can be associated with a polynomial matrix, namely the so called Smith form. This canonical form is formally defined in the following theorem whose proof also provides a systematic procedure for its computation.

Theorem 2.1 (Smith form) Let $N(s)$ be a $n \times m$ polynomial matrix and consider that $\operatorname{rank}[N(s)]:=r \leq \min [n, m]$. Then two polynomial and unimodular matrices $L(s)$ and $R(s)$ exist such that $N(s)=L(s) S(s) R(s)$ with

$$
\left.S(s)=\left[\begin{array}{ccccc}
\alpha_{1}(s) & 0 & \cdots & 0 & 0 \\
0 & \alpha_{2}(s) & \cdots & 0 & 0 \\
\vdots & \vdots & \ddots & \vdots & \vdots \\
0 & 0 & \cdots & \alpha_{r}(s) & 0 \\
0 & 0 & \cdots & 0 & 0
\end{array}\right]\right\} n-r \text { rows }
$$

where each polynomial $\alpha_{i}(s)$ is monic and divides the next one, i.e. $\alpha_{i}(s) \mid \alpha_{i+1}(s), i=$ $1,2, \cdots, r-1$. Matrix $S(s)$ is said to be the Smith form of $N(s)$.

Proof The proof of the theorem is constructive since the procedure to be described leads to the determination of the matrices $S(s), L(s)$ and $R(s)$. In the various steps which characterize such a procedure, the matrix $N(s)$ is subject to a number of manipulations resulting from suitable elementary operations on its rows and columns, i.e. pre-multiplications or post-multiplications by unimodular matrices. These operations determine the matrices $L(s)$ and $R(s)$. For simplicity, let $n_{i j}(s)$ be the ( $i, j$ ) element of the matrix which is presently considered.

1) Through two elementary operations on the rows and the columns of $N(s)$, bring a nonzero and minimum degree polynomial of $N(s)$ in position $(1,1)$.
2) Write the element $(2,1)$ of $N(s)$ as $n_{21}(s)=n_{11}(s) \gamma(s)+\beta(s)$, with $\beta(s)$ such that $\operatorname{deg}[\beta(s)]<\operatorname{deg}\left[n_{11}(s)\right]$. Now multiply the first row by $\gamma(s)$ and subtract the result from the second row. In this way the $(2,1)$ element becomes $\beta(s)$. Now, if $\beta(s)=0$ go to step 3 ), otherwise interchange the first row with the second one and repeat again this step. This causes a continuous reduction of the degree of the element $(2,1)$ so that, in a finite number of iterations, it results $n_{21}(s)=0$.
3) As exactly done in step 2), bring all the elements of the first column but element $(1,1)$ to zero.
4) Through elementary operations on the columns bring all the elements of the first row but $n_{11}(s)$ to zero.
5) If step 4) brought the elements of the first column under $n_{11}(s)$ to be nonzero, then go back to step 2). Notice that a finite number of operations through steps 2)-4) leads to a situation in which $n_{11}(s) \neq 0, n_{i 1}(s)=0, i=2,3, \cdots, n, n_{1 j}(s)=$ $0, j=2,3, \cdots, m$. If, in one of the columns aside the first one an element is not divisible by $n_{11}(s)$, add this column to the first one and go back to step 2). At each iteration of the cycle 2)-5) the degree of $n_{11}(s)$ decreases. Hence, in a finite number of cycles one arrives to the situation reported above where $n_{11}(s)$ divides each element of the submatrix $N_{1}(s)$ constituted by the last $m-1$ columns and $n-1$ rows. Assume that $n_{11}(s)$ is monic (otherwise perform an obvious elementary
operation) and let $\alpha_{1}(s)=n_{11}(s)$. Now apply to the submatrix $N_{1}(s)$ (obviously assumed to be nonzero) the entire procedure performed for matrix $N(s)$. The ( 1,1 ) element will be now $\alpha_{2}(s)$ and, in view of the adopted procedure $\alpha_{1}(s)$ will be a divisor of $\alpha_{2}(s)$. Finally $\alpha_{i}(s) \neq 0, i=1, \cdots, r$ since an elementary operation does not affect the matrix rank.

Example 2.5 Consider the polynomial matrix

$$
N(s)=\left[\begin{array}{cc}
s^{2}+1 & s^{3}+s \\
s+1 & s
\end{array}\right]
$$

and the two polynomial and unimodular matrices

$$
L(s)=\left[\begin{array}{cc}
\left(s^{2}+1\right) / 2 & s-1 \\
(s+1) / 2 & 1
\end{array}\right], \quad R(s)=\left[\begin{array}{cc}
2 & s^{3}-s^{2}+2 s \\
0 & -1 / 2
\end{array}\right]
$$

It follows that $N(s)=L(s) S(s) R(s)$ where

$$
S(s):=\left[\begin{array}{cc}
1 & 0 \\
0 & s^{4}+s^{2}
\end{array}\right]
$$

### 2.4 Proper rational matrices

This section deals with matrices $\mathrm{F}(\mathrm{s})$ whose elements are ratios of polynomials in the same unknown. Therefore, the generic element $f_{i j}(s)$ of $F(s)$ has the form

$$
\begin{gathered}
f_{i j}(s)=\frac{a(s)}{b(s)}:=\frac{\alpha_{\nu} s^{\nu}+\alpha_{\nu-1} s^{\nu-1}+\cdots+\alpha_{1} s+\alpha_{0}}{\beta_{\mu} s^{\mu}+\beta_{\mu-1} s^{\mu-1}+\cdots+\beta_{1} s+\beta_{0}} \\
\alpha_{i} \in R, i=0,1, \cdots, \nu, \beta_{i} \in R, i=0,1, \cdots \mu
\end{gathered}
$$

The relative degree reldeg $\left[f_{i j}(s)\right]$ of $f_{i j}(s)$ is defined as the difference between the degree of the two polynomials which constitute the denominator and numerator of $f_{i j}(s)$, respectively. Specifically, with reference to the above function, and assuming $\alpha_{\nu} \neq 0$ and $\beta_{\mu} \neq 0$, it is

$$
\operatorname{reldeg}\left[f_{i j}(s)\right]:=\operatorname{deg}[b(s)]-\operatorname{deg}[a(s)]=\mu-\nu
$$

A rational matrix $F(s)$ is said to be proper (resp. strictly proper) if reldeg $\left[f_{i j}(s)\right] \geq 0$ (resp. reldeg $\left[f_{i j}(s)\right]>0$ ) for all $i, j$. Throughout the section it is implicitly assumed that the rational matrices considered herein are always either proper or strictly proper.

The rank of a rational matrix $F(s)$ is, in analogy with the definition given for a polynomial matrix, the dimension of the largest square submatrix in $F(s)$ with determinant not identically equal to zero.

A rational square matrix is said to be unimodular if it has maximum rank and its determinant is a rational function with zero relative degree. Hence, a unimodular rational matrix admits a unimodular rational inverse and vice-versa.

Example 2.6 The matrix

$$
F(s)=\left[\begin{array}{cc}
\left(s^{2}+2 s+3\right) /\left(s^{2}-1\right) & 2 \\
s /(s+1) & 1
\end{array}\right]
$$

is unimodular. Actually, $\operatorname{det}[F(s)]=\left(-s^{2}+4 s+3\right) /\left(s^{2}-1\right)$. Moreover,

$$
F^{-1}(s)=\left[\begin{array}{cc}
\left(s^{2}-1\right) /\left(-s^{2}+4 s+3\right) & \left(-2 s^{2}+2\right) /\left(-s^{2}+4 s+3\right) \\
\left(-s^{2}+s\right) /\left(-s^{2}+4 s+3\right) & \left(s^{2}+2 s+3\right) /\left(-s^{2}+4 s+3\right)
\end{array}\right]
$$

The concepts of divisor and greatest common divisor, already given for polynomial matrices, are now extended to rational matrices in the definitions below.

Definition 2.4 (Right divisor) Let $F(s)$ be a rational matrix. A rational square matrix $R(s)$ is said to be a right divisor of $F(s)$ if

$$
F(s)=\bar{F}(s) R(s)
$$

with $\bar{F}(s)$ rational.
A similar definition could be given for a left divisor as well.
Definition 2.5 (Greatest common right divisor) Consider two rational matrices $F(s)$ and $G(s)$ with the same number of columns. A Greatest Common Right Divisor ( $G C R D$ ) of $(F(s), G(s))$ is a square rational matrix $R(s)$ such that
i) $R(s)$ is a right divisor of $(F(s), G(s))$, i.e.

$$
\begin{aligned}
& F(s)=\bar{F}(s) R(s) \\
& G(s)=\bar{G}(s) R(s)
\end{aligned}
$$

with $\bar{F}(s)$ and $\bar{G}(s)$ rational.
ii) If $\hat{R}(s)$ is any other right divisor of $(F(s), G(s))$, then $R(s)=W(s) \hat{R}(s)$ with $W(s)$ rational.

A similar definition holds for a Greatest Common Left Divisor (GCLD). By exploiting the properties of the rational unimodular matrices, it is easy to see that, given two rational matrices $F(s)$ and $G(s)$, there exist more than one GCRD (and GCLD). A way to compute a GCRD (resp. GCLD) of an assigned pair of rational matrices $F(s)$ and $G(s)$ calls for their manipulation via a rational unimodular matrix resulting from a sequence of elementary operations on their rows (resp. columns). The elementary operations on the rows (resp. columns) of a rational matrix $F(s)$ are:

1) Interchange of the i -th row (resp. i-th column) with the j -th row (resp. j -th column) and vice-versa
2) Multiplication of the i-th row (resp. i-th column) by a non zero rational function with zero relative degree
3) Addition of the i -th row (resp. i-th column) to the j-th row (resp. j-th column) multiplied by a rational function with zero relative degree

Obviously, each of these elementary operations reduces to premultiplying (resp. postmultiplying) matrix $F(s)$ by a suitable rational unimodular matrix $T(s)$. Moreover, matrix $T(s) F(s)(\operatorname{resp} . F(s) T(s))$ has the same rank as $F(s)$.

Remark 2.3 Given two scalar rational functions $f(s)$ and $g(s)$ with relative degree such that $\operatorname{reldeg}[f(s)] \leq \operatorname{reldeg}[g(s)]$, then reldeg $[g(s) / f(s)] \geq 0$. Hence, considering the rational unimodular matrix

$$
T(s):=\left[\begin{array}{cc}
1 & 0 \\
-\frac{g(s)}{f(s)} & 1
\end{array}\right]
$$

it follows that

$$
T(s)\left[\begin{array}{l}
f(s) \\
g(s)
\end{array}\right]=\left[\begin{array}{c}
f(s) \\
0
\end{array}\right]
$$

By recursively exploiting this fact, it is easy to convince oneself that, if a rational matrix $F(s)$ does not have more columns than rows, it is always possible to build up a rational unimodular matrix $T(s)$ such that

$$
T(s) F(s)=\left[\begin{array}{c}
R(s) \\
0
\end{array}\right]
$$

with $R(s)$ square and rational. Moreover, the null matrix vanishes when $F(s)$ is square.
A GCRD of two rational matrices can be computed in the way described in the following algorithm, which relies on the same arguments as in Algorithm 2.1.

Algorithm 2.2 Let $F(s)$ and $G(s)$ be two rational matrices with the same number, say $m$, of columns and possibly different numbers of rows, say $n_{f}$ and $n_{g}$, respectively.

1) Assume first that $m \leq n_{f}+n_{g}$, otherwise go to point 2). Let $H(s):=$ $\left[F^{\prime}(s) G^{\prime}(s)\right]^{\prime}$ and determine a rational unimodular matrix $T(s)$ such that

$$
T(s) H(s)=\left[\begin{array}{c}
R(s) \\
0
\end{array}\right]
$$

Then $R(s)$ is a GCRD
2) If $m>n_{f}+n_{g}$, then

$$
R(s)=\left[\begin{array}{c}
F(s) \\
G(s) \\
0
\end{array}\right]
$$

is a GCRD
Example 2.7 Consider the two rational matrices

$$
\begin{aligned}
& F(s)=\left[\begin{array}{cc}
(s+1) / s & 1 /(s+2)
\end{array}\right] \\
& G(s)=\left[\begin{array}{cc}
(s+1) /(s-1) & s /\left(s^{2}+1\right) \\
1 & s /(s+1)
\end{array}\right]
\end{aligned}
$$

Take now

$$
\begin{aligned}
& T_{1}(s)=\left[\begin{array}{ccc}
1 & 0 & -(s+1) / s \\
0 & 1 & -(s+1) /(s-1) \\
0 & 0 & 1
\end{array}\right], T_{2}(s)=\left[\begin{array}{lll}
1 & 0 & 1 \\
1 & 0 & 0 \\
0 & 1 & 0
\end{array}\right] \\
& T_{3}(s)=\left[\begin{array}{ccc}
1 & 0 & 0 \\
0 & 1 & 0 \\
0 & \left(s^{4}+s^{3}+4 s\right) /\left(1-s^{4}\right) & 1
\end{array}\right]
\end{aligned}
$$

It follows that

$$
T_{3}(s) T_{2}(s) T_{1}(s)\left[\begin{array}{c}
F(s) \\
G(s)
\end{array}\right]=\left[\begin{array}{cc}
1 & s /(s+1) \\
0 & -(s+1) /(s+2) \\
0 & 0
\end{array}\right]
$$

so that

$$
\begin{aligned}
& R(s)=\left[\begin{array}{cc}
1 & s /(s+1) \\
0 & -(s+1) /(s+2)
\end{array}\right], \quad \bar{F}(s)=\left[\begin{array}{ll}
(s+1) / s & 1
\end{array}\right] \\
& \bar{G}(s)=\left[\begin{array}{cc}
(s+1) /(s-1) & \left(s^{4}+s^{3}+4 s\right) /\left(s^{4}-1\right) \\
1 & 0
\end{array}\right]
\end{aligned}
$$

and $F(s)=\bar{F}(s) R(s), G(s)=\bar{G}(s) R(s)$.
Also for rational matrices it is possible to consistently introduce the concept of coprimeness.

Definition 2.6 (Right coprimeness) Two rational matrices $F(s)$ and $G(s)$ with the same number of columns are said to be right coprime if the relations

$$
\begin{aligned}
& F(s)=\bar{F}(s) T(s) \\
& G(s)=\bar{G}(s) T(s)
\end{aligned}
$$

with $\bar{F}(s)$ and $\bar{G}(s)$ rational matrices, are verified only if $T(s)$ is a rational unimodular matrix.

Example 2.8 The two matrices

$$
\begin{aligned}
& F(s)=\left[\begin{array}{cc}
(s-1) /(s+1)^{2} & \left(s^{2}+1\right) /(s+1)\left(s^{2}+3 s\right) \\
1 /(s+1) & s /\left(s^{2}-1\right)
\end{array}\right] \\
& G(s)=\left[\begin{array}{ll}
1 /(s+1) & (2 s+1) /(s+1)(s+3)
\end{array}\right]
\end{aligned}
$$

are not right coprime since it results

$$
\begin{aligned}
& F(s)=\left[\begin{array}{cc}
(s-1) /(s+1) & \left(s^{2}+1\right) /\left(s^{2}+3 s\right) \\
1 & s /(s-1)
\end{array}\right] R(s) \\
& G(s)=\left[\begin{array}{ll}
1 & (2 s+1) /(s+3)
\end{array}\right] R(s)
\end{aligned}
$$

with $R(s)=1 /(s+1)$ which is not unimodular.
An analogous definition holds for left coprimeness. From Definition 2.6 it follows that two rational matrices are right (resp. left) coprime if all their common right (resp. left) divisors are unimodular. In particular, each GCRD (resp. GCLD) of two rational right (resp. left) coprime matrices must be unimodular. Therefore, a necessary condition for matrices $F(s)$ and $G(s)$ to be right (resp. left) coprime is that the number of their columns be not greater than the sum of the number of their rows (resp. the number of their rows be not greater than the sum of the number of their columns), since, from Algorithm 2.2 point 2) in the opposite case one of their GCRD would not be unimodular. Moreover, a way to verify whether or not two rational matrices are right (resp. left) coprime consists in the computation through Algorithm 2.2 of a greatest common divisor and evaluation of its determinant. As a matter of fact, as stated in the next lemma, if a greatest common divisor is unimodular then all the greatest common divisors are unimodular as well.

Lemma 2.3 Let $F(s)$ and $G(s)$ be two rational matrices with the same number of columns. Let $R(s)$ be a $G C R D$ of $(F(s), G(s))$. Then,
i) $T(s) R(s)$ is a $G C R D$ for any rational unimodular $T(s)$
ii) If $\bar{R}(s)$ is a $G C R D$ of $(F(s), G(s))$, then there exists a rational unimodular matrix $T(s)$ such that $\bar{R}(s)=T(s) R(s)$

Proof The proof follows from that of Lemma 2.1 by substituting there the term "rational" in place of the term "polynomial" and symbols $F(s)$ and $G(s)$ in place of $N(s)$ and $D(s)$, respectively.

A further significant property of a GCRD of a pair of rational matrices $F(s)$ and $G(s)$ is stated in the following lemma, whose proof hinges on Algorithm 2.2.

Lemma 2.4 Consider two rational matrices $F(s)$ and $G(s)$ with the same number of columns and let $R(s)$ be a $G C R D$ of $(F(s), G(s))$. Then, there exist two rational matrices $X(s)$ and $Y(s)$ such that

$$
X(s) F(s)+Y(s) G(s)=R(s)
$$

Proof Let $n_{f}$ and $n_{g}$ be the number of rows of $F(s)$ and $G(s)$, respectively, and $m$ the number of their columns. Preliminarily, assume that $n_{f}+n_{g} \geq m$ and let $T(s)$ be a unimodular matrix such that

$$
T(s)\left[\begin{array}{c}
F(s)  \tag{2.10}\\
G(s)
\end{array}\right]=\left[\begin{array}{ll}
T_{11}(s) & T_{12}(s) \\
T_{21}(s) & T_{22}(s)
\end{array}\right]\left[\begin{array}{c}
F(s) \\
G(s)
\end{array}\right]=\left[\begin{array}{c}
\bar{R}(s) \\
0
\end{array}\right]
$$

Based on Algorithm 2.2, matrix $\bar{R}(s)$ turns out to be a GCRD of $(F(s), G(s))$. Hence, thanks to Lemma 2.3, there exist a unimodular matrix $U(s)$ such that $R(s)=$ $U(s) \bar{R}(s)$, that is, in view of eq. (2.10),

$$
R(s)=U(s) T_{11}(s) F(s)+U(s) T_{12}(s) G(s)=X(s) F(s)+Y(s) G(s)
$$

On the contrary, if $m>n_{f}+n_{g}$, Algorithm 2.2 entails that

$$
\bar{R}(s):=\left[\begin{array}{c}
F(s) \\
G(s) \\
0
\end{array}\right]
$$

is a GCRD of $(F(s), G(s))$. In view of Lemma 2.3, it is possible to write

$$
\begin{aligned}
R(s) & =U(s) \bar{R}(s) \\
& =U(s)\left[\begin{array}{l}
I \\
0 \\
0
\end{array}\right] F(s)+U(s)\left[\begin{array}{l}
0 \\
I \\
0
\end{array}\right] G(s) \\
& :=X(s) F(s)+Y(s) G(s)
\end{aligned}
$$

where $U(s)$ is a suitable rational and unimodular matrix.
The following result, that parallels the analogous one presented in Lemma 2.2, can now be stated.

Lemma 2.5 Let $F(s)$ and $G(s)$ be two rational matrices with the same number of columns. Then, $F(s)$ and $G(s)$ are right coprime if and only if there exist two rational matrices $X(s)$ and $Y(s)$ such that

$$
X(s) F(s)+Y(s) G(s)=I
$$

Proof Recall that two matrices are right coprime if each one of their GCRD's is unimodular. Hence, if $R(s)$ is a GCRD of $(F(s), G(s))$, thanks to Lemma 2.4, it results $R(s)=\bar{X}(s) F(s)+\bar{Y}(s) G(s)$, with $\bar{X}(s)$ and $\bar{Y}(s)$ suitable rational matrices. From this last equation it follows

$$
I=R^{-1}(s) \bar{X}(s) F(s)+R^{-1}(s) \bar{Y}(s) G(s):=X(s) F(s)+Y(s) G(s)
$$

Conversely, let $R(s)$ be a GCRD of $(F(s), G(s))$ derived according to Algorithm 2.2 point 1), as the number of their columns must be not greater than the sum of the numbers of their rows, so that

$$
T(s)\left[\begin{array}{l}
F(s) \\
G(s)
\end{array}\right]=\left[\begin{array}{c}
R(s) \\
0
\end{array}\right]
$$

where $T(s)$ is a suitable rational and unimodular matrix. Hence,

$$
\left[\begin{array}{c}
F(s) \\
G(s)
\end{array}\right]=T^{-1}(s)\left[\begin{array}{c}
R(s) \\
0
\end{array}\right]:=\left[\begin{array}{cc}
S_{11}(s) & S_{12}(s) \\
S_{21}(s) & S_{22}(s)
\end{array}\right]\left[\begin{array}{c}
R(s) \\
0
\end{array}\right]
$$

so that

$$
I=X(s) F(s)+Y(s) G(s)=\left[X(s) S_{11}(s)+Y(s) S_{21}(s)\right] R(s)
$$

shows that $R(s)$ is unimodular (its inverse is rational). Therefore, $(F(s), G(s))$ are right coprime.

Example 2.9 Consider two rational matrices

$$
\begin{aligned}
& F(s)=\left[\begin{array}{cc}
1 /\left(s^{2}+s\right) & \left(2 s^{2}+2 s-2\right) /\left(s^{2}+2 s\right)
\end{array}\right] \\
& G(s)=\left[\begin{array}{cc}
-s /(s+1) & s /(s+1) \\
\left(s^{2}-s-1\right) /\left(s^{2}+s\right) & \left(s^{2}+2 s+2\right) /\left(s^{2}+2 s\right)
\end{array}\right]
\end{aligned}
$$

These matrices are right coprime. Actually,

$$
T(s)\left[\begin{array}{c}
F(s) \\
G(s)
\end{array}\right]=\left[\begin{array}{c}
R(s) \\
0
\end{array}\right]
$$

where

$$
T(s)=\frac{1}{p(s)}\left[\begin{array}{ccc}
s^{3}-s^{2}-2 s-1 & -s^{3}+s & s^{3}+s^{2}-s-1 \\
s^{3}+2 s^{2}+s & s^{3}+s^{2} & -s^{3}-2 s^{2}-s \\
-s^{3} & s^{3}-s^{2}-2 s & s^{3}
\end{array}\right]
$$

with $p(s):=2 s^{3}-3 s-2$ and

$$
R(s)=\left[\begin{array}{cc}
s /(s+1) & (s+1) /(s+2) \\
-1 & 1
\end{array}\right]
$$

that is a rational and unimodular matrix. Moreover, taking

$$
X(s)=\frac{1}{p(s) q(s)}\left[\begin{array}{c}
-2 s^{4}-9 s^{3}-13 s^{2}-8 s-2 \\
2 s^{5}+6 s^{4}+2 s^{3}-7 s^{2}-7 s-2
\end{array}\right]
$$

and

$$
Y(s)=\frac{1}{p(s) q(s)}\left[\begin{array}{cc}
-2 s^{5}-6 s^{4}-4 s^{3}+2 s^{2}+2 s & 2 s^{5}+8 s^{4}+10 s^{3}+2 s^{2}-4 s-2 \\
s^{3}+3 s^{2}+2 s & -s^{3}-4 s^{2}-5 s-2
\end{array}\right]
$$

with $q(s):=2 s^{2}+4 s+\mathbf{1}$, it follows $X(s) F(s)+Y(s) G(s)=I$.
Also for rational matrices there exist a particularly useful canonical form, which is called Smith-McMillan form. This form is precisely defined in the following theorem, whose proof also provides a procedure for its computation.

Theorem 2.2 (Smith-McMillan form) Let $G(s)$ be a proper rational matrix with $n$ rows, $m$ columns and $\operatorname{rank}[G(s)]=r \leq \min [n, m]$. Then there exist two polynomial and unimodular matrices $L(s)$ and $R(s)$ such that $G(s)=L(s) M(s) R(s)$, where

$$
\left.M(s)=\left[\begin{array}{ccccc}
f_{1}(s) & 0 & \cdots & 0 & 0 \\
0 & f_{2}(s) & \cdots & 0 & 0 \\
\vdots & \vdots & \ddots & \vdots & \vdots \\
0 & 0 & \cdots & f_{r}(s) & 0 \\
0 & 0 & \cdots & 0 & 0
\end{array}\right]\right\} n-r \text { rows }
$$

with

$$
f_{i}(s)=\frac{\varepsilon_{i}(s)}{\psi_{i}(s)}, i=1,2, \cdots, r
$$

and

- $\varepsilon_{i}(s)$ and $\psi_{i}(s)$ are monic, $i=1,2, \cdots r$
- $\varepsilon_{i}(s)$ and $\psi_{i}(s)$ are coprime $i=1,2, \cdots r$
- $\varepsilon_{i}(s)$ divides $\varepsilon_{i+1}(s), i=1,2, \cdots r$
- $\psi_{i+1}$ divides $\psi_{i}, i=1,2, \cdots r$

Matrix $M(s)$ is the Smith-McMillan form of $F(s)$.
Proof Let $\psi(s)$ be the least common multiple of all polynomials at the denominators of the elements of $F(s)$. Therefore, matrix $N(s):=\psi(s) F(s)$ is polynomial. If $S(s)$ is the Smith form of $N(s)$ (recall Theorem 2.1) it follows that

$$
F(s)=\frac{1}{\psi(s)} N(s)=\frac{1}{\psi(s)} L(s) S(s) R(s)
$$

Hence,

$$
M(s)=\frac{S(s)}{\psi(s)}
$$

once all the possible simplifications between the elements of $S(s)$ and the polynomial $\psi(s)$ have been performed. This matrix obviously has the properties claimed in the statement.

Remark 2.4 The result stated in Theorem 2.2 allows one to represent a generic rational $p \times m$ matrix $G(s)$ with $\operatorname{rank}[G(s)]:=r \leq \min [m, p]$ in the two forms

$$
G(s)=N(s) D^{-1}(s)=\hat{D}^{-1}(s) \hat{N}(s)
$$

where the polynomial matrices $N(s)$ and $D(s)$ are right coprime, while the polynomial matrices $\hat{D}(s)$ and $\hat{N}(s)$ are left coprime. Actually, observe that letting

$$
\begin{aligned}
& \Psi(s):=\left[\begin{array}{cccc}
\psi_{1}(s) & 0 & \cdots & 0 \\
0 & \psi_{2}(s) & \cdots & 0 \\
\vdots & \vdots & \ddots & \vdots \\
0 & 0 & \cdots & \psi_{r}(s)
\end{array}\right] \\
& E(s):=\left[\begin{array}{cccc}
\varepsilon_{1}(s) & 0 & \cdots & 0 \\
0 & \varepsilon_{2}(s) & \cdots & 0 \\
\vdots & \vdots & \ddots & \vdots \\
0 & 0 & \cdots & \varepsilon_{r}(s)
\end{array}\right]
\end{aligned}
$$

it follows

$$
\begin{aligned}
M(s) & =\left[\begin{array}{cc}
\Psi^{-1}(s) & 0 \\
0 & I
\end{array}\right]\left[\begin{array}{cc}
E(s) & 0 \\
0 & 0
\end{array}\right] \\
& =\left[\begin{array}{cc}
E(s) & 0 \\
0 & 0
\end{array}\right]\left[\begin{array}{cc}
\Psi^{-1}(s) & 0 \\
0 & I
\end{array}\right]
\end{aligned}
$$

Now, defining

$$
\begin{aligned}
& N(s):=L(s)\left[\begin{array}{cc}
E(s) & 0 \\
0 & 0
\end{array}\right], \quad D(s):=R^{-1}(s)\left[\begin{array}{cc}
\Psi(s) & 0 \\
0 & I
\end{array}\right] \\
& \hat{N}(s):=\left[\begin{array}{cc}
E(s) & 0 \\
0 & 0
\end{array}\right] R(s), \quad \hat{D}(s):=\left[\begin{array}{cc}
\Psi(s) & 0 \\
0 & I
\end{array}\right] L^{-1}(s)
\end{aligned}
$$

one can easily check that $G(s)=N(s) D^{-1}(s)=\hat{D}^{-1}(s) \hat{N}(s)$. In order to verify that $N(s)$ and $D(s)$ are right coprime, one can resort to Lemma 2.2. Actually, considering the two matrices $X(s)$ and $Y(s)$ defined by

$$
\begin{aligned}
& X(s):=\left[\begin{array}{ccccc}
x_{1}(s) & 0 & \cdots & 0 & 0 \\
0 & x_{2}(s) & \cdots & 0 & 0 \\
\vdots & \vdots & \ddots & x_{r}(s) & 0 \\
0 & 0 & \cdots & 0 & 0
\end{array}\right] L^{-1}(s) \\
& Y(s):=\left[\begin{array}{ccccc}
y_{1}(s) & 0 & \cdots & 0 & 0 \\
0 & y_{2}(s) & \cdots & 0 & 0 \\
\vdots & \vdots & \ddots & y_{r}(s) & 0 \\
0 & 0 & \cdots & 0 & I
\end{array}\right] R(s)
\end{aligned}
$$

it turns out that $X(s) N(s)+Y(s) D(s)=I$ for a suitable choice of the polynomials $x_{i}(s)$ and $y_{i}(s)$ (recall that the polynomials $\psi_{i}(s)$ and $\varepsilon_{i}(s)$ are coprime and take in mind Theorem A.1). From Lemma 2.2 it follows that the two matrices $N(s)$ and $D(s)$ are right coprime. Analogously, one can verify that $\hat{N}(s)$ and $\hat{D}(s)$ are left coprime.

Example 2.10 Consider the rational matrix

$$
F(s)=\frac{1}{s^{3}-s^{2}}\left[\begin{array}{cc}
3 s^{3}+7 s^{2}+3 s-1 & 4 s^{3}+8 s^{2}+2 s-2 \\
2 s^{3}+4 s^{2}+s-1 & 3 s^{3}+5 s^{2}-2
\end{array}\right]
$$

Now applying what has been indicated in the proof of Theorems 2.1, 2.2 it follows that $F(s)=L(s) M(s) R(s)$, with

$$
\begin{aligned}
L(s) & =\left[\begin{array}{cc}
\left(3 s^{2}+4 s-1\right) / 3 & (6 s+5) / 12 \\
\left(2 s^{2}+2 s-1\right) / 3 & (2 s+1) / 6
\end{array}\right] \\
R(s) & =\left[\begin{array}{cc}
3 & -2 s^{3}-3 s^{2}+2 s+6 \\
0 & 4
\end{array}\right] \\
M(s) & =\left[\begin{array}{cc}
(s+1) /\left(s^{3}-s^{2}\right) & 0 \\
0 & (s+1)^{2}(s+2) / s
\end{array}\right]
\end{aligned}
$$

Moreover, by following the arguments in Remark 2.4,

$$
\begin{aligned}
& E(s)=\left[\begin{array}{cc}
s+1 & 0 \\
0 & (s+1)^{2}(s+2)
\end{array}\right] \\
& \Psi(s)=\left[\begin{array}{cc}
s^{2}(s-1) & 0 \\
0 & s
\end{array}\right]
\end{aligned}
$$

it turns out that

$$
\begin{aligned}
& N(s)=L(s) E(s)=\frac{1}{12}\left[\begin{array}{cc}
4\left(3 s^{3}+7 s^{2}+3 s-1\right) & 6 s^{4}+29 s^{3}+50 s^{2}+37 s+10 \\
4\left(2 s^{3}+4 s^{2}+s-1\right) & 2\left(2 s^{4}+9 s^{3}+14 s^{2}+9 s+2\right)
\end{array}\right] \\
& D(s)=R^{-1}(s) \Psi(s)=\frac{1}{12}\left[\begin{array}{cc}
4\left(s^{3}-s\right) & s\left(2 s^{3}+3 s^{2}-2 s-6\right) \\
0 & 3 s
\end{array}\right] \\
& \hat{N}(s)=E(s) R(s)=\left[\begin{array}{cc}
3(s+1) & -2 s^{4}-5 s^{3}-s^{2}+8 s+6 \\
0 & 4\left(s^{3}+4 s^{2}+5 s+2\right)
\end{array}\right] \\
& \hat{D}(s)=\Psi(s) L^{-1}(s)=\left[\begin{array}{cc}
2 s^{2}\left(2 s^{2}-s-1\right) & -s^{2}\left(6 s^{2}-s-5\right) \\
-4 s\left(2 s^{2}+2 s-1\right) & 4 s\left(3 s^{2}+4 s-1\right)
\end{array}\right]
\end{aligned}
$$

Many of the results provided till now can be straightforwardly extended to the subset of proper, rational and stable functions, namely the subset constituted by the matrices whose generic element $f_{i j}(s)$ is a proper, rational function with poles in the open left half plane only. This extension calls for the introduction of a suitable scalar associated with a generic proper rational scalar function $f(s)$. Precisely, rhpdeg $[f(s)]$ will indicate the number of finite nonnegative real part zeros of $f(s)$ plus reldeg $[f(s)]$. For example

$$
f(s)=\frac{(s-1)(s+2)}{(s+4)^{2}}
$$

is such that rhpdeg $[f(s)]=3$ since reldeg $[f(s)]=2$ and $f(s)$ has a zero in $s=1$. Moreover, the function

$$
f(s)=\frac{s}{s+1}
$$

is such that $\operatorname{rhpdeg}[f(s)]=1$, since $f(s)$ has a zero in $s=0$ and $\operatorname{reldeg}[f(s)]=0$. It will conventionally be set rhpdeg $[0]=-\infty$. Preliminarily, observe that the definitions
of divisor and greatest common divisors of rational matrices (Definitions 2.4 and 2.5) can be trivially generalized to the subset of stable matrices by actually requiring the stability property. As for the generalization of the concept of unimodular matrix $T(s)$, it suffices requiring, besides the stability of $T(s)$, also that $\operatorname{rhpdeg}[\operatorname{det}[T(s)]]=0$. In this way, the stable matrix $T(s)$ has a stable inverse as well.

The three elementary operations on the rows (columns) of a rational matrix are extended to stable rational matrices by simply requiring that in the second operation the multiplying function $f(s)$ be stable with $\operatorname{rhpdeg}[f(s)]=0$ and simply that in the third operation $f(s)$ be stable.

Lemma 2.6 Let $f(s)$ and $g(s)$ be two stable rational scalar functions with $g(s) \neq 0$ and $\operatorname{rhpdeg}[f(s)] \geq \operatorname{rhpdeg}[g(s)]$. Then, there exists a stable rational function $q(s)$ such that

$$
\begin{equation*}
\operatorname{rhpdeg}[f(s)-g(s) q(s)]<\operatorname{rhpdeg}[g(s)] \tag{2.11}
\end{equation*}
$$

Proof If rhpdeg $[g(s)]=0$, then $g^{-1}(s)$ is rational, proper and stable so that equation (2.11) is obviously satisfied with $q(s)=g^{-1}(s) f(s)$. Therefore, suppose that $\operatorname{rhpdeg}[g(s)]:=\nu \neq 0$ and write

$$
g(s):=\frac{n g(s)}{d g(s)}:=\frac{n g^{+}(s) n g^{-}(s)}{d g(s)}
$$

where the polynomials $n g(s)$ and $d g(s)$ are coprime whereas $n g^{+}(s)$ has roots in the closed right half plane only and $n g^{-}(s)$ in the open left half plane only. Moreover, let

$$
h(s):=\frac{(s+1)^{\nu} n g^{-}(s)}{d g(s)}
$$

so that both $h(s)$ and $h^{-1}(s)$ are proper stable rational functions. Of course, it results

$$
g(s)=\frac{h(s) n g^{+}(s)}{(s+1)^{\nu}}
$$

Also, write

$$
f(s):=\frac{n f(s)}{d f(s)}
$$

where the two polynomials $n f(s)$ and $d f(s)$ are coprime. Notice that, being $f(s)$ stable, the zeros of $d f(s)$ are in the open left half plane. This entails that $d f(s)$ and $n g^{+}(s)$ are coprime. By exploiting Lemma A.2, one can claim that there exist two polynomials $\varphi(s)$ and $\psi(s)$ with $\operatorname{deg}[\varphi(s)]<\operatorname{deg}[d f(s)]$ such that

$$
n g^{+}(s) \varphi(s)+d f(s) \psi(s)=(s+1)^{\nu-1} n f(s)
$$

From this relation it follows

$$
\frac{n g^{+}(s) \varphi(s)}{d f(s)(s+1)^{\nu-1}}+\frac{\psi(s)}{(s+1)^{\nu-1}}=\frac{n f(s)}{d f(s)}
$$

Let now

$$
q(s):=\frac{(s+1) \phi(s)}{d f(s) h(s)}
$$

and observe that such function is rational, proper and stable. This conclusion derives from properness and stability of $h^{-1}(s)$ and the fact that $\varphi(s) / d f(s)$ is strictly proper and stable. Moreover, let

$$
r(s):=\frac{\psi(s)}{(s+1)^{\nu-1}}
$$

then,

$$
\begin{aligned}
q(s) g(s)+r(s) & =\frac{(s+1) \varphi(s) d g(s) n g^{+}(s) n g^{-}(s)}{(s+1)^{\nu} d f(s) n g^{-}(s) d g(s)}+\frac{\psi(s)}{(s+1)^{\nu-1}} \\
& =\frac{(s+1) \varphi(s) n g^{+}(s)}{(s+1)^{\nu} d f(s)}+\frac{\psi(s)}{(s+1)^{\nu-1}} \\
& =\frac{n f(s)}{d f(s)}=f(s)
\end{aligned}
$$

Thus, being $f(s), g(s)$ and $q(s)$ proper, rational and stable, the function $r(s)=$ $f(s)-q(s) g(s)$ is rational, proper and stable as well. Finally, recalling the definition of $r(s)$, one can conclude that

$$
\operatorname{rhpdeg}[r(s)] \leq \nu-1<\nu=\operatorname{rhpdeg}[g(s)]
$$

Remark 2.5 Lemma 2.6 allows one to discuss further what has been shown in Remark 2.3 , in the context of stable matrices. As a matter of fact, let $f(s)$ and $g(s)$ be two rational stable scalar functions with

$$
\operatorname{rhpdeg}[f(s)] \geq \operatorname{rhpdeg}[g(s)]
$$

In view of Lemma 2.6 there exists a stable rational matrix $q(s)$ such that

$$
\operatorname{rhpdeg}[f(s)-q(s) g(s)]<\operatorname{rhpdeg}[g(s)]
$$

Then, the unimodular stable matrix

$$
T_{1}(s)=\left[\begin{array}{cc}
1 & -q(s) \\
0 & 1
\end{array}\right]
$$

is such that

$$
T_{1}(s)\left[\begin{array}{l}
f(s) \\
g(s)
\end{array}\right]=\left[\begin{array}{c}
f(s)-q(s) g(s) \\
g(s)
\end{array}\right]:=\left[\begin{array}{l}
f^{1}(s) \\
g^{1}(s)
\end{array}\right]
$$

By iterating this operation with stable unimodular matrices of the given form (or, alternatively, of the form corresponding to its transpose), one get

$$
\prod_{i=0}^{n-1} T_{n-i}(s)\left[\begin{array}{l}
f(s) \\
g(s)
\end{array}\right]=\left[\begin{array}{l}
f^{n}(s) \\
g^{n}(s)
\end{array}\right]
$$

with either $\operatorname{rhpdeg}\left[f^{n}(s)\right]=0$ or $\operatorname{rhpdeg}\left[g^{n}(s)\right]=0$. Assuming, for instance, that the first situation has occurred, it follows that

$$
\prod_{i=-1}^{n-1} T_{n-i}(s)\left[\begin{array}{l}
f(s) \\
g(s)
\end{array}\right]=\left[\begin{array}{c}
f^{n}(s) \\
0
\end{array}\right]
$$

where

$$
T_{n+\mathrm{I}}=\left[\begin{array}{cc}
1 & 0 \\
-\left(f^{n}(s)\right)^{-1} g^{n}(s) & 1
\end{array}\right]
$$

Then, one can conclude that, given two stable rational scalar functions $f(s)$ and $g(s)$, there exists a unimodular stable matrix $T(s)$ such that

$$
T(s)\left[\begin{array}{l}
f(s) \\
g(s)
\end{array}\right]=\left[\begin{array}{c}
r(s) \\
0
\end{array}\right]
$$

with $r(s)$ stable and $\operatorname{rhpdeg}[r(s)]=0$.
The arguments used in Remark 2.5 are well suited to be extended in the context of proper stable rational matrices. The same occurs for Lemmas 2.3, 2.4 provided that Definition 2.6 on coprimeness is adapted to this new setting. With these arguments in mind, it is possible to state the following lemma, which specializes Lemma 2.5 to the case of stable rational matrices.

Lemma 2.7 Let $F(s)$ and $G(s)$ be two stable rational matrices with the same number of columns. Then $F(s)$ and $G(s)$ are right coprime (in the setting of proper stable matrices) if and only if there exist two stable rational matrices $X(s)$ and $Y(s)$ such that

$$
X(s) F(s)+Y(s) G(s)=I
$$

Example 2.11 Consider the two stable rational functions

$$
f(s)=\frac{s+2}{s+1}, \quad g(s)=\frac{(s+2)^{2}}{(s+1)^{2}}
$$

They are (right) coprime. Actually, taking

$$
x(s):=\frac{s+1}{2(s+2)}, y(s):=\frac{(s+1)^{2}}{2(s+2)^{2}}
$$

it follows

$$
x(s) f(s)+y(s) g(s)=1
$$

On the other side, consider the stable unimodular matrix

$$
T(s):=\left[\begin{array}{cc}
1 & 0 \\
-(f(s))^{-1} g(s) & 1
\end{array}\right]
$$

Then

$$
T(s)\left[\begin{array}{c}
f(s) \\
g(s)
\end{array}\right]=\left[\begin{array}{c}
f(s) \\
0
\end{array}\right]
$$

with $f(s)$ stable and unimodular.

### 2.5 Poles and zeros

This section is devoted to a schematic presentation of the main properties of the poles and zeros of a linear and time-invariant dynamic system $\Sigma$,

$$
\begin{aligned}
\dot{x} & =A x+B u \\
y & =C x+D u
\end{aligned}
$$

with $n$ states, $m$ inputs and $p$ outputs. It will be of main importance in the sequel to distinguish between two cases. In the first one, reference is only made to an inputoutput description of the system, i.e to its transfer function

$$
G(s)=C(s I-A)^{-1} B+D
$$

whereas in the second case a state-space description of the system is considered. In order to rule out trivialities and make simpler the exposition, it is assumed, throughout all the section, that $G(s)$ has full rank, i.e. $\operatorname{rank}[G(s)]=\min [p, m]:=\mathrm{r}$.

Definition 2.7 (Zeros and poles of a rational matrix) Consider the rational matrix $G(s)$ and its associated Smith-McMillan form

$$
M(s)=\left[\begin{array}{ccccc}
f_{1}(s) & 0 & \cdots & 0 & 0 \\
0 & f_{2}(s) & \cdots & 0 & 0 \\
\vdots & \vdots & \ddots & \vdots & \vdots \\
0 & 0 & \cdots & f_{r}(s) & 0 \\
0 & 0 & \cdots & 0 & 0
\end{array}\right]
$$

with

$$
f_{i}(s):=\frac{\varepsilon_{i}(s)}{\psi_{i}(s)}, i=1,2, \cdots, r
$$

and define the polynomials $\pi_{p}(s)$ and $\pi_{z t}(s)$ as

$$
\begin{aligned}
\pi_{p}(s) & :=\psi_{1}(s) \psi_{2}(s) \cdots \psi_{r}(s) \\
\pi_{z t}(s) & :=\varepsilon_{1}(s) \varepsilon_{2}(s) \cdots \varepsilon_{r}(s)
\end{aligned}
$$

The poles of $G(s)$ are defined as the roots of $\pi_{p}(s)$ and the zeros of $G(s)$ as those of $\pi_{z t}(s)$.

The definition of zeros and poles of $G(s)$ coincides with that of transmission zeros and transmission poles of a system having $G(s)$ as transfer function. As customary, the transmission poles will be simply referred to as poles of the system.

Definition 2.8 Consider a linear time invariant system $\Sigma$ with transfer function $G(s)$. The transmission zeros (resp. poles) of $\Sigma$ are the zeros (resp. poles) of $G(s) . \square$

Example 2.12 Consider the linear system $\Sigma(A, B, C, D)$ defined by matrices

$$
\begin{gathered}
A=\left[\begin{array}{rrrr}
0 & 1 & 0 & 0 \\
-10 & 7 & 0 & 0 \\
0 & 0 & 5 & 0 \\
1 & -1 & 1 & 0
\end{array}\right], \quad B=\left[\begin{array}{l}
0 \\
1 \\
2 \\
0
\end{array}\right] \\
C=\left[\begin{array}{rrrr}
-13 & 5 & 0 & 0 \\
0 & 0 & 1 & 0
\end{array}\right], \quad D=\left[\begin{array}{l}
1 \\
1
\end{array}\right]
\end{gathered}
$$

The transfer function and its Smith-McMillan form are given by

$$
G(s)=\left[\begin{array}{c}
(s-3)(s+1) /(s-5)(s-2) \\
(s-3) /(s-5)
\end{array}\right]=\left[\begin{array}{ll}
s+1 & 1 \\
s-2 & 1
\end{array}\right]\left[\begin{array}{c}
(s-3) /(s-5)(s-2) \\
0
\end{array}\right]
$$

Therefore there is only one transmission zero in $s=3$.

Remark 2.6 In general, two polynomials $\varepsilon_{j}(s)$ and $\psi_{i}(s), i \neq j$, can have common roots. This entails that a multivariable system may admit coincident poles and zeros even in case of minimallity of its state-space description.

The Smith-McMillan form of $G(s)$ allows one for an alternative characterization of the transmission zeros in terms of vectors belonging to the kernel of $G(s)$ or $G^{\prime}(s)$, as proved in the following lemma.

Lemma 2.8 (Rank property of transmission zeros) Consider a transfer function $G(s)$ of rank $r:=\min [p, m]$. The complex number $\lambda$ is a transmission zero of $G(s)$ if and only if there exists a non zero vector $z$ such that

$$
\left\{\begin{array}{ccc}
\lim _{s \rightarrow \infty} & G(s) z=0 & \text { if } p \geq m \\
\lim _{s \rightarrow \infty} & G^{\prime}(s) z=0 & \text { if } p \leq m
\end{array}\right.
$$

Proof Consider first the case $p \geq m$ and let $M(s)$ be the Smith-McMillan form of $G(s)$, so that $G(s)=L(s) M(s) R(s)$, where $L(s)$ and $R(s)$ are suitable polynomial and unimodular matrices of dimensions $p$ and $m$, respectively. It is possible to write (recall Remark 2.4)

$$
M(s)=\left[\begin{array}{c}
E(s) \Psi^{-1}(s) \\
0
\end{array}\right]
$$

where $E(s):=\operatorname{diag}\left\{\varepsilon_{1}(s), \cdots, \varepsilon_{m}(s)\right\}, \Psi(s):=\operatorname{diag}\left\{\psi_{1}(s), \cdots, \psi_{m}(s)\right\}$. Further, denote by $e_{k}(h)$ the $k$-th column of the $h$-dimensional identity matrix and let $\lambda$ be a zero of $G(s)$, root of the polynomial $\varepsilon_{k}(s)$ of $E(s)$. Since $\psi_{k}(\lambda) \neq 0$ and $R(\lambda)$ is nonsingular, it then follows that $z=R^{-1}(\lambda) e_{k}(m)$ satisfies the condition of the theorem.

Conversely, if there exists $z \neq 0$ such that $G(s) z \rightarrow 0$ as $z \rightarrow \lambda$, then, necessarily,

$$
E(\lambda) \lim _{s \rightarrow \lambda} \Psi^{-1}(s) R(\lambda) z=0
$$

so that $\lambda$ is a root of at least one of the polynomials $\varepsilon_{i}(s), i=1, \cdots r$. The proof of the lemma in the converse case ( $p \leq m$ ) formally proceeds along the same route by replacing $G(s)$ with $G^{\prime}(s)$.

A quite different definition of transmission zeros and poles makes reference to the minors of $G(s)$. A $k$-degree minor of a matrix $A$ is the determinant of any square $k$-dimensional submatrix of A. It is possible to prove that the polynomial $\pi_{p}(s)$ of the poles of $G(s)$ is given by the least common denominator of all the non zero minors of any order of $G(s)$. Analogously, the polynomial $\pi_{z t}(s)$ of the transmission zeros of $G(s)$ is the greatest common divisor of all numerators of the minors of order $r$ of $G(s)$ provided that they have been adjusted so as to present the polynomial $\pi_{p}(s)$ as their denominator.

Remark 2.7 (Transmission zeros of a square system) In the particular case where $G(s)$ is square, it follows (recall Theorem 2.2)

$$
\operatorname{det}[G(s)]=\frac{\pi_{z t}(s)}{\pi_{p}(s)}
$$

Notice, however, that the presence of possible cancellations avoids in general to catch all transmission zeros and poles of $G(s)$ from its determinant.

The transmission poles and zeros of a system $\Sigma(A, B, C, D)$ enjoy an important input output characterization. As for the transmission zeros, reference is made to the so called blocking property, which deals with the possibility of getting identically zero forced output when the input is suitably chosen in the class of exponential and impulsive signals. Before formally stating the relevant result, it is advisable to stress that the transmission zeros of $G(s)$ and $G^{\prime}(s)$ actually coincides (recall the appropriate definition and Lemma 2.8). The same occurs for (transmission) poles. Let now indicate with $\delta(t):=\delta^{(0)}(t)$ the impulsive "function" and with $\delta^{(k)}(t)$ its $k$-th order derivative (recall that $\delta_{L}^{(k)}=s^{k}$ ). Moreover define

$$
\hat{\Sigma}:=\left\{\begin{array}{ll}
\Sigma & \text { if } p \geq m  \tag{2.12}\\
\Sigma^{\prime} & \text { if } p \leq m
\end{array} \quad, \quad \hat{G}(s):= \begin{cases}G(s) & \text { if } p \geq m \\
G^{\prime}(s) & \text { if } p \leq m\end{cases}\right.
$$

Theorem 2.3 (Time domain characterization of transmission zeros and poles) Let $G(s)$ be the transfer function of a system $\Sigma$. Then
i) The complex number $\lambda$ is a pole of $\Sigma$ if and only if there exists an impulsive input

$$
u(t)=\sum_{i=0}^{\nu} \alpha_{i} \delta^{(i)}(t)
$$

of $\hat{\Sigma}$, with $\alpha_{i}$ suitable constants, $\nu \geq 0$, such that the forced output $y_{f}(\cdot)$ of $\hat{\Sigma}$ is

$$
y_{f}(t)=y_{0} e^{\lambda t}, \forall t>0
$$

ii) The complex number $\lambda$ is a transmission zero of $\Sigma$ if and only if there exists an exponential/impulsive input

$$
u(t)=u_{0} e^{\lambda t}+\sum_{i=0}^{\nu} \alpha_{i} \delta^{(i)}(t)
$$

of $\hat{\Sigma}$, with $\alpha_{i}$ suitable constants, $\nu \geq 0, u_{0} \neq 0$, such that the forced output $y_{f}(\cdot)$ of $\hat{\Sigma}$ is

$$
y_{f}(t)=0, \quad \forall t>0
$$

Proof Consider the Smith-McMillan form $\hat{M}(s)$ of $\hat{G}(s)$, introduced in Theorem 2.2, so that $\hat{G}(s)=L(s) \hat{M}(s) R(s)$, where $\hat{M}(s)$ is the $\max [m, p] \times r$ matrix $\hat{M}(s)=$ $\left[\operatorname{diag}\left\{\varepsilon_{i}(s) / \psi_{i}(s)\right\} \quad 0\right]^{\prime}$. The polynomial matrices $L(s)$ and $R(s)$ are unimodular. Denote by $l_{i}(s)$ and $r_{i}^{\prime}(s)$ the $i$-th column and $i$-th row of $L(s)$ and $R(s)$, respectively. It turns out that

$$
y_{f L 0}=\sum_{i=1}^{r} l_{i}(s) \frac{\varepsilon_{i}(s)}{\psi_{i}(s)} r_{i}^{\prime}(s) u_{L}
$$

Point i) Assume that $\lambda$ is a pole of $\hat{G}(s)$, i.e. a root of the polynomial $\psi_{k}(s)$. Hence $\gamma(s):=(s-\lambda)^{-1} \psi_{k}(s)$ is a polynomial. Then the input $u_{L}$ defined as

$$
u_{L}:=R(s)^{-1}\left[\begin{array}{c}
0  \tag{2.13}\\
\vdots \\
0 \\
\gamma(s) \\
0 \\
\vdots \\
0
\end{array}\right] \longleftarrow k \text {-th row }
$$

is a polynomial vector since $R^{-1}(s)$ is a polynomial matrix. Therefore $r_{i}^{\prime}(s) u_{L}=$ $0, i \neq k$ and $r_{k}^{\prime}(s) u_{L}=\gamma(s)$. It then turns out that $y_{f L 0}=(s-\lambda)^{-1} l_{k}(s) \varepsilon_{k}(s)$. Since $\lambda$ is not a root of $\varepsilon_{k}(s)$,

$$
y_{f L 0}=y_{0}(s-\lambda)^{-1}+\beta(s)
$$

where $y_{0}$ is a suitable constant vector and $\beta(s)$ a suitable polynomial vector. Transforming back this expression in the time domain for $t>0$, the conclusion follows.

Conversely, assume that there exists an input, with polynomial Laplace transform $u_{L}$, such that the Laplace transform of the forced output is $y_{f L 0}=y_{0}(s-\lambda)^{-1}+\beta(s)$, being $\beta(s)$ a polynomial vector. Then

$$
y_{f L 0}=\sum_{i=1}^{r} l_{i}(s) \frac{\varepsilon_{i}(s)}{\psi_{i}(s)} r_{i}^{\prime}(s) u_{L}=y_{0}(s-\lambda)^{-1}+\beta(s)
$$

This means that at least one polynomial $\psi_{i}(s)$ must possess $\lambda$ as a root.
Point ii) Assume now that $\lambda$ is a zero of $\hat{G}(s)$, root of the polynomial $\varepsilon_{k}(s)$. Choose $u_{L}$ as in eq. (2.13) with $\gamma(s):=(s-\lambda)^{-1} \psi_{k}(s)$. Since $\lambda$ is not a root of $\psi_{k}(s)$ such an input matches the form given in the statement. Moreover $r_{i}^{\prime}(s) u_{L}=0, i \neq k$ and $r_{k}^{\prime}(s) u_{L}=(s-\lambda)^{-1} \psi_{k}(s)$, so that $y_{f L 0}=(s-\lambda)^{-1} l_{k}(s) \varepsilon_{k}(s)$ is a polynomial vector whose inverse Laplace transform is zero for strictly positive time instants.

Conversely assume that there exists an input of the form $u_{L}=u_{0}(s-\lambda)^{-1}+\beta(s)$, with $u_{0} \neq 0$ constant and $\beta(s)$ polynomial such that

$$
y_{f L 0}=\sum_{i=1}^{r} l_{i}(s) \frac{\varepsilon_{i}(s)}{\psi_{i}(s)} r_{i}^{\prime}(s)\left(u_{0}(s-\lambda)^{-1}+\beta(s)\right)
$$

is polynomial. A little thought shows that, besides other things, $y_{f L 0}$ may well be polynomial only if $\lambda$ is a root of at least one polynomial $\varepsilon_{i}(s)$.

The terminology adopted for the transmission zeros derives from the fact that they basically make reference to the transfer function (transmittance) of the system at hand. In the simple case of single input single output systems, the transfer function can be given the form

$$
G(s)=\frac{C \operatorname{adj}[s I-A] B}{\operatorname{det}[s I-A]}+D
$$

where $\operatorname{adj}[s I-A]$ is the matrix whose generic element $(i, j)$ is given by the determinant, multiplied by $(-1)^{i+j}$, of the matrix obtained by $(s I-A)$ ruling out its $j$-th row and $i$-th column. The transmission zeros coincide with the roots of the numerator once all the possible cancellations between the polynomial $C$ adj $[s I-A] B$ and the characteristic polynomial have been actually performed. As shown in the sequel, all the roots of $C \operatorname{adj}[s I-A] B+D \operatorname{det}[s I-A]$ are still properly called zeros of the system. These roots actually constitutes the so called invariant zeros. In the general multivariable framework, the definition of such zeros calls for the introduction of the polynomial matrix

$$
P(s):=\left[\begin{array}{cc}
s I-A & -B \\
C & D
\end{array}\right]
$$

which is referred to as system matrix.

Definition 2.9 (Invariant zeros) Consider system $\Sigma$ and let $P(s)$ be the associated system matrix with $v:=\operatorname{rank}[P(s)]$. Moreover, let $S(s)$ be the Smith form of $P(s)$, i.e.

$$
S(s)=\left[\begin{array}{ccccc}
\alpha_{1}(s) & 0 & \cdots & 0 & 0 \\
0 & \alpha_{2}(s) & \cdots & 0 & 0 \\
\vdots & \vdots & \ddots & \vdots & \vdots \\
0 & 0 & \cdots & \alpha_{v}(s) & 0 \\
0 & 0 & \cdots & 0 & 0
\end{array}\right]
$$

A complex number $\lambda$ is said to be an invariant zero of $\Sigma$ if it is a root of the polynomial

$$
\pi_{z i}(s):=\alpha_{1}(s) \alpha_{2}(s) \cdots \alpha_{v}(s)
$$

Remark 2.8 Notice that $v=\operatorname{rank}[P(s)]=n+\operatorname{rank}[G(s)]$. As a matter of fact, simple computations show that

$$
\begin{aligned}
P(s) & :=\left[\begin{array}{cc}
s I-A & -B \\
C & D
\end{array}\right] \\
& =\left[\begin{array}{cc}
s I-A & 0 \\
C & I
\end{array}\right]\left[\begin{array}{cc}
(s I-A)^{-1} & 0 \\
0 & G(s)
\end{array}\right]\left[\begin{array}{cc}
s I-A & -B \\
0 & I
\end{array}\right]
\end{aligned}
$$

Therefore, the claim on the rank of $P(s)$ is proved by noticing that in the right hand side of the above equation the first and last matrices above are nonsingular, while the remaining one has rank equal to $n+\operatorname{rank}[G(s)]$.

Like the transmission zeros, also the invariant zeros admit a rank characterization, which in this case concerns the kernel of either $P(s)$ or $P^{\prime}(s)$.

Lemma 2.9 (Rank property of invariant zeros) Let $P(s)$ be the system matrix of a system $\Sigma$ with transfer function $G(s)=C(s I-A)^{-1} B+D$ with $\operatorname{rank}[G(s)]=$ $\min [p, m]$. The complex number $\lambda$ is an invariant zero of the system if and only if $P(s)$ looses rank in $s=\lambda$, i.e. if and only if there exists a nonzero vector $z$ such that

$$
\begin{cases}P(\lambda) z=0 & \text { if } p \geq m \\ P^{\prime}(\lambda) z=0 & \text { if } p \leq m\end{cases}
$$

Proof Consider first the case $p \geq m$ and let $S(s)$ be the Smith form of $P(s)$, so that $P(s)=L(s) S(s) R(s)$, where $L(s)$ and $R(s)$ are suitable polynomial and unimodular matrices of dimensions $p+n$ and $n+m$, respectively, whereas $S(s)$ is as in Definition 2.9. Let $e_{k}(h)$ be the $k$-th column of the $h$-dimensional identity matrix. Recall also that matrix $R^{-1}(s)$ is polynomial and unimodular as well. Hence, if $\lambda$ is an invariant zero of the system, root of the polynomial $\alpha_{k}(s)$ in $S(s)$, then it is $P(\lambda) z=0$ with $z=R^{-1}(\lambda) e_{k}(n+m)$. Conversely, if there exists $z \neq 0$ such that $P(\lambda) z=0$, then, necessarily, $S(\lambda) R(\lambda) z=0$, which in turn implies that $\lambda$ is a root of at least one polynomial $\alpha_{i}(s)$, since $R(\lambda) z \neq 0$.

The proof in the case $p \leq m$ develops along the same lines, once $P(s)$ has been replaced by $P^{\prime}(s)$.

Example 2.13 Consider again the system defined in Example 2.12. It results

$$
P(s)=\left[\begin{array}{rcccr}
s & -1 & 0 & 0 & 0 \\
10 & s-7 & 0 & 0 & -1 \\
0 & 0 & s-5 & 0 & -2 \\
-1 & 1 & -1 & s-6 & 0 \\
-13 & 5 & 0 & 0 & 1 \\
0 & 0 & 1 & 0 & 1
\end{array}\right]
$$

As for the effective computation of the invariant zeros, one can actually utilize the result stated in Lemma 2.9, by looking for the vectors $z=\left[\begin{array}{llll}z_{1} & z_{2} & \cdots & z_{5}\end{array}\right]^{\prime}$ in the kernel of $P(\lambda)$, i.e. those vectors $z$ such that $P(\lambda) z=0$. A nonzero solution of the relevant equations can be found only if $\lambda=3$ or $\lambda=6$, which are therefore invariant zeros. Recall (Example 2.12) that only $\lambda=3$ is a transmission zero.

As apparent from their definition, the invariant zeros are not affected by a change of basis in the state space, as stated in the following lemma.

Lemma 2.10 (Invariant zeros vs. changes of basis) The set of the invariant zeros of a system $\Sigma$ is invariant with respect to a change of basis.

Proof If the triple $(\bar{A}, \bar{B}, \bar{C})$, with $\bar{A}=T A T^{-1}, \bar{B}=T B, \bar{C}=C T^{-1}$, describes, together with matrix $D$, system $\Sigma$ in a new basis, it follows

$$
\bar{P}(s)=\left[\begin{array}{cc}
s I-\bar{A} & -\bar{B} \\
\bar{C} & D
\end{array}\right]=\left[\begin{array}{ll}
T & 0 \\
0 & I
\end{array}\right]\left[\begin{array}{cc}
s I-A & -B \\
C & D
\end{array}\right]\left[\begin{array}{cc}
T^{-1} & 0 \\
0 & I
\end{array}\right]
$$

so that $\bar{P}(s)$ and $P(s)$ have the same Smith form. Hence both systems $\Sigma(A, B, C, D)$ and $\Sigma(\bar{A}, \bar{B}, \bar{C}, D)$ have the same invariant zeros.

Also the invariant zeros enjoy a blocking property, which stems on the existence of an exponential input yielding identically zero forced output.

Theorem 2.4 (Time domain characterization of invariant zeros) The complex number $\lambda$ is an invariant zero of $\Sigma$ if and only if at least one of the two following conditions holds
i) $\lambda$ is an eigenvalue of the unobservable part of $\hat{\Sigma}$;
ii) there exist two vectors $x_{0}$ and $u_{0} \neq 0$ such that the forced output of $\hat{\Sigma}$ corresponding to the input $u=u_{0} e^{\lambda t}, t \geq 0$ and initial state $x(0)=x_{0}$ is identically zero for $t \geq 0$.

Proof It is sufficient to prove the theorem in the case $p \geq m$, since the proof in the converse case easily follows by replacing $\Sigma$ with $\Sigma^{\prime}$. Hence assume $\hat{\Sigma}=\Sigma$.

Let now $\lambda$ be an invariant zero of $\Sigma$ and let $P(s)$ be the system matrix. Thanks to Lemma 2.9 there exists a non zero vector $z=\left[v^{\prime} w^{\prime}\right]^{\prime}$ such that $P(\lambda) z=0$, i.e.

$$
\begin{align*}
(\lambda I-A) v & =B w  \tag{2.14}\\
C v+D w & =0 \tag{2.15}
\end{align*}
$$

Letting $x(0):=v$ and $u_{0}:=w$, it is now verified that the input $u(t)=w e^{\lambda t}$ produces, together with the initial state $x(0)=v$, an identically zero output for $t \geq 0$. The Laplace transform of the input is $u_{L}=(s-\lambda)^{-1} w$ and that of the state is $x_{L}=$ $(s I-A)^{-1}\left[v+B w(s-\lambda)^{-1}\right]$. Eq. (2.14) entails $(s-\lambda)^{-1} B w=(s-\lambda)^{-1}(\lambda I-A) v$, so
that $x_{L}=(s-\lambda)^{-1} v$ and, from eq. (2.15), $\left[C x_{L}+D w(s-\lambda)^{-1}\right](s-\lambda)=y_{L}(s-\lambda)=0$. Since $y(0)=0$ (eq. (2.15)), the conclusion is drawn that $y(t)=0, t \geq 0$. In particular, if $w \neq 0$ then condition $i i$ ) is verified, whereas, if $w=0$, then eqs. (2.14),(2.15) entail, in view of the PBH test, that condition i) holds.

Conversely, assume without any loss of generality (recall Lemma 2.10), that the system at hand is from the very beginning in the standard Kalman canonical form for observability, i.e.

$$
A=\left[\begin{array}{cc}
A_{1} & 0 \\
A_{2} & A_{3}
\end{array}\right], \quad B=\left[\begin{array}{l}
B_{1} \\
B_{2}
\end{array}\right], \quad C=\left[\begin{array}{ll}
C_{1} & 0
\end{array}\right]
$$

where the pair $\left(A_{1}, C_{1}\right)$ is observable. If condition $i$ ) holds (namely $\lambda$ is an eigenvalue of $A_{3}$ ) choose $z=\left[0 \xi^{\prime} 0\right]^{\prime}$, where $\xi \neq 0$ is such that $\left(\lambda I-A_{3}\right) \xi=0$. Then, obviously, $P(\lambda) z=0$ so that $\lambda$ is an invariant zero of $\Sigma$. If condition $i i)$ holds, let, according to the structure of $A, x_{0}:=\left[x_{01}^{\prime} x_{02}^{\prime}\right]^{\prime}$. Being $y_{L}=0$, it follows

$$
C_{1}\left(s I-A_{1}\right)^{-1}\left[x_{01}+\frac{B_{1} u_{0}}{(s-\lambda)}\right]+\frac{D u_{0}}{(s-\lambda)}=0
$$

By noticing that

$$
\begin{aligned}
\left(s I-A_{1}\right)^{-1} x_{01} & =\frac{I-\left(s I-A_{1}\right)^{-1}\left(\lambda I-A_{1}\right)}{(s-\lambda)} x_{01} \\
y(0) & =C_{1} x_{01}+D u_{0}
\end{aligned}
$$

it then follows

$$
C_{1}\left(s I-A_{1}\right)^{-1}\left[\left(\lambda I-A_{1}\right) x_{01}-B_{1} u_{0}\right]=0
$$

The first term of such an equation is the Laplace transform of the (free) output of the system $\Sigma\left(A_{1}, 0, C_{1}, 0\right)$ when the initial state is $\left(\lambda I-A_{1}\right) x_{01}-B_{1} u_{0}$. Since this system is observable it follows that $\left(\lambda I-A_{1}\right) x_{01}-B_{1} u_{0}=0$. Choose

$$
z:=\left[\begin{array}{lll}
x_{01}^{\prime} & \xi^{\prime} & u_{0}^{\prime}
\end{array}\right]
$$

where $\xi:=\left(\lambda I-A_{3}\right)^{-1}\left(A_{2} x_{01}+B_{2} u_{0}\right)$. Obviously, $P(\lambda) z=0$ so that $\lambda$ is an invariant zero of $\Sigma$.

The theorem above points out the circumstances under which the output of $\hat{\Sigma}$ is zero for all $t \geq 0$. Actually, a part from the trivial case of zero initial state and input, the output of $\hat{\Sigma}$ can be such if and only if $\Sigma$ possesses invariant zeros. As already said, for SISO systems the invariant zeros are the roots of $C \operatorname{adj}\left[(s I-A)^{-1}\right] B+$ $D \operatorname{det}[(s I-A)]$ whereas (in general) only a part of these roots are transmission zeros. This relationship holds for MIMO systems as well.

Theorem 2.5 (Invariant vs. transmission zeros) A transmission zero of a system $\Sigma$ is also an invariant zero of $\Sigma$.

Proof Consider first the case $p \geq m$ and let $\lambda$ be a transmission zero of $\Sigma$. Thanks to Theorem 2.3 there exists an exponential/impulsive input

$$
u(t)=u_{0} e^{\lambda t}+\sum_{i=0}^{\nu} \alpha_{i} \delta^{(i)}(t)
$$

with $\alpha_{i}$ and $u_{0} \neq 0$ suitable constants, such that the forced output of $\Sigma$ is, $\forall t>0$,

$$
y_{f}(t)=C \int_{0}^{t} e^{A(t-\tau)} B\left[u_{0} e^{\lambda \tau}+\sum_{i=0}^{\nu} \alpha_{i} \delta^{(i)}(\tau)\right] d \tau+D u_{0} e^{\lambda t}=0
$$

Letting

$$
x_{0}:=\left[\begin{array}{llll}
B & A B & \cdots & A^{\nu} B
\end{array}\right]\left[\begin{array}{c}
\alpha_{0} \\
\alpha_{1} \\
\vdots \\
\alpha_{\nu}
\end{array}\right]
$$

it follows that

$$
C e^{A t} x_{0}=C \int_{0}^{t} e^{A(t-\tau)} B \sum_{i=0}^{\nu} \alpha_{i} \delta^{(i)}(\tau) d \tau
$$

so that

$$
y_{f}(t)=C e^{A t} x_{0}+\int_{0}^{t} e^{A(t-\tau)} B u_{0} e^{\lambda \tau} d \tau+D u_{0} e^{\lambda t}=0, \quad \forall t>0
$$

This last expression coincides with the output response of system $\Sigma$ when the initial state is $x(0)=x_{0}$ and the input is $u(t)=u_{0} e^{\lambda t}$. Such response is obviously continuous from the right, so that the output is zero at $t=0$ as well. Theorem 2.4 ensures that $\lambda$ is an invariant zero of $\Sigma$.

The proof of the Theorem in the case $p \leq m$ can be derived in complete analogy by considering system $\Sigma^{\prime}$ instead of $\Sigma$.

The following result clarifies and completes the relationships between transmission and invariant zeros.

Theorem 2.6 ( Transmission vs. invariant zeros of a system in minimal form ) The transmission and invariant zeros of a reachable and observable system do coincide.

Proof As already seen in Theorem 2.5, a transmission zero is also an invariant zero. It is then left to show the converse statement when $(A, B)$ is reachable and $(A, C)$ is observable. Consider the case $p \geq m$ since the other case is easily proved by transposition. Let $\lambda$ be an invariant zero. Thanks to Lemma 2.9 there exists a nonzero vector $z=\left[v^{\prime} w^{\prime}\right]^{\prime}$ such that $P(\lambda) z=0$. Notice that if $w=0$ and $v \neq 0$, then this condition implies that $A v=\lambda v$ and $C v=0$, contrary to the observability assumption of $(A, C)$ (recall Lemma D.1). Hence $w \neq 0$. Moreover, thanks to Theorem 2.4, the system response when the initial state is $x(0)=v$ and the input is $u(t)=w e^{\lambda t}$ is identically zero, i.e.

$$
\begin{equation*}
C\left[e^{A t} v+\int_{0}^{t} e^{A(t-\tau)} B w e^{\lambda \tau} d \tau\right]+D w e^{\lambda t}=0, \quad \forall t \geq 0 \tag{2.16}
\end{equation*}
$$

Recalling Theorem 2.3, $\lambda$ is a transmission zero if there exists an input of the form

$$
u(t)=u_{0} e^{\lambda t}+\sum_{i=0}^{\nu} \alpha_{i} \delta^{(i)}(t)
$$

which yields an identically zero forced output (for $t>0$ ), i.e. if

$$
C \int_{0}^{t} e^{A(t-\tau)} B\left[u_{0} e^{\lambda \tau}+\sum_{i=0}^{\nu} \alpha_{i} \delta^{(i)}(\tau)\right] d \tau+D u_{0} e^{\lambda t}=0, \quad \forall t>0
$$

Hence the theorem is proved if one shows that there exist real coefficients $\alpha_{i}$ such that

$$
\begin{align*}
C e^{A t} v & =\int_{0}^{t} C e^{A(t-\tau)} B \sum_{i=0}^{\nu} \alpha_{i} \delta^{(i)}(\tau) d \tau \\
& =C e^{A t} \sum_{i=0}^{\nu} A^{i} B \alpha_{i}, \quad \forall t>0 \tag{2.17}
\end{align*}
$$

Any vector $\alpha:=\left[\begin{array}{llll}\alpha_{0}^{\prime} & \alpha_{1}^{\prime} & \cdots & \alpha_{\nu}^{\prime}\end{array}\right]^{\prime}$ such that

$$
\left[\begin{array}{llll}
B & A B & \cdots & A B^{\nu}
\end{array}\right] \alpha=v
$$

satisfies eq. (2.17). Notice that a vector $\alpha$ exists corresponding to $\nu=n-1$ since $(A, B)$ is reachable so that the Grammian matrix $\left[B A B \cdots A B^{n-1}\right]$ has full row rank.

The invariant and transmission zeros do not exhaust the totality of zeros which can be defined for a system. Actually, consider an unobservable system with $p<m$. It may well happen that an eigenvalue of the unobservable part, say $\lambda$, is such that $P(\lambda)$ does not loose rank. Associated with such an eigenvalue there exists an eigenvector (initial state) $x(0)$ such that the free motion of the output $y(\cdot)$ is identically zero. Therefore the complex number $\lambda$ can be still considered as a zero of the system, whose nature is different from that of the zeros previously introduced. In complete analogy, an unreachable system $\Sigma$, with $p>m$, can admit an eigenvalue of the unreachable part, say $\lambda$, which is such that the associated system matrix $P(\lambda)$ does not loose rank. Hence $\lambda$ is not an invariant zero. However, it is well known that there exists an initial state for $\Sigma^{\prime}$ (eigenvector associated with $\lambda$ ) capable of zeroing the free output of $\Sigma^{\prime}$. Again, $\lambda$ can be fairly considered as a zero of system $\Sigma$. Such zeros will be referred to as decoupling zeros.

Definition 2.10 (Output decoupling zeros) Consider $\Sigma(A, B, C, D)$ a n-dimensional system and the polynomial matrix

$$
P_{C}(s):=\left[\begin{array}{c}
s I-A \\
C
\end{array}\right]
$$

Let

$$
S_{C}(s)=\left[\begin{array}{c}
\operatorname{diag}\left\{\alpha_{i}^{C}(s)\right\} \\
0
\end{array}\right]
$$

be the Smith form of $P_{C}(s)$. A complex number $\lambda$ is said to be an output decoupling zero if it is a root of the polynomial

$$
\alpha^{C}(s):=\alpha_{1}^{C}(s) \alpha_{2}^{C}(s) \cdots \alpha_{n}^{C}(s)
$$

Definition 2.11 (Input decoupling zeros) Consider $\Sigma(A, B, C, D)$ a n-dimensional system and the polynomial matrix

$$
P_{B}(s):=\left[\begin{array}{ll}
s I-A & -B
\end{array}\right]
$$

Let

$$
S_{B}(s)=\left[\operatorname{diag}\left\{\alpha_{i}^{B}(s)\right\} \quad 0\right]
$$

be the Smith form of $P_{B}(s)$. A complex number $\lambda$ is said to be an input decoupling zero if it is a root of the polynomial

$$
\alpha^{B}(s):=\alpha_{1}^{B}(s) \alpha_{2}^{B}(s) \cdots \alpha_{n}^{B}(s)
$$

Definition 2.12 (Input-output decoupling zeros) Consider $\Sigma(A, B, C, D)$ a system, its associated polynomial matrices $P_{C}(s), P_{B}(s)$ with their Smith forms $S_{C}(s)$ and $S_{B}(s)$, and the polynomials $\alpha^{C}(s)$ and $\alpha^{B}(s)$, respectively. A complex number $\lambda$ is said to be an input-output decoupling zero if it is a root of both polynomials $\alpha^{C}(s)$ and $\alpha^{B}(s)$.

The decoupling zeros are not affected by a change of basis in the state-space of the system as it can be checked by resorting to the same arguments exploited in the proof of Lemma 2.10. Further, they can be characterized in terms of the kernels of $P_{C}(\lambda)$ and $P_{B}^{\prime}(\lambda)$. The relevant results are presented in the following lemmas given without proofs since completely similar to that of Lemma 2.9.

Lemma 2.11 (Rank property of the output decoupling zeros) A complex number $\lambda$ is an output decoupling zero if and only if there exists $z \neq 0$ such that

$$
P_{C}(\lambda) z=0
$$

Lemma 2.12 (Rank property of the input decoupling zeros) A complex number $\lambda$ is an input decoupling zero if and only if there exists $w \neq 0$ such that

$$
P_{B}^{\prime}(\lambda) w=0
$$

Lemma 2.13 (Rank property of the input-output decoupling zeros) A complex number $\lambda$ is an input-output decoupling zero if and only if there exist $z \neq 0$ and $w \neq 0$ such that

$$
\begin{aligned}
P_{C}(\lambda) z & =0 \\
P_{B}^{\prime}(\lambda) w & =0
\end{aligned}
$$

In Tables 2.1 and 2.2 the definitions and basic properties of the zeros introduced so far are schematically illustrated.

Remark 2.9 Based on Lemmas 2.11-2.13, and on the PBH tests relative to observability and reachability (Lemmas D.1- D.2), it is straightforward to realize that a system in minimal form does not possess decoupling zeros.

It is worth pointing out that the given definitions put in relief possible relations between invariant and decoupling zeros. In fact, if the number of inputs does not exceed the number of outputs, it is immediately seen that the output decoupling zeros are invariant zeros as well. Analogously, in the converse case, i.e. when the number of outputs is not greater than the number of inputs, the input decoupling zeros are particular invariant zeros. However, as shown in the example below, there may well happen that a system has decoupling zeros which are not invariant.

| Type | Definition | Rank property |
| :--- | :--- | :--- |
| Transmission | $\pi_{z t}(\lambda)=0$ | $\lim _{s \rightarrow \lambda} \quad G(s) z=0 \quad$ if $p \geq m$ <br> $\lim _{s \rightarrow \lambda}$$\quad G^{\prime}(s) z=0 \quad$ if $p \leq m$ |
| Invariant | $\pi_{z i}(\lambda)=0$ | $P(\lambda) z=0 \quad$ if $p \geq m$ <br> $P^{\prime}(\lambda) z=0 \quad$ if $p \leq m$ |
| Output <br> decoupling | $\alpha^{C}(\lambda)=0$ | $P_{C}(\lambda) z=0$ |
| Input <br> decoupling | $\alpha^{B}(\lambda)=0$ | $P_{B}^{\prime}(\lambda) z=0$ |
| Input-output <br> decoupling | $\alpha^{C}(\lambda)=0$ <br> $\alpha^{B}(\lambda)=0$ | $P_{C}(\lambda) z=0$ <br> $P_{B}^{\prime}(\lambda) z=0$ |

Table 2.1: Rank properties of the zeros

| Type | Definition | Output property |
| :---: | :---: | :---: |
| Transmission | $\pi_{z t}(\lambda)=0$ | $\begin{aligned} & x(0)=0, \exists u(\cdot) \\ & y_{f}(t)=0, \forall t>0 \end{aligned} \quad \text { on } \hat{\Sigma}$ |
| Invariant | $\pi_{z i}(\lambda)=0$ | $\begin{aligned} & x(0)=0, \exists u(\cdot) \\ & y(t)=0, \forall t \geq 0 \end{aligned} \quad \text { on } \hat{\Sigma}$ |
| Output decoupling | $\alpha^{C}(\lambda)=0$ | $\begin{aligned} & \exists x(0) \neq 0, u(\cdot)=0 \\ & y(t)=0, \forall t \geq 0 \end{aligned} \quad \text { on } \Sigma$ |
| Input decoupling | $\alpha^{B}(\lambda)=0$ | $\begin{aligned} & \exists x(0) \neq 0, u(\cdot)=0 \quad \text { on } \Sigma^{\prime} \\ & y(t)=0, \forall t \geq 0 \end{aligned}$ |
| Input - output decoupling | $\begin{aligned} & \alpha^{C}(\lambda)=0 \\ & \alpha^{B}(\lambda)=0 \\ & \hline \end{aligned}$ | $\begin{aligned} & \exists x(0) \neq 0, u(\cdot)=0 \quad \text { on } \quad \Sigma \\ & y(t)=0, \forall t \geq 0 \end{aligned} \quad \begin{aligned} & \Sigma^{\prime} \\ & \hline \end{aligned}$ |

Table 2.2: Output properties of the zeros

Example 2.14 Consider again the system defined in Example 2.13. It is obvious that the invariant zero $\lambda=6$ is also an output decoupling zero. However, there exists an input decoupling zero, $\lambda=5$, which is not invariant. Actually, matrix $P(5)$ has full rank (equal to five), even though the first four rows are linearly dependent, so that $P_{B}(s)=[s I-A-B]$ looses rank for $s=5$.

It should be now evident the relation existing between invariant and decoupling zeros when the system at hand is square.
Lemma 2.14 (Decoupling vs. invariant zeros for square systems) Consider a system with the same number of inputs and outputs (square system). Then the set of decoupling zeros is a subset of the set of invariant zeros.

Proof If $\lambda$ is a decoupling zero, one or both of the two matrices $P_{C}(s)$ and $P_{B}(s)$ must loose rank for $s=\lambda$. Hence matrix $P(s)$ looses rank in $s=\lambda$ as well.

In case of nonminimal systems, the set of invariant zeros does not coincide with that of transmission zeros. Moreover, there may be decoupling zeros which are not
invariant. These fact motivates the definition below.
Definition 2.13 (System zeros) The set $\mathcal{Z}_{s}$ of system zeros is defined as

$$
\mathcal{Z}_{s}=\mathcal{Z}_{t} \cup \mathcal{Z}_{i-o} \cup \mathcal{Z}_{o-i} \cup \mathcal{Z}_{i o}
$$

where $\mathcal{Z}_{t}$ is the set of transmission zeros, $\mathcal{Z}_{i-o}$ is the set of input decoupling zeros which are not also output decoupling, $\mathcal{Z}_{o-i}$ is the set of output decoupling zeros which are not also input decoupling, and $\mathcal{Z}_{i o}$ is the set of input-output decoupling zeros.

Lemma 2.15 The set of invariant zeros of a system $\Sigma(A, B, C, D)$ is a subset of the set of the system zeros.

Proof The lemma will be proved in the case $p \geq m$, as the converse case being easily handled by transposition. One has to show that if $\lambda$ is an invariant zero, it is also either a transmission or a decoupling zero. Without any loss of generality, assume that the system is decomposed accordingly to the Kalman canonical decomposition, i.e.

$$
A=\left[\begin{array}{cccc}
A_{1} & A_{2} & A_{3} & A_{4} \\
0 & A_{5} & 0 & A_{6} \\
0 & 0 & A_{7} & A_{8} \\
0 & 0 & 0 & A_{9}
\end{array}\right], \quad B=\left[\begin{array}{c}
B_{1} \\
B_{2} \\
0 \\
0
\end{array}\right], \quad C=\left[\begin{array}{llll}
0 & C_{1} & 0 & C_{2}
\end{array}\right]
$$

where $\Sigma\left(A_{5}, B_{2}, C_{1}, D\right)$ constitutes a subsystem which is completely observable and reachable. Of course, if $\lambda$ belongs to the unreachable and/or unobservable part of the system, then it is a decoupling zero, so that the proof would be over.

Let $\lambda$ be an invariant zero and assume that it is not an eigenvalue of any of the matrices $A_{1}, A_{7}, A_{9}$. Then, there exists $z \neq 0$ such that $P(\lambda) z=0$, with $z:=\left[\begin{array}{lllll}z_{1}^{\prime} & z_{2}^{\prime} & z_{3}^{\prime} & z_{4}^{\prime} & u_{0}^{\prime}\end{array}\right]^{\prime}$. Since $\lambda$ is not an eigenvalue of $A_{9}$, it turns out that $z_{4}=0$ so that

$$
\left(\lambda I-A_{5}\right)^{-1} z_{2}-B_{2} u_{0}=0, \quad C_{1} z_{2}+D u_{0}=0
$$

Hence $\lambda$ is a transmission zero of the reachable and observable part $\Sigma\left(A_{5}, B_{2}, C_{1}, D\right)$ of the system, and hence of the system itself, provided that $\left[z_{2}^{\prime} u_{0}^{\prime}\right]^{\prime} \neq 0$. If it were not so, it would happen that $\left(\lambda I-A_{1}\right) z_{1}-A_{3} z_{3}=0$ and $\left(\lambda I-A_{7}\right) z_{3}=0$. Since $\lambda$ is not an eigenvalue neither of $A_{7}$ nor of $A_{1}$, it would follow $z_{1}=0$ and $z_{3}=0$ so that $z=0$, a contradiction.

Remark 2.10 In view of Lemma 2.15 and Theorem 2.4 it is easy to conclude that the set of systems zeros coincides with that of invariant zeros, relative to square systems. If in addition the system is in minimal form, the three sets (of transmission, invariant and system zeros) do actually coincide.

At the light of what has previously been said, the set of system zeros constitutes the totality of the zeros defined till now. It can be expressively partitioned in its subsets, as shown in fig. 2.5 with reference to the case $m>p, m=p$, and $m<p$. In the figure, the symbols $t, i, d_{i}, d_{o}, d_{i o}$ denote, respectively, the transmission, invariant, input decoupling, output decoupling and input-output decoupling zeros. The presence of one of this symbols in a part of the figure indicates that this part is contained in the set of the zeros under consideration.

This section ends with a brief discussion on the concept of inverse system in the simple case where $p=m$ and $\operatorname{det}[D] \neq 0$. For the more general case, the reader is referred to specialized texts.
$m>p$

| $t, i$ | $i, d_{i}$ |
| :--- | :---: |
|  |  |
|  | $i, d_{i o}$ |
|  | $i, d_{o}$ |
| $d_{o}$ |  |

\[

\]

\[

\]

Figure 2.1: The zeros of a system

Definition 2.14 (Inverse system) Consider a square system $\Sigma(A, B, C, D)$ such that $\operatorname{det}[D] \neq 0$. Then, the inverse system is

$$
\Sigma_{i n v}:=\left[\begin{array}{c|c}
F & G \\
\hline H & E
\end{array}\right]
$$

where $F:=A-B D^{-1} C, G:=B D^{-1}, H:=-D^{-1} C, E:=D^{-1}$.
The reason why system $\Sigma_{i n v}$ is called the inverse system can be simply explained as follows. Let $\Sigma$ and $\Sigma_{i n v}$ be described by

$$
\begin{aligned}
\dot{x} & =A x+B u \\
y & =C x+D u
\end{aligned}
$$

and

$$
\begin{aligned}
& \dot{\xi}=\left(A-B D^{-1} C\right) \xi+B D^{-1} \nu \\
& \theta=-D^{-1} C \xi+D^{-1} \nu
\end{aligned}
$$

respectively. Now, build up the series connection of the systems, according to the following two cases:
i) $\nu=y$, so that $\Sigma_{i n v}$ follows $\Sigma$ (system $\Sigma_{i n v} \Sigma$ )
ii) $u=\theta$, so that $\Sigma_{i n v}$ precedes $\Sigma$ (system $\Sigma \Sigma_{i n v}$ )

In the first case, letting $z:=\xi-x$ it follows

$$
\begin{aligned}
& \dot{z}=\left(A-B D^{-1} C\right) z \\
& \theta=-D^{-1} C z+u
\end{aligned}
$$

whereas in the second

$$
\begin{aligned}
& \dot{z}=A z \\
& y=-C z+\nu
\end{aligned}
$$

In both cases, the transfer function of the resulting system is the identity. Hence, if $G(s)$ is the transfer function of $\Sigma$, the transfer function of $\Sigma_{i n v}$ is exactly $G^{-1}(s)$.
Theorem 2.7 Consider a square system $\Sigma(A, B, C, D)$ where $D$ is nonsingular and let $\Sigma_{i n v}$ be its associated inverse system. Then:
i) The set of transmission zeros of $\Sigma$ coincides with the set of poles of $\Sigma_{i n v}$ and, conversely, the set of poles of $\Sigma$ coincides with the set of transmission zeros of $\Sigma_{i n v}$.
ii) The set of eigenvalues of $\Sigma$ coincides with the set of invariant zeros of $\Sigma_{i n v}$ and, conversely, the set of invariant zeros of $\Sigma$ coincides with the set of eigenvalues of $\Sigma_{i n v}$.

Proof Point i) It suffices to verify that, if $M(s)$ is the Smith-McMillan form of the transfer function of $\Sigma$, then $T M^{-1}(s) T$ is the Smith-McMillan form of the transfer function of $\Sigma_{i n v}$, where

$$
T:=\left[\begin{array}{ccccc}
0 & 0 & \cdots & 0 & 1 \\
0 & 0 & \cdots & 1 & 0 \\
\vdots & \vdots & \ddots & \vdots & \vdots \\
0 & 1 & \cdots & 0 & 0 \\
1 & 0 & \cdots & 0 & 0
\end{array}\right]
$$

The conclusion then follows from the definition of poles and transmission zeros (Definitions 2.7 and 2.8).

Point ii) Assume that $\lambda$ is an invariant zero of $\Sigma_{i n v}$, i.e. $\left(\lambda I-\left(A-B D^{-1} C\right)\right) w_{1}-$ $B D^{-1} w_{2}=0$ and $-D^{-1} C w_{1}+D^{-1} w_{2}=0$, with $w:=\left[w_{1}^{\prime} w_{2}^{\prime}\right]^{\prime} \neq 0$. These equations imply that $w_{1} \neq 0$ and $A w_{1}=\lambda w_{1}$, so that $\lambda$ is an eigenvalue of $\Sigma$. Also these considerations can be easily reversed. Finally, assume that $\lambda$ is an invariant zero of $\Sigma$, i.e. $(\lambda I-A) w_{1}-B w_{2}=0$ and $C w_{1}+D w_{2}=0$ with $w:=\left[w_{1}^{\prime} w_{2}^{\prime}\right]^{\prime} \neq 0$. Hence $w_{2}=-D^{-1} C w_{1}, w_{1} \neq 0$, so that $\left(A-B D^{-1} C\right) w_{1}=\lambda w_{1}$ implies that $\lambda$ is an eigenvalue of $\Sigma_{i n v}$. Reversing the procedure proves the validity of the converse statement as well.

### 2.6 Singular values

In this section, some of the most significant properties of the singular values of a matrix are reported along with the so-called singular value decomposition. Reference is made to constant vectors and matrices with complex elements. Recall that the symbol " $\sim$ " denotes the operation of conjugate transposition. The norm adopted for vectors is the one induced by the usual inner product in $C^{n}$.

Definition 2.15 (Singular values of a matrix) Let $A$ be a $n \times m$ complex matrix. The square roots $\sigma_{i}(A), i=1, \cdots m$ of the eigenvalues of $A^{\sim} A$ are called singular values of $A$.
Notice that the singular values of $A$ are real and nonnegative, since $A^{\sim} A$ is a hermitian and positive semidefinite matrix. Of course, $m-\operatorname{rank}[A]$ singular values of $A$ are in fact zero. The singular values of $A$ can be put into evidence by the so-called singular value decomposition of $A$, whose existence is established in the theorem below. Its proof, which is reported in Appendix B, is based on a few preliminarily steps that provide a way, although not optimal for the computational burden, to determine such a decomposition.

Theorem 2.8 (Singular value decomposition) Let $A$ be a $n \times m$ matrix such that $\operatorname{rank}[A]=k$. Then, there exist two unitary matrices $U$ and $V$ such that

$$
U^{\sim} A V=S
$$

where the only nonzero elements of $S$ are those in positions $(i, i), i=1, \cdots k$. Such elements are positive and nonincreasing.

The elements of $S$ in positions $i=1, \cdots, \min [n, m]$, are the singular values of $A$. Further,

$$
A^{\sim} A=V S^{\sim} U^{\sim} U S V^{\sim}=V S^{\sim} S V^{\sim}=V D V^{\sim}
$$

where $D=\operatorname{diag}\left\{\sigma_{i}^{2}(A), i=1, \cdots, m\right\}$. Of course, if $m>n$, matrix $A$ possesses at least $m-n$ singular values at the origin. The greatest and least singular values of $A$ are indicated with $\bar{\sigma}(A)$ and $\underline{\sigma}(A)$, respectively. The unitary matrices $U$ and $V$ specify the singular value decomposition of $A$.

Remark 2.11 Based on Theorem 2.8, matrix $S$ has the same dimensions as $A$ and exhibits the following structure

$$
S=\left\{\begin{array}{ccc}
{\left[\begin{array}{cc}
\Delta & 0
\end{array}\right]} & \text { if } & n<m \\
\Delta & \text { if } & n=m \\
{\left[\begin{array}{c}
\Delta \\
0
\end{array}\right]} & \text { if } & n>m
\end{array}\right.
$$

The "meaningful" part of $S$ is therefore constituted by the diagonal matrix $\Delta$ with dimension $\min [n, m]$. For such a reason, the singular value decomposition is sometimes presented in a different way, distinguishing two different situations:

1) Case $n>m$ : There exist two matrices $U_{1}$ and $V$ of dimensions $n \times m$ and $m \times m$, respectively, such that $A=U_{1} \Delta V^{\sim}$, with $V^{\sim} V=V V^{\sim}=I$ and $U_{1}^{\sim} U_{1}=I$. The diagonal matrix $\Delta$ is $m$-dimensional and contains the singular values of $A$.
2) Case $n<m$ : There exist two matrices $U$ and $V_{1}$ of dimensions $n \times n$ and $m \times n$, respectively, such that $A=U \Delta V_{1}^{\sim}$, with $U^{\sim} U=U U^{\sim}=I$ and $V_{1} V_{1}^{\sim}=I$. The diagonal matrix $\Delta$ is $n$-dimensional and contains the singular values of $A$.
With reference to what said in Theorem 2.8, it is immediate to check that $U=\left[\begin{array}{ll}U_{1} & U_{2}\end{array}\right]$, $S=[\Delta 0]^{\prime}$, if $n>m$ whereas $V=\left[V_{1} V_{2}\right], S=[\Delta 0]$, if $n<m$.

Remark 2.12 The pseudoinverse of $A$ can be found in a very simple way once the singular value decomposition of $A$ is known. Actually, let

$$
U^{\sim} A V=S=\left[\begin{array}{ll}
\Sigma & 0 \\
0 & 0
\end{array}\right]
$$

with $\Sigma$ square and nonsingular. Letting

$$
B:=V\left[\begin{array}{cc}
\Sigma^{-1} & 0 \\
0 & 0
\end{array}\right] U^{\sim}:=V T U^{\sim}
$$

it follows that

$$
\begin{aligned}
& B A B=V T U^{\sim} U S V^{\sim} V T U^{\sim}=V T S T U^{\sim}=V T U^{\sim}=B \\
& A B A=U S V^{\sim} V T U^{\sim} U S V^{\sim}=U S T S V^{\sim}=U S V^{\sim}=A
\end{aligned}
$$

These relations show that $B$ is the pseudoinverse of $A$.

The most significant properties of the singular values of a matrix are reported in the following lemmas. The symbol $\|x\|$ indicates the usual norm induced by the inner product in $C^{n}$, i.e. $\|x\|^{2}=x^{\sim} x$ and recall that $\lambda_{i}(A)$ is the $i$-th eigenvalue of $A$ and $r_{s}(A)$ is the spectral radius of $A$. The proof of the following lemmas are reported in Appendix B.

Lemma 2.16 Given a matrix A, the following relations hold:
1)

$$
\bar{\sigma}(A)=\max _{x \neq 0} \frac{\|A x\|}{\|x\|}
$$

2) 

$$
\underline{\sigma}(A)=\min _{x \neq 0} \frac{\|A x\|}{\|x\|}
$$

Lemma 2.17 Let $A$ be a matrix. Then $\bar{\sigma}(A)=\bar{\sigma}\left(A^{\sim}\right)$
Lemma 2.18 Let $A$ be a square matrix. Then
1)

$$
\underline{\sigma}(A) \leq\left|\lambda_{i}(A)\right| \leq \bar{\sigma}(A)
$$

2) 

$$
r_{s}(A) \leq \bar{\sigma}(A)
$$

Moreover, if $A$ is nonsingular, then
3)

$$
\underline{\sigma}(A)=\frac{1}{\bar{\sigma}\left(A^{-1}\right)}
$$

4) 

$$
\bar{\sigma}(A)=\frac{1}{\underline{\sigma}\left(A^{-1}\right)}
$$

Finally, if $A$ is hermitian, then
5)

$$
r_{s}(A)=\bar{\sigma}(A)
$$

Lemma 2.19 Let $\alpha \in C$ be an arbitrary scalar and $A$ a matrix. Then,

$$
\sigma_{i}(\alpha A)=|\alpha| \sigma_{i}(A)
$$

Lemma 2.20 Let $A$ and $B$ be two matrices with the same dimensions. Then,

$$
\bar{\sigma}(A+B) \leq \bar{\sigma}(A)+\bar{\sigma}(B)
$$

Lemma 2.21 Let $A$ and $B$ be two matrices such that $A B$ makes sense. Then,

$$
\bar{\sigma}(A B) \leq \bar{\sigma}(A) \bar{\sigma}(B)
$$

Lemma 2.22 Let $A$ and $B$ be two matrices with the same dimensions. Then,

$$
\underline{\sigma}(A)-\bar{\sigma}(B) \leq \underline{\sigma}(A+B) \leq \underline{\sigma}(A)+\bar{\sigma}(B)
$$

Lemma 2.23 Let $A$ and $B$ be two matrices with the same number of rows. Then,

$$
\max [\bar{\sigma}(A), \bar{\sigma}(B)] \leq \bar{\sigma}\left(\left[\begin{array}{ll}
A & B
\end{array}\right]\right) \leq \sqrt{2} \max [\bar{\sigma}(A), \bar{\sigma}(B)]
$$

Lemma 2.24 Let $A$ be a square matrix. Then,

$$
\sum_{i} \sigma_{i}^{2}(A)=\operatorname{trace}\left[A^{\sim} A\right]
$$

The quantity $\|A\|_{F}:=\sqrt{\operatorname{trace}\left(A^{\sim} \bar{A}\right)}$ is the so called Frobenius norm of $A$.
Lemma 2.25 Let $m$ be the number of columns of a matrix $A$ and denote by $A_{i j}$ its element in position $(i, j)$. Then,

$$
\max _{i, j}\left|A_{i j}\right| \leq \bar{\sigma}(A) \leq m \max _{i, j}\left|A_{i j}\right|
$$

### 2.7 Basic facts on linear operators

In this section some facts on the theory of linear operators are recalled. Since no confusion can arise in the present context, the term linear will be often disregarded.

Definition 2.16 (Operator) A (linear) operator $T$ is a linear map acting between two linear spaces.

For an operator $T: X \rightarrow Y$ acting between two normed linear spaces $X$ and $Y$, it is possible to introduce the notion of boundedness, norm and continuity.

Definition 2.17 (Boundedness of an operator) Let $X$ and $Y$ be two normed linear spaces and $T: X \rightarrow Y$ an operator. This operator is said to be bounded if there exists a constant $M$ such that

$$
\begin{equation*}
\|T x\| \leq M\|x\|, \forall x \in X \tag{2.18}
\end{equation*}
$$

Definition 2.18 Let $X$ and $Y$ be two normed linear spaces and $T: X \rightarrow Y$ a bounded operator. The smallest constant $M$ satisfying eq. (2.18) is called norm of $T$ and denoted by $\|T\|$.

Remark 2.13 Thanks to linearity and Definitions 2.17, 2.18, the norm of an operator can be significantly characterized as follows

$$
\|T\|=\sup _{x \neq 0} \frac{\|T x\|}{\|x\|}=\sup _{\|x\|=1}\|T x\|
$$

In the particular case where $T$ is actually a complex matrix, it follows, in view of Lemma 2.16 , that $\|T\|=\bar{\sigma}(T)$.

Definition 2.19 (Continuous operator) Let $X$ and $Y$ be two normed linear spaces and $T: X \rightarrow Y$ an operator. If, chosen a vector $\xi \in X$, for any $\varepsilon>0$ there exists $a$ $\delta>0$ such that

$$
\|x-\xi\|<\delta \Longrightarrow\|T x-T \xi\|<\varepsilon
$$

then $T$ is said to be continuous.

The concept of rank, which is well understood for operators defined in finite dimensional spaces, can be extended to the general case as follows.

Definition 2.20 (Rank of an operator) Let $X$ and $Y$ be two normed linear spaces and $T: X \rightarrow Y$ an operator. The rank of $T$ is the dimension of the closure of the range of $T$.

Also the concept of eigenvalue and eigenvector, which are familiar in the context of linear operators in finite dimensional spaces (matrices), can be extended in an analogous way.

Definition 2.21 (Eigenvalue and eigenvector of an operator) Let $T: X \rightarrow X$ be an operator and $\lambda$ a complex number. If

$$
T x=\lambda x, x \neq 0
$$

then $\lambda$ and $x$ are called eigenvalue and eigenvector of $T$, respectively.
Remark 2.14 It is well known that, in the case of finite dimensional spaces over the field of complex numbers, an operator admits at least one eigenvalue. This is not true, in general, if $X$ is not finite dimensional. For example, the operator $J$ (integral) acting on the space of polynomials $p(s)$ with real coefficients defined as

$$
J: \sum_{i=0}^{n} \alpha_{i} s^{i} \mapsto \sum_{i=0}^{n} \alpha_{i} \frac{s^{i+1}}{i+1}
$$

does not admit eigenvalues. Actually $J p(s)=\lambda p(s)$ holds only for $p(s)=0$. On the contrary, the (derivative) operator $D$ acting again in the space of polynomial with real coefficients, defined as

$$
D: \sum_{i=0}^{n} \alpha_{i} s^{n} \mapsto \sum_{i=0}^{n} i \alpha_{i} s^{i-1}
$$

admits $\lambda=0$ as an eigenvalue, since $D p(s)=0$ for $p(s)=\alpha_{0} \neq 0$. Furthermore, if the rank of the operator is finite, then it admits at least one eigenvalue even if $X$ is not finite dimensional.

If $T: X \rightarrow Y$ is a bounded operator and $X$ and $Y$ are Hilbert spaces (so that suitable inner products are there defined along with the induced norms), it is possible to define the adjoint operator of $T$, hereafter indicated with $T^{*}$. To this aim, consider an operator $T: X \rightarrow Y$ and the map $S:=Y \rightarrow X$ defined by the equation

$$
<T x, y>=<x, S y>, \forall x \in X, \forall y \in Y
$$

The map $S$ is easily shown to be a linear and bounded operator. Hence the following definition is in order.

Definition 2.22 (Adjoint operator) Let $X$ and $Y$ be two normed linear Hilbert spaces with the norm induced by the (relevant) inner products. Let $T: X \rightarrow Y$ be a bounded operator. The adjoint operator $T^{*}$ is defined by

$$
<T x, y>=<x, T^{*} y>, \forall x \in X, \forall y \in Y
$$

The linear operators enjoy a number of useful properties, some of them gathered in the following theorem, whose proof can be easily found in any specialized text.

Theorem 2.9 Let $X$ and $Y$ be two linear Hilbert spaces and let $T: X \rightarrow Y$ and $S: Y \rightarrow X$ two bounded operators. Then,

1) $T^{*}$ is bounded and $\|T\|=\left\|T^{*}\right\|$.
2) $\left\|T^{*} T\right\|=\left\|T T^{*}\right\|=\|T\|^{2}$.
3) $(T+S)^{*}=T^{*}+S^{*}$.
4) If $\alpha \in C,(\alpha T)^{*}=\alpha^{\sim} T^{*}$.
5) $(T S)^{*}=S^{*} T^{*}$.
6) If $T^{-1}$ exists and is bounded, then $\left(T^{-1}\right)^{*}=\left(T^{*}\right)^{-1}$.
7) $\left(T^{*}\right)^{*}=T$.

A particular case occurs when an operator $T: X \rightarrow X$ coincides with its adjoint. In this case $T$ is said to be self-adjoint. For example, the operator $T^{*} T$ is self-adjoint: its eigenvalues, if any, are real and nonnegative.

### 2.8 Functional spaces of rational matrices

The present book often refers to some functional spaces whose elements are proper rational matrices endowed with peculiar properties. To say the true, these elements could be viewed as belonging to suitable subspaces of more general linear spaces whose elements are not necessarily rational functions. However, undertaking this broader perspective is not strictly necessary in the present context, so that the exposition is restricted to the smaller world of rational functions.

Being the elements of the spaces under consideration rational matrices, it should be necessary, for a more rigorous notation, to indicate time by time the dimensions of the relevant matrices. However, this choice is completely useless, since the context widely clarifies the dimensions of the matrices under consideration. Therefore, it will be said that the $n \times m$ matrix $A$ belongs to $X$ instead of $X^{n \times m}$, and so on and so forth. Consistently, the identity (resp. null) matrix will be simply indicated by $I$ (resp. 0).

Definition 2.23 (The space $R L_{\infty}$ ) The set of the rational matrices $F(s)$ such that

$$
\sup _{\omega}\|F(j \omega)\|=\sup _{\omega} \bar{\sigma}[F(j \omega)]<\infty
$$

constitutes the space $R L_{\infty}$.
Definition 2.24 (Norm in $R L_{\infty}$ ) Let $F(s)$ be an element of $R L_{\infty}$. The norm of $F(s)$ is the scalar

$$
\|F(s)\|_{\infty}:=\sup _{\omega} \bar{\sigma}[F(j \omega)]
$$

Remark 2.15 Based on the given definitions it should be obvious that a matrix belongs to $R L_{\infty}$ if any element is a proper rational function without poles on the imaginary axis. In the scalar case, the norm coincides with the peak value of the frequency gain.

An important subspace of $R L_{\infty}$ is that of the rational matrices which are bounded in the right half plane.

Definition 2.25 (The space $R H_{\infty}$ ) The set of rational functions $F(s)$ such that

$$
\sup _{\operatorname{Re}(s) \geq 0}\|F(s)\|<\infty
$$

constitutes the space $R H_{\infty}$.
Remark 2.16 In view of the given definitions, it results that the elements of $R H_{\infty}$ are rational, proper and stable matrices. Moreover, being $R H_{\infty}$ a subspace of $R L_{\infty}$, the norm adopted for the former can be the same as that utilized for the elements of the latter. Further, thanks to a well known property of analytic functions, if $F(s) \in R H_{\infty}$, then

$$
\|F(s)\|_{\infty}=\sup _{\operatorname{Re}(s) \geq 0}\|F(s)\|
$$

Example 2.15 Let consider the functions

$$
\begin{aligned}
& F_{1}(s)=\frac{s+1}{(s+2)(s-3)}, \quad F_{2}(s)=\frac{s-1}{s}, \quad F_{3}(s)=\frac{1}{s-1} \\
& F_{4}(s)=\frac{s}{(s+2)(s+3)}, \quad F_{5}(s)=\frac{s-1}{s+1}, \quad F_{6}(s)=\frac{1}{s+1} \\
& F_{7}(s)=\left[\begin{array}{cc}
(s+1) /\left(s^{2}-4\right) & (s+1) /\left(s^{4}+4\right) \\
1 & (s-1) /(s+2)
\end{array}\right], \quad F_{8}(s)=\left[\begin{array}{l}
s /(s+3) \\
1 /(s-1)
\end{array}\right] \\
& F_{9}(s)=\left[\begin{array}{lll}
1 & (s-1) /(s+1) & 1 /(s+2)
\end{array}\right]
\end{aligned}
$$

Among these functions, $F_{1}(s), F_{3}(s), F_{4}(s), F_{5}(s), F_{6}(s), F_{8}(s)$ and $F_{9}(s)$ belong to $R L_{\infty}$, whereas $F_{4}(s), F_{5}(s), F_{6}(s)$ and $F_{9}(s)$ belong also to $R H_{\infty}$.

In the space $R H_{\infty}$ there are functions which are particularly meaningful. Among them, the so called inner and outer functions are the object of the two definitions below.

Definition 2.26 (Inner function) A function $F(s) \in R H_{\infty}$ is said to be inner if

$$
F^{\sim}(s) F(s)=I, \forall s
$$

Definition 2.27 (Outer function) A function $F(s) \in R H_{\infty}$ is said to be outer if there exists a function $X(s)$, analytic in $\operatorname{Re}(s)>0$, such that

$$
F(s) X(s)=I, \forall s
$$

Remark 2.17 The definition of outer function can be equivalently formulated by requiring that $F(s)$ has full row rank for any $s$ with $\operatorname{Re}(s)>0$. Obviously, if $F(s)$ is a square matrix belonging to $R H_{\infty}$ together with its inverse, then it is outer. Moreover, if $F(s)$ and $G(s)$ belong to $R H_{\infty}$ and are left coprime in the setting of proper stable matrices, then matrix [ $F(s) G(s)]$ is outer.

In the restricted case of scalar functions, it is readily seen that a function $F(s) \in R H_{\infty}$ is inner if $\lim _{s \rightarrow \infty} F(s)=1$ and the zeros and poles are symmetrically positioned, with respect to the origin, in the complex plane. On the other hand, a function $F(s) \in R H_{\infty}$ is outer (minimum phase) if it has no zeros with positive real part.

It is possible to prove that any matrix $F(s) \in R H_{\infty}$ can be decomposed in the product of an inner matrix and an outer matrix (inner outer factorization). Precisely, given $F(s)$ it is always possible to find two matrices $F_{i}(s)$ and $F_{o}(s)$ such that $F(s)=$ $F_{i}(s) F_{o}(s)$. In the scalar case, such a decomposition is simply performed in the way indicated in the proof of the following theorem.

Theorem 2.10 (Inner outer factorization of scalar functions) Let $F(s) \in R H_{\infty}$ be a scalar function. Then, there exist a scalar inner function $F_{i}(s)$ and a scalar outer function $F_{o}(s)$ such that

$$
F(s)=F_{i}(s) F_{o}(s)
$$

Moreover, if $F(j \omega) \neq 0,0 \leq \omega \leq \infty$, then $F_{o}^{-1}(s) \in R H_{\infty}$.
Proof Let $z_{i}, i=1,2, \cdots, \nu$ be the zeros of $F(s)$ with positive real parts (taken with their multiplicity). Moreover, let

$$
\begin{aligned}
& F_{i}(s):=\prod_{i=1}^{\nu} \frac{s-z_{i}}{s+z_{i}^{\sim}} \\
& F_{o}(s):=\frac{F(s)}{F_{i}(s)}
\end{aligned}
$$

Observe that the zeros of $F_{o}(s)$ are the zeros of $F(s)$ with nonpositive real part and the opposite of the zeros of $F(s)$ positive real part. Hence, if $F(s)$ does not have zeros on the extended imaginary axis, then $\operatorname{rhpdeg}\left[F_{o}(s)\right]=0$, so that $F_{o}^{-1}(s) \in R H_{\infty}$.
Example 2.16 Let consider the functions

$$
\begin{aligned}
& F_{1}(s)=1, \quad F_{2}(s)=\frac{s}{s+1}, \quad F_{3}(s)=\frac{s-1}{s+1} \\
& F_{4}(s)=\frac{s-1}{s+2}, \quad F_{5}(s)=\frac{s+1}{s+3}
\end{aligned}
$$

$F_{1}(s)$ and $F_{3}(s)$ are inner, whereas $F_{1}(s), F_{2}(s)$ and $F_{5}(s)$ are outer. The function $F_{4}(s)$ is neither inner nor outer. For this function, take $F_{4 i}(s)=(s-1) /(s+1)$ and $F_{40}=$ $(s+1) /(s+2)$, so that $F_{4}(s)=F_{4 i}(s) F_{4 o}(s)$. Since $F(s)$ has no zeros on the extended imaginary axis, the function $F_{4 o}(s)$ admits as inverse an element of $R H_{\infty}$.

Another meaningful subspace of $R L_{\infty}$ is now defined.
Definition 2.28 (The space $R L_{2}$ ) The set of rational functions $F(s)$ such that

$$
\frac{1}{2 \pi} \int_{-\infty}^{\infty} \operatorname{trace}\left[F^{\sim}(j \omega) F(j \omega)\right] d \omega<\infty
$$

constitutes the space $R L_{2}$.
It is easy to figure out that, taken two rational functions $F(s)$ and $G(s)$, both belonging to $R L_{2}$, the scalar

$$
\alpha:=\frac{1}{2 \pi} \int_{-\infty}^{\infty} \operatorname{trace}\left[G^{\sim}(j \omega) F(j \omega)\right] d \omega<\infty
$$

satisfies all axioms which characterize the inner product, so that the following definition is in order.

Definition 2.29 (Inner product in $R L_{2}$ ) The inner product of two functions of $R L_{2}$ is defined as

$$
<G(s), F(s)>=\frac{1}{2 \pi} \int_{-\infty}^{\infty} \operatorname{trace}\left[G^{\sim}(j \omega) F(j \omega)\right] d \omega
$$

The space $R L_{2}$ is therefore a pre-Hilbert space. To say the true, it is possible to show that $R L_{2}$ is complete, and this implies that it is actually a Hilbert space. Accordingly to what said before, the norm in $R L_{2}$ can be induced by its inner product.

Definition 2.30 (Norm in $R L_{2}$ ) Let $F(s)$ be an element of $R L_{2}$. The norm of $F(s)$ is the scalar

$$
\|F(s)\|_{2}:=\left[\frac{1}{2 \pi} \int_{-\infty}^{\infty} \operatorname{trace}\left[F^{\sim}(j \omega) F(j \omega)\right] d \omega\right]^{1 / 2}
$$

The space $R L_{2}$ can be decomposed into the direct sum of two subspaces, $R H_{2}$ and $R H_{2}^{\perp}$, i.e

$$
R L_{2}=R H_{2} \oplus R H_{2}^{\perp}
$$

where the subspaces $R H_{2}$ and $R H_{2}^{\perp}$ are defined as follows.
Definition 2.31 (The subspaces $R H_{2}$ and $R H_{2}^{\perp}$ ) The subspace $R H_{2}$ is constituted by the functions of $R L_{2}$ which are analytic in the right half plane. Conversely, the subspace $R H_{2}^{\perp}$ is constituted by the functions $R L_{2}$ which are analytic in the left half plane.

Remark 2.18 The given definitions imply that a matrix belongs to $R L_{2}$ if its elements are strictly proper rational functions without poles on the imaginary axis. It belongs to $\mathrm{RH}_{2}$ (resp. $R H_{2}^{\perp}$ ) if its elements are strictly proper rational functions without poles in the closed right (resp. left) half plane. Obviously, it turns out that $R H_{2} \subset R H_{\infty}$.

Remark 2.19 The elements in the spaces $R L_{2}, R H_{2}$ and $R H_{2}^{\perp}$ (which are rational functions of complex variable) can be related to the elements of the spaces $R L_{2}(-\infty \infty)$, $R L_{2}[0 \infty)$, and $R L_{2}(-\infty 0]$, which are functions of the real variable $t$. Such functions are characterized by having rational Fourier transform and being square integrable in the intervals $(-\infty \infty),[0 \infty)$, and $(-\infty 0]$, respectively. Moreover, the elements of $R L_{2}(-\infty 0]$ and $R L_{2}[0 \infty)$ are zero for $t>0$ and $t<0$, respectively, so that both of them belong to $R L_{2}(-\infty \infty)$. Notice that any element of $R L_{2}(-\infty \infty)$ can be uniquely written as the sum of an element of $R L_{2}[0 \infty)$ and one of $R L_{2}(-\infty 0]$. Hence $R L_{2}(-\infty \infty)=$ $R L_{2}[0 \infty) \oplus R L_{2}(-\infty 0]$.

For example, the matrix of functions for which there exists the integral

$$
\int_{0}^{\infty} e^{-j \omega} f(t) d t:=F(j \omega)
$$

and such that

$$
\int_{0}^{\infty} \operatorname{trace}\left[f^{\prime}(t) f(t)\right] d t<\infty
$$

is an element of $R L_{2}[0 \infty)$, which can be put into correspondence with the element $F(s) \in$ $R H_{2}$, provided that this matrix is rational.

Example 2.17 With reference to the functions $F_{i}(s)$ defined in Example 2.15, it follows that $F_{1}(s), F_{3}(s), F_{4}(s)$ and $F_{6}(s)$ are elements of $R L_{2}$. Moreover, $F_{4}(s)$ and $F_{6}(s)$ are also elements of $R H_{2}$ and $F_{3}(s)$ is an element of $R H_{2}^{\perp}$. Finally, notice that

$$
\begin{aligned}
F_{1}(s) & =\frac{s+1}{(s+2)(s-3)}=\frac{1}{5(s+2)}+\frac{4}{5(s-3)} \\
& :=F_{1 s}(s)+F_{1 a}(s)
\end{aligned}
$$

where $F_{1 s}(s) \in R H_{2}$ and $F_{1 a}(s) \in R H_{2}^{\perp}$. Hence, an element of $R L_{2}$ has been shown to equal the sum of a stable function and an antistable one.

Remark 2.20 In linear system theory, the norm of a function in $R H_{2}$ lends itself to particularly significant characterizations. To put them into light, consider the time-invariant system

$$
\begin{aligned}
& \dot{x}=A x+B w, x(0)=0 \\
& z=C x
\end{aligned}
$$

and assume that it is completely reachable and observable. If the matrix $A$ is stable, its transfer function $F(s)$ is an element of $R H_{2}$. Now, consider the $m$-dimensional input vector $w^{(i)}(t)=\delta(t) e_{i}(m)$, where $i=1, \cdots m, e_{i}(m)$ is the $i$-th column of the $m$-dimensional identity matrix and $\delta(t)$ is the "impulsive" function. Let $z^{(i)}$ be the corresponding forced response of the output $z$ and $Z^{(i)}(j \omega)$ its Fourier transform. One wants to evaluate the quantity

$$
J_{1}:=\sum_{i=1}^{m} \int_{0}^{\infty} z^{(i)^{\prime}} z^{(i)} d t
$$

By exploiting the well known Parceval theorem and the properties of the trace operator, it follows

$$
\begin{aligned}
\|F(s)\|_{2}^{2} & =\frac{1}{2 \pi} \int_{-\infty}^{\infty} \operatorname{trace}\left[F^{\prime}(-j \omega) F(j \omega)\right] d \omega \\
& =\frac{1}{2 \pi} \int_{-\infty}^{\infty} \sum_{i=1}^{m} e_{i}(m)^{\prime} F^{\prime}(-j \omega) F(j \omega) e_{2}(m) d \omega \\
& =\frac{1}{2 \pi} \int_{-\infty}^{\infty} \sum_{i=1}^{m} Z^{(i)^{\prime}}(-j \omega) Z^{(i)}(j \omega) d \omega \\
& =\sum_{i=1}^{m} \int_{0}^{\infty} z^{(i) \prime} z^{(i)} d t=J_{1}
\end{aligned}
$$

As for the computation of the scalar $J_{1}$, observe that

$$
\int_{0}^{\infty} z^{(i)^{\prime}} z^{(i)} d t=\int_{0}^{\infty} \operatorname{trace}\left[e_{i}(m)^{\prime} B^{\prime} e^{A^{\prime} t} C^{\prime} C e^{A t} B e_{i}(m)\right] d t
$$

so that

$$
\begin{aligned}
J_{1} & =\sum_{i=1}^{m} \int_{0}^{\infty} \operatorname{trace}\left[e_{i}(m)^{\prime} B^{\prime} e^{A^{\prime} t} C^{\prime} C e^{A t} B e_{i}(m)\right] d t \\
& =\sum_{i=1}^{m} \int_{0}^{\infty} \operatorname{trace}\left[e^{A^{\prime} t} C^{\prime} C e^{A t} B e_{i}(m) e_{i}(m)^{\prime} B^{\prime}\right] d t \\
& =\int_{0}^{\infty} \operatorname{trace}\left[e^{A^{\prime} t} C^{\prime} C e^{A t} \sum_{i=1}^{m} B e_{i}(m) e_{i}(m)^{\prime} B^{\prime}\right] d t
\end{aligned}
$$

$$
\begin{align*}
& =\int_{0}^{\infty} \operatorname{trace}\left[e^{A^{\prime} t} C^{\prime} C e^{A t} B B^{\prime}\right] d t \\
& =\operatorname{trace}\left[B^{\prime} \int_{0}^{\infty} e^{A^{\prime} t} C^{\prime} C e^{A t} d t B\right]  \tag{2.19}\\
& =\operatorname{trace}\left[C \int_{0}^{\infty} e^{A t} B B^{\prime} e^{A^{\prime} t} d t C^{\prime}\right] \tag{2.20}
\end{align*}
$$

As well known, the value of the last two integrals which appear in eqs. (2.19),(2.20) can be computed by resorting to two suitable Lyapunov equations. More precisely, it follows

$$
\begin{align*}
& \operatorname{trace}\left[B^{\prime} \int_{0}^{\infty} e^{A^{\prime} t} C^{\prime} C e^{A t} d t B\right]=\operatorname{trace}\left[B^{\prime} P_{o} B\right] \\
& \operatorname{trace}\left[C \int_{0}^{\infty} e^{A t} B B^{\prime} e^{A^{\prime} t} d t C^{\prime}\right]=\operatorname{trace}\left[C P_{r} C^{\prime}\right] \tag{2.21}
\end{align*}
$$

where $P_{o}$ and $P_{r}$ are the unique solutions of the Lyapunov equations (in the unknown $P$ )

$$
\begin{aligned}
& 0=A^{\prime} P+P A+C^{\prime} C \\
& 0=A P+P A^{\prime}+B B^{\prime}
\end{aligned}
$$

respectively. As a matter of fact, stability of $A$ implies that these equations admit a unique solution, which is also positive semidefinite (recall Lemma C.1).

Consider again the system defined at the beginning of this remark and let now the input $w$ be a white noise with identity intensity. Associated with such a system, consider the quantity

$$
J_{2}:=\lim _{t \rightarrow \infty} \mathrm{E}\left[z^{\prime}(t) z(t)\right]
$$

It follows that

$$
\begin{aligned}
J_{2} & =\lim _{t \rightarrow \infty} \mathrm{E}\left[\operatorname{trace}\left[z^{\prime}(t) z(t)\right]\right] \\
& =\lim _{t \rightarrow \infty} \mathrm{E}\left[\operatorname{trace}\left[z(t) z^{\prime}(t)\right]\right] \\
& =\lim _{t \rightarrow \infty} \operatorname{trace}\left[\mathrm{E}\left[\int_{0}^{t} \Phi(t, \tau) w(\tau) d \tau \int_{0}^{t} w^{\prime}(\sigma) \Phi^{\prime}(t, \sigma) d \sigma\right]\right] \\
& =\lim _{t \rightarrow \infty} \operatorname{trace}\left[\int_{0}^{t} \Phi(t, \tau) \int_{0}^{t} \mathrm{E}\left[w(\tau) w^{\prime}(\sigma)\right] \Phi^{\prime}(t, \sigma) d \sigma d \tau\right] \\
& =\lim _{t \rightarrow \infty} \operatorname{trace}\left[\int_{0}^{t} \Phi(t, \tau) \Phi^{\prime}(t, \tau) d \tau\right]
\end{aligned}
$$

where $\Phi(\xi, \vartheta):=C e^{A(\xi-\vartheta)} B$. Letting $\eta=t-\tau$ one obtains

$$
\begin{aligned}
J_{2} & =\lim _{t \rightarrow \infty} \operatorname{trace}\left[C \int_{0}^{t} e^{A \eta} B B^{\prime} e^{A^{\prime} \eta} d \eta C^{\prime}\right] \\
& =\operatorname{trace}\left[C \int_{0}^{\infty} e^{A \eta} B B^{\prime} e^{A^{\prime} \eta} d \eta C^{\prime}\right]=J_{1}
\end{aligned}
$$

where the last equality follows from what previously shown. Finally, consider again the system fed by a white noise $w$ with identity intensity, and let

$$
J_{3}:=\lim _{T \rightarrow \infty} \frac{1}{T} \mathrm{E}\left[\int_{0}^{T} z^{\prime} z d t\right]
$$

From the analysis performed for the scalar $J_{2}$, one can easily check that

$$
\mathrm{E}\left[z^{\prime}(t) z(t)\right]=\operatorname{trace}\left[C \int_{0}^{t} e^{A^{\prime} \eta} B B^{\prime} e^{A \eta} d \eta C^{\prime}\right]=\operatorname{trace}\left[C P(t) C^{\prime}\right]
$$

with

$$
\begin{equation*}
P(t):=\int_{0}^{t} e^{A \eta} B B^{\prime} e^{A^{\prime} \eta} d \eta \tag{2.22}
\end{equation*}
$$

Therefore,

$$
\begin{aligned}
P(0) & =0 \\
\dot{P}(t) & =e^{A t} B B^{\prime} e^{A^{\prime} t}
\end{aligned}
$$

so that, taking into account eq. (2.22), it follows

$$
\dot{P}(t)=A P(t)+P(t) A^{\prime}+B B^{\prime}
$$

By integrating both sides of this equation from 0 to $T$ and recalling that $P(0)=0$, one gets

$$
P(T)=\int_{0}^{T} \dot{P}(t) d t=A \int_{0}^{T} P(t) d t+\int_{0}^{T} P(t) d t A^{\prime}+B B^{\prime} T
$$

Letting

$$
X(T):=\int_{0}^{T} P(t) d t
$$

it then follows that

$$
\begin{equation*}
\frac{P(T)}{T}=A \frac{X(T)}{T}+\frac{X(T)}{T} A^{\prime}+B B^{\prime} \tag{2.23}
\end{equation*}
$$

Observe that, thanks to eq. (2.22), it is

$$
\lim _{T \rightarrow \infty} P(T)=\lim _{T \rightarrow \infty} \int_{0}^{T} e^{A \eta} B B^{\prime} e^{A^{\prime} \eta} d \eta=P_{r}
$$

so that, taking the limit as $T \rightarrow \infty$ of both members of eq. (2.23), it turns out that

$$
\begin{equation*}
0=A Y+Y A^{\prime}+B B^{\prime} \tag{2.24}
\end{equation*}
$$

where

$$
Y:=\lim _{T \rightarrow \infty} \frac{X(T)}{T}
$$

Since the solution of eq. (2.24) is unique, one can conclude that $Y=P_{r}$ and consequently $J_{3}=J_{1}$.

Remark 2.20 also indicates how the computation of the norm of a function in the space $\mathrm{RH}_{2}$ can be actually performed. As a matter of fact, it is sufficient to solve a Lyapunov equation. Of course, the problem of computing the norm of a function $F(s)$ in $R H_{2}^{\perp}$, is easily solved by observing that $G(s):=F^{\sim}(s) \in R H_{2}$ and that

$$
\begin{aligned}
\|F(s)\|_{2}^{2} & =\frac{1}{2 \pi} \int_{-\infty}^{\infty} \operatorname{trace}\left[F^{\sim}(j \omega) F(j \omega)\right] d \omega \\
& =\frac{1}{2 \pi} \int_{-\infty}^{\infty} \operatorname{trace}\left[F(j \omega) F^{\sim}(j \omega)\right] d \omega \\
& =\frac{1}{2 \pi} \int_{-\infty}^{\infty} \operatorname{trace}\left[G^{\sim}(j \omega) G(j \omega)\right] d \omega \\
& =\|G(s)\|_{2}^{2}
\end{aligned}
$$

Hence, the norm of $F(s) \in R H_{2}$ can be computed by solving a Lyapunov equation associated with $G(s)=F^{\sim}(s) \in R H_{2}$.

Finally, the problem of computing the norm of a function $F(s)$ in the space $R L_{2}$ is easily solvable by writing a decomposition of such a function into the sum of an element of $R \mathrm{H}_{2}$ and another one in $\mathrm{RH}_{2}^{\perp}$. The following result is then provided, whose proof is obvious and then omitted.

Theorem 2.11 Let $F(s)=F_{a}(s)+F_{s}(s)$ with $F_{a}(s) \in R H_{2}^{\perp}$ and $F_{s}(s) \in R H_{2}$. Then,
i) $<F_{a}(s), F_{s}(s)>=0$
ii) $\|F(s)\|_{2}^{2}=\left\|F_{a}(s)\right\|_{2}^{2}+\left\|F_{s}(s)\right\|_{2}^{2}$

Example 2.18 Consider the function

$$
F(s)=\frac{s}{(s+1)(s-1)} \in R L_{2}
$$

Letting

$$
F_{a}(s):=\frac{1}{2(s-1)}, \quad F_{s}(s):=\frac{1}{2(s+1)}
$$

it follows that $F(s)=F_{s}(s)+F_{a}(s)$ and $G_{a}(s):=F_{a}^{\sim}(s)=-1 / 2(s+1)$. With the functions $F_{s}(s)$ and $G_{a}(s)$ let associate the two (minimal) realizations $\Sigma\left(A_{s}, B_{s}, C_{s}, D_{s}\right)$ and $\Sigma\left(A_{a}, B_{a}, C_{a}, D_{a}\right)$, respectively, with $A_{a}=A_{s}=-1, B_{a}=B_{s}=0.5, C_{s}=1, C_{a}=-1$ and $D_{s}=D_{a}=0$. It turns out that

$$
\left\|F_{s}(s)\right\|_{2}^{2}=\left\|G_{a}(s)\right\|_{2}^{2}=\frac{1}{8}
$$

so that

$$
\|F(s)\|_{2}=\frac{1}{2}
$$

Remark 2.21 It is worth noticing that the norm of a generic function $F(s) \in R L_{\infty}$ coincides with that of a suitable function $\hat{F}(s) \in R H_{\infty}$, which is easily derived from $F(s)$. Actually, notice that the least common multiple of all denominators of the elements of $F(s)$, denoted by $\psi(s)$, can be always factorized as $\psi(s):=\psi_{a}(s) \psi_{s}(s)$, where, since $F(s) \in R L_{\infty}$, the polynomial $\psi_{a}(s)$ has all its roots in the open right half plane and $\psi_{s}(s)$ has all its roots in the open left half plane. Hence, if

$$
F(s):=\frac{P(s)}{\psi_{a}(s) \psi_{s}(s)}
$$

being $P(s)$ a polynomial matrix, then

$$
\hat{F}(s):=\frac{P(s)}{\psi_{a}(-s) \psi_{s}(s)}
$$

is such that, $\hat{F}(s) \in R H_{\infty}$ and $F^{\sim}(s) F(s)=\hat{F}^{\sim}(s) \hat{F}(s)$. Consequently, $\|F(s)\|_{\infty}=$ $\|\hat{F}(s)\|_{\infty}$.

The norm of a function $F(s) \in R L_{\infty}$, thought of as the transfer function of a linear system, can be meaningfully related to the norm of the input and output signals, both assumed to be square integrable. Precisely, the following result, besides further characterizing the concept of norm in $R L_{\infty}$ for dynamical systems, also provides an alternative definition of such a norm.

Theorem 2.12 Let $F(s) \in R L_{\infty}$. Then

$$
\|F(s)\|_{\infty}=\sup _{\substack{u \neq 0 \\ u \in R L_{2}}} \frac{\|F(s) u\|_{2}}{\|u\|_{2}}
$$

Proof Observe first that, thanks to Remark 2.21, one can assume, without any loss of generality, that $F(s) \in R H_{\infty}$. As a consequence, it results that $F(s) u \in R L_{2}$ if $u \in R L_{2}$. Now, notice that

$$
\begin{equation*}
\sup _{\substack{u \neq 0 \\ u \in R L_{2}}} \frac{\|F(s) u\|_{2}^{2}}{\|u\|_{2}^{2}} \geq \sup _{\substack{R e(\lambda)>0 \\ u \neq 0}} \frac{\left\|F(s) \frac{\mu}{s-\lambda}\right\|_{2}^{2}}{\left\|\frac{\mu}{s-\lambda}\right\|_{2}^{2}} \tag{2.25}
\end{equation*}
$$

where $\mu$ is a generic constant vector with suitable dimension. Recall that the time response of a system with transfer function $F(s)$ fed by the input $u=\mu e^{\lambda t}$ is, if $\lambda$ does not coincide with any poles of $F(s), y=F(\lambda) \mu e^{\lambda t}$, provided that a suitable initial state is chosen. Then,

$$
F(s) \frac{\mu}{s-\lambda}=F(\lambda) \frac{\mu}{s-\lambda}-f(s)
$$

where $f \in \mathrm{RH}_{2}$ is the Laplace transform of the output free response, whereas the left hand side of the equation corresponds to the output forced response. Since $F(\lambda) \frac{\mu}{s-\lambda} \in$ $R H_{2}^{\perp}$, it follows that (recall Theorem 2.11)

$$
\begin{equation*}
\left\|F(s) \frac{\mu}{s-\lambda}\right\|_{2}^{2}=\left\|F(\lambda) \frac{\mu}{s-\lambda}\right\|_{2}^{2}+\|f(s)\|_{2}^{2} \geq\left\|F(\lambda) \frac{\mu}{s-\lambda}\right\|_{2}^{2} \tag{2.26}
\end{equation*}
$$

From eqs. (2.25), (2.26) it follows that

$$
\begin{align*}
\sup _{\substack{u \neq 0 \\
u \in R L_{2}}} \frac{\|F(s) u\|_{2}^{2}}{\|u\|_{2}^{2}} & \geq \sup _{\substack{R e(\lambda)>0 \\
\mu \neq 0}} \frac{\left\|F(s) \frac{\mu}{s-\lambda}\right\|_{2}^{2}}{\|s-\lambda\|_{2}^{2}} \\
& \geq \sup _{\substack{R e(\lambda)>0 \\
\mu \neq 0}} \frac{\|F(\lambda) \mu\|^{2}}{\|\mu\|^{2}} \\
& \geq \sup _{R e(\lambda)>0}\|F(\lambda)\|^{2} \\
& \geq\|F(s)\|_{\infty} \tag{2.27}
\end{align*}
$$

In getting the above expression, the following two facts have been exploited:

1) If $K$ is a constant, then $\left\|\frac{K}{s-\lambda}\right\|_{2}^{2}=K^{\sim} K\left\|\frac{1}{s-\lambda}\right\|_{2}^{2}=\|K\|^{2}\left\|\frac{1}{s-\lambda}\right\|_{2}^{2}$
2) $\sup _{\mu \neq 0} \frac{\|F(\lambda) \mu\|^{2}}{\|\mu\|^{2}}=\|F(\lambda)\|^{2}$

On the other hand, letting $e_{i}$ denote the $i$-th column of the identity matrix (the context will clarify the relevant dimension), and defining

$$
u:=\left[\begin{array}{lll}
u e_{1} & \cdots & u e_{\nu}
\end{array}\right]
$$

it follows that

$$
\sup _{\substack{u \neq 0 \\ u \in R L_{2}}} \frac{\|F(s) u\|_{2}^{2}}{\|u\|_{2}^{2}}=\sup _{\substack{u \neq 0 \\ u \in R L_{2}}} \frac{\int_{-\infty}^{\infty} \operatorname{trace}\left[u^{\sim}(j \omega) F^{\sim}(j \omega) F(j \omega) u(j \omega)\right] d \omega}{\int_{-\infty}^{\infty} \operatorname{trace}\left[u^{\sim}(j \omega) u(j \omega)\right] d \omega}
$$

$$
\begin{align*}
& =\sup _{\substack{u \neq 0 \\
u \in R L_{2}}} \frac{\int_{-\infty}^{\infty} \sum_{i=1}^{\nu} e_{i}^{\prime} u^{\sim}(j \omega) F^{\sim}(j \omega) F(j \omega) u(j \omega) e_{i} d \omega}{\int_{-\infty}^{\infty} \sum_{i=1}^{\nu} e_{i}^{\prime} u^{\sim}(j \omega) u(j \omega) e_{i} d \omega} \\
& =\sup _{\substack{u \neq 0 \\
u \in R L_{2}}} \frac{\sum_{i=1}^{\nu} \int_{-\infty}^{\infty}\left\|F(j \omega) u(j \omega) e_{i}\right\|^{2} d \omega}{\sum_{i=1}^{\nu} \int_{-\infty}^{\infty}\left\|u(j \omega) e_{i}\right\|^{2} d \omega} \\
& \leq \sup _{\substack{u \neq 0 \\
u \in R L_{2}}} \frac{\sum_{i=1}^{\nu} \int_{-\infty}^{\infty}\|F(j \omega)\|^{2}\left\|u(j \omega) e_{i}\right\|^{2} d \omega}{\sum_{i=1}^{\nu} \int_{-\infty}^{\infty}\left\|u(j \omega) e_{i}\right\|^{2} d \omega} \\
& \leq \sup _{\substack{u \neq 0 \\
u \in R L_{2}}} \frac{\|F(s)\|_{\infty}^{2} \sum_{i=1}^{\nu} \int_{-\infty}^{\infty}\left\|u(j \omega) e_{i}\right\|^{2} d \omega}{\sum_{i=1}^{\nu} \int_{-\infty}^{\infty}\left\|u(j \omega) e_{i}\right\|^{2} d \omega} \\
& \leq\|F(s)\|_{\infty}^{2} \tag{2.28}
\end{align*}
$$

From equations (2.27), (2.28), the result follows.
The computation of the norm of $F(s) \in R L_{\infty}$ can be easily performed by exploiting the result provided in the following theorem, which refers to functions in $R H_{\infty}$, only. As explained in Remark 2.21, this does not entail any loss of generality.
Theorem 2.13 Let $F(s)=C(s I-A)^{-1} B+D$, with $A$ stable and $\bar{\sigma}(D)<\gamma$. Moreover, let $\hat{A}(\gamma):=A+B\left(\gamma^{2} I-D^{\prime} D\right)^{-1} D^{\prime} C$ and

$$
Z(\gamma):=\left[\begin{array}{cc}
\hat{A}(\gamma) & B\left(\gamma^{2} I-D^{\prime} D\right)^{-1} B^{\prime} \\
-C^{\prime}\left(I-\gamma^{-2} D D^{\prime}\right)^{-1} C & -\hat{A}^{\prime}(\gamma)
\end{array}\right]
$$

Then, the following conditions are equivalent:
a) $\|F(s)\|_{\infty}<\gamma$
b) All the eigenvalues of $Z(\gamma)$ do not lie on the imaginary axis
c) The subspace generated by the (generalized) eigenvectors of $Z(\gamma)$, associated with the eigenvalues with negative real parts, is complementary to $\left.\operatorname{Im}\left[\begin{array}{ll}0 & I\end{array}\right]^{\prime}\right]$
d) There exists a symmetric, positive semidefinite and stabilizing solution $S(\gamma)$ of the algebraic Riccati equation (in the unknown $S$ ) associated with $Z(\gamma)$

$$
0=S \hat{A}(\gamma)+\hat{A}^{\prime}(\gamma) S+C^{\prime}\left(I-\gamma^{-2} D D^{\prime}\right)^{-1} C+S B\left(\gamma^{2} I-D^{\prime} D\right)^{-1} B^{\prime} S
$$

namely, such that matrix $\hat{A}(\gamma)+B\left(\gamma^{2} I-D^{\prime} D\right)^{-1} B^{\prime} S(\gamma)$ is stable.
Proof Preliminarily observe that the assumption on the norm of $D$ ensures, thanks to Lemma B. 11 that both matrices $\left(I-\gamma^{-2} D D^{\prime}\right)$ and $\left(\gamma^{2} I-D^{\prime} D\right)$ are positive definite so that the statement of the theorem is well defined.
$a) \Longrightarrow b$ ) Assume, by contradiction, that $Z(\gamma)$ has an eigenvalue in $j \omega$ and let $\xi:=\left[x^{\prime} y^{\prime}\right]^{\prime} \neq 0$ be an associated eigenvector. From $Z(\gamma) \xi=j \omega \xi$ it follows

$$
\begin{aligned}
(j \omega-\hat{A}(\gamma)) x & =B\left(\gamma^{2} I-D^{\prime} D\right)^{-1} B^{\prime} y \\
-\left(j \omega+\hat{A}^{\prime}(\gamma)\right) y & =C^{\prime}\left(I-\gamma^{-2} D D^{\prime}\right)^{-1} C x
\end{aligned}
$$

By recalling the definition of $\hat{A}(\gamma)$ and the identity $\left(I-\gamma^{-2} D D^{\prime}\right)^{-1}=I+D\left(\gamma^{2} I-\right.$ $\left.D^{\prime} D\right)^{-1} D^{\prime}$ (Lemma B.9), these expressions can be rewritten as follows

$$
\begin{align*}
(j \omega-A) x & =B\left(\gamma^{2} I-D^{\prime} D\right)^{-1}\left(B^{\prime} y+D^{\prime} C x\right)  \tag{2.29}\\
-\left(j \omega+A^{\prime}\right) y & =C^{\prime} C x+C^{\prime} D\left(\gamma^{2} I-D^{\prime} D\right)^{-1}\left(D^{\prime} C x+B^{\prime} y\right) \tag{2.30}
\end{align*}
$$

Let now

$$
\begin{align*}
& v:=\left(\gamma^{2} I-D^{\prime} D\right)^{-1}\left(B^{\prime} y+D^{\prime} C x\right)  \tag{2.31}\\
& z:=C x+D v \tag{2.32}
\end{align*}
$$

In view of eqs. (2.29),(2.30) and thanks to the stability of $A$, these last expressions become

$$
\begin{align*}
& x=(j \omega I-A)^{-1} B v  \tag{2.33}\\
& y=-\left(j \omega I+A^{\prime}\right)^{-1} C^{\prime} z \tag{2.34}
\end{align*}
$$

which, substituted into eqs. (2.31),(2.32) yield

$$
z=F(j \omega) v, \gamma^{2} v=F^{\prime}(-j \omega) z
$$

that is

$$
F^{\prime}(-j \omega) F(j \omega) v=\gamma^{2} v
$$

Being $v \neq 0$ (otherwise, both $x$ and $z$ would be zero and thus also $y$ ), it results

$$
\begin{aligned}
\gamma^{2} & =\frac{v^{\sim} F^{\sim}(j \omega) F(j \omega) v}{v^{\sim} v} \\
& =\frac{\|F(j \omega) v\|^{2}}{\|v\|^{2}} \\
& \leq \sup _{v \neq 0} \frac{\|F(j \omega) v\|^{2}}{\|v\|^{2}} \\
& \leq\|F(j \omega)\|^{2} \\
& \leq \sup _{\omega}\|F(j \omega)\|^{2}=\|F(s)\|_{\infty}^{2}
\end{aligned}
$$

which is a contradiction.
$b) \Longrightarrow a)$ Conversely, suppose that $\|F(s)\|_{\infty} \geq \gamma$. Since, by assumption, $\|D\|=$ $\lim _{\omega \rightarrow \infty}\|F(\omega)\|<\gamma$, by a continuity argument there exists a real number $\omega$ such that $\|F(j \omega)\|=\gamma$. Therefore, $\gamma^{2}$ is an eigenvalue (actually the maximum one) of $F^{\sim}(j \omega) F(j \omega)$ so that there exists a nonzero vector $v$ such that $F^{\sim}(j \omega) F(j \omega) v=\gamma^{2} v$. Being $v \neq 0$, it also follows that $z:=F(j \omega) v \neq 0$. Now define two vectors $x$ and $y$ as in eqs. (2.33),(2.34). With such definitions, simple computations show that also eqs. (2.29)-(2.32) hold. From eq. (2.31) one can conclude that $\xi:=\left[x^{\prime} y^{\prime}\right]^{\prime} \neq 0$ since, on the contrary, $v$ would be zero. Finally, eqs. (2.29),(2.30) imply that $Z(\gamma) \xi=j \omega \xi$, which is a contradiction.
$d) \Longrightarrow b$ ) The existence of a stabilizing solution of the Riccati equation introduced in the statement implies, by a well known result, that the eigenvalues of $Z(\gamma)$ do not lie on the imaginary axis.
$b) \Longrightarrow d)$ Notice first that, as obvious, stability of $A$ guarantees stabilizability of the pair $(A, B)$ and, recalling the definition of $\hat{A}(\gamma)$, also that of the pair $(\hat{A}(\gamma), B)$. This condition implies also the stabilizability of the pair $\left(\hat{A}(\gamma), B\left(\gamma^{2} I-D^{\prime} D\right)^{-1} B^{\prime}\right)$. To see this, suppose by contradiction that, for a certain $x \neq 0$ and $\lambda, \operatorname{Re}(\lambda)>0$, it is $\hat{A}^{\prime}(\gamma) x=\lambda x$ and $B\left(\gamma^{2} I-D^{\prime} D\right)^{-1} B^{\prime} x=0$ (recall the $P B H$ test). It follows that $x^{\sim} B\left(\gamma^{2} I-D^{\prime} D\right)^{-1} B^{\prime} x=0$ so that $B^{\prime} x=0$ thanks to the fact that $\left(\gamma^{2} I-D^{\prime} D\right)^{-1}>0$. The conditions $x \neq 0, \operatorname{Re}(\lambda)>0, \hat{A}^{\prime}(\gamma) x=\lambda x$ and $B^{\prime} x=0$, finally violate the stabilizability of $(\hat{A}(\gamma), B)$. Lemma C. 3 can now be applied to the Riccati equation,
ensuring the existence of its symmetric and stabilizing solution $S(\gamma)$. It is only left to show that such a solution is actually positive semidefinite. To this aim, rewrite the Riccati equation as follows

$$
\begin{aligned}
0 & =S(\gamma) A+A^{\prime} S(\gamma)+C^{\prime} C+\left[S(\gamma) B+C^{\prime} D\right]\left(\gamma^{2} I-D^{\prime} D\right)^{-1}\left[S(\gamma) B+C^{\prime} D\right]^{\prime} \\
& :=S(\gamma) A+A^{\prime} S(\gamma)+W(\gamma)
\end{aligned}
$$

where $W(\gamma) \geq 0$. Thanks to Lemma C.1, it can be concluded that $S(\gamma) \geq 0$.
$c) \Longrightarrow b$ ) Condition $c$ ) implies that the subspace generated by the (generalized) eigenvectors of $Z(\gamma)$ associated with the eigenvalues with negative real part has the same dimension $n$ as the system. Hence $Z(\gamma)$ has $n$ eigenvalues with negative real part and, obviously, has no eigenvalues on the imaginary axis.
$b) \Longrightarrow c)$ Since condition $b$ ) has already been proved to be equivalent to condition d), it follows that the stabilizing solution $S(\gamma)$ is such that $\operatorname{Im}\left[[I S(\gamma)]^{\prime}\right]$, which is obviously complementary to $\operatorname{Im}\left[[0 I]^{\prime}\right]$, is actually the subspace generated by $n$ (generalized) linear independent eigenvectors of $Z(\gamma)$ associated with the $n$ eigenvalues with negative real part.

Theorem 2.14 Let $F(s):=C(s I-A)^{-1} B+D$ and $\gamma$ a positive scalar. Then, the following two conditions are equivalent:
a) The matrix $A$ is stable and $\|F(s)\|_{\infty}<\gamma$
b) $\bar{\sigma}(D)<\gamma$ and there exists the positive semidefinite stabilizing solution $S(\gamma)$ of the algebraic Riccati equation (in the unknown $S$ )

$$
0=S \hat{A}(\gamma)+\hat{A}^{\prime}(\gamma) S+C^{\prime}\left(I-\gamma^{-2} D D^{\prime}\right)^{-1} C+S B\left(\gamma^{2} I-D^{\prime} D\right)^{-1} B^{\prime} S
$$

i.e. such that

$$
\hat{A}(\gamma)+B\left(\gamma^{2} I-D^{\prime} D\right)^{-1} B^{\prime} S(\gamma)=A+B\left(\gamma^{2} I-D^{\prime} D\right)^{-1}\left(B^{\prime} S(\gamma)+D^{\prime} C\right)
$$

is stable.
Proof $a) \Longrightarrow b$ ) If $\|F(s)\|_{\infty}<\gamma$, then, obviously $\bar{\sigma}(D)<\gamma$ since $D=\lim _{\omega \rightarrow \infty} F(j \omega)$. Therefore, Theorem 2.13 guarantees the existence of $S(\gamma)$.
$b) \Longrightarrow a$ ) Preliminarily observe that $\bar{\sigma}(D)<\gamma$ implies, thanks to Lemma B.11, that matrices $\left(I-\gamma^{-2} D D^{\prime}\right)$ and $\left(\gamma^{2} I-D^{\prime} D\right)$ are positive definite so that the Riccati equation is well defined. Such an equation can be equivalently rewritten as (recall the definition for $\hat{A}(\gamma)$ and Lemma B.9)

$$
0=S A+A^{\prime} S+\left(S B+C^{\prime} D\right)\left(\gamma^{2} I-D^{\prime} D\right)^{-1}\left(B^{\prime} S+D^{\prime} C\right)+C^{\prime} C
$$

The stabilizing solution $S(\gamma)$ is also a solution of the Lyapunov equation

$$
0=S A+A^{\prime} S+\Lambda^{\prime}(\gamma) \Lambda(\gamma)
$$

where

$$
\Lambda^{\prime}(\gamma) \Lambda(\gamma):=\left[S(\gamma) B+C^{\prime} D\right]\left(\gamma^{2} I-D^{\prime} D\right)^{-1}\left[B^{\prime} S(\gamma)+D^{\prime} C\right]+C^{\prime} C
$$

The proof is concluded by showing that the pair $(A, \Lambda(\gamma))$ is detectable. Actually, in this case, Lemma C. 1 implies that $A$ is stable, so that the proof follows directly from

Theorem 2.13. Therefore, assume by contradiction that this pair is not detectable. Based on the PBH test (recall Lemma D.2) it follows that

$$
\begin{aligned}
& A x=\lambda x, \operatorname{Re}(\lambda) \geq 0, x \neq 0 \\
& \Lambda(\gamma) x=0
\end{aligned}
$$

Recalling the definition of $\Lambda^{\prime}(\gamma) \Lambda(\gamma)$, the second equation implies that

$$
x^{\sim} \Lambda^{\prime}(\gamma) \Lambda(\gamma) x=x^{\sim}\left[S(\gamma) B+C^{\prime} D\right]\left(\gamma^{2} I-D^{\prime} D\right)^{-1}\left[B^{\prime} S(\gamma)+D^{\prime} C\right] x+x^{\sim} C^{\prime} C x=0
$$

so that, in particular,

$$
\left[B^{\prime} S(\gamma)+D^{\prime} C\right] x=0
$$

thanks to the fact that $C^{\prime} C \geq 0$ and $\left(\gamma^{2} I-D^{\prime} D\right)>0$. Hence,

$$
\left\{A+B\left(\gamma^{2} I-D^{\prime} D\right)^{-1}\left[B^{\prime} S(\gamma)+D^{\prime} C\right]\right\} x=A x=\lambda x
$$

This is obviously a contradiction, since $x \neq 0, \operatorname{Re}(\lambda) \geq 0$ and $S(\gamma)$ is stabilizing.
Example 2.19 Consider the rational function

$$
F(s)=\frac{10}{s^{2}+2 s+10}=\left[\begin{array}{rr|r}
0 & 1 & 0 \\
-10 & -2 & 1 \\
\hline 10 & 0 & 0
\end{array}\right]
$$

It turns out that $|F(j \omega)|^{2}=100 /\left(\omega^{4}-16 \omega^{2}+100\right)$ so that $\|F(s)\|_{\infty}=5 / 3 \simeq 1.66$. Taken $\gamma=1.67$, the eigenvalues of $Z(\gamma)$ are $\pm 0.067 \pm 2.83 j$ whereas for $\gamma=1.66$ the eigenvalues of $Z(\gamma)$ are $\pm 2.92 j$ and $\pm 2.73 j$.

Remark 2.22 Theorems 2.13 and 2.14 call for a Riccati equation which resembles the one utilized in the context of optimal LQ control. In view of Lemma 2.16 and Remark 2.13, it results that

$$
\|F(j \omega)\|=\bar{\sigma}[F(j \omega)]=\bar{\sigma}\left[F^{\prime}(-j \omega)\right]=\left\|F^{\prime}(j \omega)\right\|
$$

so that

$$
\|F(s)\|_{\infty}=\left\|F^{\prime}(s)\right\|_{\infty}
$$

Hence, the two relevant theorems can be equivalently reformulated with reference to $F^{\prime}(s)=$ $B^{\prime}\left(s I-A^{\prime}\right)^{-1} C^{\prime}+D^{\prime}$ instead of $F(s)$. As a consequence, the relevant Riccati equation becomes the following

$$
0=S \hat{A}^{\prime}(\gamma)+\hat{A}(\gamma) S+B\left(I-\gamma^{-2} D^{\prime} D\right)^{-1} B^{\prime}+S C^{\prime}\left(\gamma^{2} I-D D^{\prime}\right)^{-1} C S
$$

which resembles the equation involved in the optimal filtering problem.
The previous results on the characterization of an upper bound of $\|F(s)\|_{\infty}$ for $F(s)$ being a function in $R H_{\infty}$ are all expressed in terms of the existence of a symmetric, positive semidefinite and stabilizing solution of an algebraic Riccati equation (recall Theorems 2.13 and 2.14) which are of particular importance to get the results of Chapter 5 . However, in Chapter 6 reference is made to the following similar results which are alternatively expressed in terms of Riccati inequalities with no additional requirement concerning the stability of their solutions.

Theorem 2.15 Let $F(s):=C(s I-A)^{-1} B$ and $\gamma$ a positive scalar. Then, the following conditions are equivalent:
a) Matrix $A$ is stable and $\|F(s)\|_{\infty}<\gamma$
b) There exists a symmetric and positive definite matrix $S$ satisfying the Riccati inequality

$$
\begin{equation*}
S A+A^{\prime} S+\gamma^{-2} S B B^{\prime} S+C^{\prime} C<0 \tag{2.35}
\end{equation*}
$$

c) There exists a symmetric and positive definite matrix $P$ satisfying the Riccati inequality

$$
\begin{equation*}
P A^{\prime}+A P+\gamma^{-2} P C^{\prime} C P+B B^{\prime}<0 \tag{2.36}
\end{equation*}
$$

Proof That points $b$ ) and $c$ ) are equivalent is straightforward. Indeed, the positive definite solutions (if any) of both inequalities are related one to the other by $P=$ $\gamma^{2} S^{-1}$.
$b) \Longrightarrow a)$ Since there exists $S>0$ satisfying (2.35), it also satisfies the Riccati equation

$$
\begin{equation*}
0=S A+A^{\prime} S+\gamma^{-2} S B B^{\prime} S+\bar{C}^{\prime} \bar{C} \tag{2.37}
\end{equation*}
$$

where $\bar{C}^{\prime}=\left[\begin{array}{ll}C^{\prime} E^{\prime}\end{array}\right]$ for some matrix $E$ such that $E^{\prime} E>0$. From the Extended Lyapunov lemma the stability of matrix $A$ follows. With $S>0$ being a solution of (2.37), define the auxiliary Riccati equation

$$
\begin{equation*}
0=X\left(A+\gamma^{-2} B B^{\prime} S\right)^{\prime}+\left(A+\gamma^{-2} B B^{\prime} S\right)^{\prime} X-\gamma^{-2} X B B^{\prime} X+E^{\prime} E \tag{2.38}
\end{equation*}
$$

Using the fact that $A$ is stable and $E^{\prime} E>0$, from Lemma C. 3 it is readily verified that it admits a symmetric and positive semidefinite solution $X \geq 0$ such that matrix

$$
\begin{align*}
A_{X} & :=A+\gamma^{-2} B B^{\prime} S-\gamma^{-2} B B^{\prime} X \\
& =A+\gamma^{-2} B B^{\prime}(S-X) \tag{2.39}
\end{align*}
$$

is stable. Defining $W:=S-X$ and using (2.37) together with (2.38) we get

$$
0=W A+A^{\prime} W+\gamma^{-2} W B B^{\prime} W+C^{\prime} C
$$

Moreover, since $A$ is stable, this means that $W \geq 0$ and from (2.39) matrix $A+$ $\gamma^{-2} B B^{\prime} W$ is stable. From Theorem 2.14 the conclusion is that part $a$ ) holds indeed.
$a) \Longrightarrow b$ ) Let us define the transfer function

$$
\bar{F}(s):=\left[\begin{array}{c}
C \\
\sqrt{\epsilon} I
\end{array}\right](s I-A)^{-1} B
$$

where $\epsilon>0$ is a scalar to be determined. It is a simple matter to verify that for all $\omega \in R$

$$
\bar{F}^{\prime}(-j \omega) \bar{F}(j \omega)=F^{\prime}(-j \omega) F(j \omega)+\epsilon G^{\prime}(-j \omega) G(j \omega)
$$

where $G(s):=(s I-A)^{-1} B$. Hence

$$
\|\bar{F}(s)\|_{\infty}^{2} \leq\|F(s)\|_{\infty}^{2}+\epsilon\|G(s)\|_{\infty}^{2}
$$

and choosing the scalar $\epsilon$ such that

$$
0<\epsilon<\frac{\gamma^{2}-\|F(s)\|_{\infty}^{2}}{\|G(s)\|_{\infty}^{2}}
$$

which is always possible since by assumption, matrix $A$ is stable and $\|F(s)\|_{\infty}<\gamma$ then we get $\|\bar{F}(s)\|_{\infty}<\gamma$. Using again Theorem 2.14 the conclusion is that there exists a symmetric and positive semidefinite solution $S \geq 0$ to the Riccati equation

$$
0=S A+A^{\prime} S+\gamma^{-2} S B B^{\prime} S+C^{\prime} C+\epsilon I
$$

which obviously satisfies the Riccati inequality (2.35) and in view of Lemma C. 1 is actually positive definite. The proof of the theorem proposed is complete.

Theorem 2.16 Let $F(s):=C(s I-A)^{-1} B$ and $\gamma$ a positive scalar. Under the assumption that the pair $(A, B)$ is reachable, the following conditions are equivalent:
a) Matrix $A$ is stable and $\|F(s)\|_{\infty} \leq \gamma$
b) There exists a symmetric and positive definite matrix $S$ satisfying the Riccati inequality

$$
\begin{equation*}
S A+A^{\prime} S+\gamma^{-2} S B B^{\prime} S+C^{\prime} C \leq 0 \tag{2.40}
\end{equation*}
$$

c) There exists a symmetric and positive definite matrix $P$ satisfying the Riccati inequality

$$
\begin{equation*}
P A^{\prime}+A P+\gamma^{-2} P C^{\prime} C P+B B^{\prime} \leq 0 \tag{2.41}
\end{equation*}
$$

Proof The equivalence of points $b$ ) and $c$ ) is immediate. The positive definite matrices satisfying inequalities (2.40) and (2.41), if any, are related one to the other by $S=$ $\gamma^{2} P^{-1}$.
$b) \Longrightarrow a)$ Assuming (2.40) admits a positive definite feasible solution then $P=$ $\gamma^{2} S^{-1}$ is feasible for inequality (2.41) which together with the reachability of the pair ( $A, B$ ) implies that $A$ is stable. On the other hand, for all $\omega \in R$, inequality (2.40) can be rewritten as

$$
\left(-j \omega-A^{\prime}\right) S+S(j \omega-A)-\gamma^{-2} S B B^{\prime} S-C^{\prime} C \geq 0
$$

from which and $G(s):=\gamma I-\gamma^{-1} B^{\prime} S(s I-A)^{-1} B$ we have

$$
\begin{aligned}
F^{\prime}(-j \omega) F(j \omega) & \leq \gamma^{2} I-G^{\prime}(-j \omega) G(j \omega) \\
& \leq \gamma^{2} I, \quad \forall \omega \in R
\end{aligned}
$$

consequently $\|F(s)\|_{\infty} \leq \gamma$ which is the desired result.
$a) \Longrightarrow b$ ) From Theorem 2.15 we only need to prove that if $\|F(s)\|_{\infty}=\gamma$ then the Riccati inequality (2.40) is still feasible for some positive definite matrix. To this end, consider the sequence of matrices $C_{n}:=\sqrt{\epsilon_{n}} C$ with $\epsilon_{n}$ being an arbitrary element of an increasing sequence of scalars such that $0<\epsilon_{n}<1$ and $\epsilon_{n}$ goes to 1 as $n$ goes to infinite. By virtue of

$$
\begin{align*}
\left\|C_{n}(s I-A)^{-1} B\right\|_{\infty} & =\sqrt{\epsilon_{n}}\|F(s)\|_{\infty} \\
& =\sqrt{\epsilon_{n}} \gamma<\gamma \tag{2.42}
\end{align*}
$$

we already know that the Riccati equation (in the unknown $S$ )

$$
\begin{equation*}
0=S A+A^{\prime} S+\gamma^{-2} S B B^{\prime} S+C_{n}^{\prime} C_{n} \tag{2.43}
\end{equation*}
$$

admits an unique symmetric, positive semidefinite solution $S_{n}$ such that matrix $A_{n}:=A+\gamma^{-2} B B^{\prime} S_{n}$ is stable. Additionally, the sequence $S_{n} \geq 0, n=1,2, \cdots$ is nondecreasing because as it can be verified

$$
\begin{aligned}
0= & \left(S_{n+1}-S_{n}\right) A_{n}+A_{n}^{\prime}\left(S_{n+1}-S_{n}\right)+ \\
& +\gamma^{-2}\left(S_{n+1}-S_{n}\right) B B^{\prime}\left(S_{n+1}-S_{n}\right)+\left(\epsilon_{n+1}-\epsilon_{n}\right) C^{\prime} C
\end{aligned}
$$

Once again, due to (2.42), the Riccati equation (in the unknown $P$ )

$$
\begin{equation*}
0=P A^{\prime}+A P+\gamma^{-2} P C_{n}^{\prime} C_{n} P+B B^{\prime} \tag{2.44}
\end{equation*}
$$

also admits an unique symmetric, positive semidefinite and stabilizing solution $P_{n}$ which in fact is positive definite since $A$ stable and $(A, B)$ reachable yield

$$
P_{n} \geq \int_{0}^{\infty} e^{A t} B B^{\prime} e^{A^{\prime} t} d t>0
$$

Using the same reasoning adopted before we can show that the sequence $P_{n}>0$, $n=1,2, \cdots$ is nondecreasing as well. Hence the sequence of positive matrices $\bar{S}_{n}:=$ $\gamma^{2} P_{n}^{-1}$ is nonincreasing and satisfy the Riccati equation (2.43). Moreover, defining matrix $\bar{A}_{n}:=A+\gamma^{-2} B B^{\prime} \bar{S}_{n}$ we have from (2.44) with $P=P_{n}$

$$
\begin{aligned}
P_{n} \bar{A}_{n}^{\prime} P_{n}^{-1} & =P_{n}\left(A+B B^{\prime} P_{n}^{-1}\right)^{\prime} P_{n}^{-1} \\
& =-\left(A+\gamma^{-2} P_{n} C_{n}^{\prime} C_{n}\right)
\end{aligned}
$$

implying that $-\bar{A}_{n}$ is stable. Finally taking into account that $S_{n}$ and $\bar{S}_{n}$ solve the Riccati equation (2.43) we get

$$
0=\left(\bar{S}_{n}-S_{n}\right) A_{n}+A_{n}^{\prime}\left(\bar{S}_{n}-S_{n}\right)+\gamma^{-2}\left(\bar{S}_{n}-S_{n}\right) B B^{\prime}\left(\bar{S}_{n}-S_{n}\right)
$$

which shows from the stability of matrix $A_{n}$ that $S_{n} \leq \bar{S}_{n}$. This last inequality together with the fact that the sequence $S_{n} \geq 0, n=1,2, \cdots$ is nondecreasing and the sequence $\bar{S}_{n}>0, n=1,2, \cdots$ is nonincreasing allow the conclusion that both sequences converge to some matrices such that

$$
0 \leq S_{\infty}:=\lim _{n \rightarrow \infty} S_{n} \leq \lim _{n \rightarrow \infty} \bar{S}_{n}:=\bar{S}_{\infty}
$$

At this point, it remains to prove that $\bar{S}_{\infty}$ is positive definite even though some eigenvalues of matrix $\bar{A}_{\infty}$ lie on the imaginary axis. Indeed, assume by contradiction that $\bar{S}_{\infty} \geq 0$. Since it solves the Riccati equation

$$
0=\bar{S}_{\infty} A+A^{\prime} \bar{S}_{\infty}+\gamma^{-2} \bar{S}_{\infty} B B^{\prime} \bar{S}_{\infty}+C^{\prime} C
$$

from Lemma C. 1 and the $P B H$ test there exits $x \neq 0$ such that $A x=\lambda x, B^{\prime} \bar{S}_{\infty} x=0$ and $C x=0$, that is $\bar{A}_{\infty} x=\lambda x$. However, being $A$ stable, this is impossible since, as proved before, all eigenvalues of matrix $\bar{A}_{\infty}$ are located on the right part (including the imaginary axis) of the complex plane.

Example 2.20 Consider the rational function $F(s)$ defined in Example 2.19. We notice that for this transfer function the pair $(A, B)$ is reachable and the pair $(A, C)$ is observable. Taking $\gamma=5 / 3$, matrices $S_{\infty}$ and $\bar{S}_{\infty}$ (recall the proof of Theorem 2.16) are found to be

$$
S_{\infty}=\bar{S}_{\infty}=\left[\begin{array}{rr}
55.56 & 5.56 \\
5.56 & 5.56
\end{array}\right]
$$

It is interesting to see that, in this case, the eigenvalues of matrices $A_{\infty}=A+\gamma^{-2} B B^{\prime} S_{\infty}=$ $\bar{A}_{\infty}$ are $\pm 2.83 j$. That is, the corresponding Riccati equation solution $\bar{S}_{\infty}$ is no more stabilizing but still, for $\gamma=5 / 3=\|F(s)\|_{\infty}$, there exists a symmetric and positive definite matrix satisfying the Riccati inequality (2.40). One of such matrices is exactly $\bar{S}_{\infty}$.

Example 2.21 Consider the rational function $F(s)$

$$
F(s)=\frac{1}{s+2}=\frac{s+1}{s^{2}+3 s+2}=\left[\begin{array}{rr|r}
0 & 1 & 0 \\
-2 & -3 & 1 \\
\hline 1 & 1 & 0
\end{array}\right]
$$

so that $\|F(s)\|_{\infty}=1 / 2$. It is apparent that the pair $(A, B)$ is reachable but the pair $(A, C)$ is not observable. With $\gamma=0.5$ the following matrices appearing in the proof of Theorem 2.16 have been calculated

$$
S_{\infty}=\left[\begin{array}{ll}
0.50 & 0.50 \\
0.50 & 0.50
\end{array}\right] \leq\left[\begin{array}{ll}
0.50 & 0.50 \\
0.50 & 1.00
\end{array}\right]=\bar{S}_{\infty}
$$

In this case, matrix $S_{\infty}$ is positive semidefinite and the eigenvalues of $A_{\infty}$ are 0 and -1 . On the other hand, matrix $\bar{S}_{\infty}$ is positive definite and the eigenvalues of $\bar{A}_{\infty}$ are 0 and 1. By construction, matrix $\bar{S}_{\infty}$ satisfies the Riccati inequality (2.40).

An important role in the development of the discussion in Chapter 5 will be played by a few linear operators acting on the spaces previously defined.

Definition 2.32 (The Laurent operator with symbol $\mathbf{F}$ ) Let $F(s) \in R L_{\infty}$. The $\operatorname{map} \Lambda_{F}: R L_{2} \rightarrow R L_{2}$ defined as

$$
\Lambda_{F}: G(s) \mapsto \Lambda_{F} G(s):=F(s) G(s)
$$

is called the Laurent operator with symbol $F(s)$.
Notice that the operator $\Lambda_{F}$ is obviously linear and, thanks to Theorem 2.12, $\left\|\Lambda_{F}\right\|=$ $\|F(s)\|_{\infty}$ so that $\Lambda_{F}$ is bounded.

Definition 2.33 (Orthogonal stable and antistable projections) Let $G(s)$ be a generic element of $R L_{2}$ such that $G(s):=G_{s}(s)+G_{a}(s)$ with $G_{s}(s) \in R H_{2}$ and $G_{a}(s) \in R H_{2}^{\perp}$. The map $\Pi_{s}: R L_{2} \rightarrow R H_{2}$, defined by

$$
\Pi_{s}: G(s) \mapsto \Pi_{s} G(s):=G_{s}(s)
$$

is called the stable (orthogonal) projection, whereas the map $\Pi_{a}: R L_{2} \rightarrow R H_{2}^{\perp}$, defined by

$$
\Pi_{a}: G(s) \mapsto \Pi_{a} G(s):=G_{a}(s)
$$

is called the antistable (orthogonal) projection.
It is immediate to check that $\Pi_{s}$ and $\Pi_{a}$ are actually linear operators.
Definition 2.34 (The Hankel operator with symbol F) Let $F(s) \in R L_{\infty}$. The $\operatorname{map} \Gamma_{F}: R H_{2}^{\perp} \rightarrow R H_{2}$ defined by

$$
\Gamma_{F}: G(s) \mapsto \Gamma_{F} G(s):=\Pi_{s} \Lambda_{F} G(s)
$$

is called the Hankel operator with symbol $F$.
Remark 2.23 It should be evident from the above definition that the Hankel operator is the result of the composition of the Laurent operator with the stable projection, i.e. $\Gamma_{F}=\Pi_{s} \Lambda_{F}$. Hence, $\Gamma_{F}$ is indeed an operator.

Remark 2.24 Sometimes the Hankel operator is defined as mapping $R H_{2}$ to $R H_{2}^{\perp}$ in such a way that $\Gamma_{F}=\Pi_{a} \Lambda_{F}$. This should not be surprising. Actually, given the existing isomorphism between $\mathrm{RH}_{2}$ and $\mathrm{RH}_{2}^{\perp}$, the deep understanding of the basic facts of the underlying theory is certainly not blurred by the presence of these different (although equivalent) definitions.

Remark 2.25 Thanks to Definition 2.34, it turns out that only the strictly proper and stable part of $F(s)$ contributes to $\Gamma_{F} G(s)$. As a matter of fact, letting $F(s)=F_{\infty}+F_{s}(s)+$ $F_{a}(s)$, where $F_{\infty}=\lim _{s \rightarrow \infty} F(s), F_{s}(s) \in R H_{2}, F_{a}(s) \in R H_{2}^{\perp}$, it follows

$$
\begin{aligned}
\Gamma_{F} G(s) & =\Pi_{s} \Lambda_{F} G(s) \\
& =\Pi_{s} F(s) G(s) \\
& =\Pi_{s}\left[F_{\infty}+F_{s}(s)+F_{a}(s)\right] G(s) \\
& =\Gamma_{F_{s}} G(s)
\end{aligned}
$$

Hence, $\Gamma_{F}=\Gamma_{F_{s}}$.
Example 2.22 A way to compute the components $F_{\infty}, F_{a}(s)$ and $F_{s}(s)$ of $F(s) \in R L_{\infty}$, calls for a (minimal) realization $\Sigma(A, B, C, D)$ of it. Let $n_{s}$ (resp. $n_{a}$ ) be the number of eigenvalues of $A$ with negative (resp. positive) real parts and define $\mathcal{X}_{s}$ (resp. $\mathcal{X}_{a}$ ) as the $n_{s}$ dimensional ( $n_{a}$ dimensional) subspace generated by $n_{s}$ (resp. $n_{a}$ ) linear independent generalized eigenvectors of $A$ associated with the eigenvalues with negative (resp. positive) real part. Such subspaces are obviously complementary and can be identified by the image of suitable full rank matrices $X_{s}$ and $X_{a}$, namely $\mathcal{X}_{s}=\operatorname{Im}\left[X_{s}\right]$ and $\mathcal{X}_{a}=\operatorname{Im}\left[X_{a}\right]$. Then, matrix $\left[\begin{array}{ll}X_{s} & X_{a}\end{array}\right]$ is invertible and, letting $T:=\left[\begin{array}{ll}X_{s} & X_{a}\end{array}\right]^{-1}$, it results

$$
\hat{A}:=T A T^{-1}:=\left[\begin{array}{cc}
A_{s} & 0 \\
0 & A_{a}
\end{array}\right]
$$

where $A_{s}$ is stable and $A_{a}$ antistable. Now, let

$$
\hat{B}:=T B:=\left[\begin{array}{l}
B_{s} \\
B_{a}
\end{array}\right], \hat{C}:=C T^{-1}:=\left[\begin{array}{ll}
C_{s} & C_{a}
\end{array}\right]
$$

so that

$$
\begin{aligned}
F(s) & =\hat{C}(s I-\hat{A})^{-1} \hat{B}+D \\
& =\left[\begin{array}{ll}
C_{s} & C_{a}
\end{array}\right]\left[\begin{array}{cc}
s I-A_{s} & 0 \\
0 & s I-A_{a}
\end{array}\right]^{-1}\left[\begin{array}{l}
B_{s} \\
B_{a}
\end{array}\right]+D \\
& =C_{s}\left(s I-A_{s}\right)^{-1} B_{s}+C_{a}\left(s I-A_{a}\right)^{-1} B_{a}+D \\
& :=F_{s}(s)+F_{a}(s)+F_{\infty}
\end{aligned}
$$

Specifically, consider the function

$$
F(s)=\Sigma(A, B, C, D)=\frac{s^{2}+s+1}{s^{2}-1} \in R L_{\infty}
$$

with

$$
A=\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right], \quad B=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad C=\left[\begin{array}{ll}
2 & 1
\end{array}\right], \quad D=1
$$

It follows

$$
X_{s}=\left[\begin{array}{c}
1 \\
-1
\end{array}\right], \quad X_{a}=\left[\begin{array}{l}
1 \\
1
\end{array}\right]
$$

and $A_{s}=-1, A_{a}=1, B_{s}=-1 / 2, B_{a}=1 / 2, C_{s}=1, C_{a}=3$ so that $F_{\infty}=1, F_{s}(s)=$ $-1 / 2(s+1)$ and $F_{a}(s)=3 / 2(s-1)$. Let now

$$
G(s):=\sum_{i=1}^{k} \frac{\alpha_{i}}{\left(s-p_{i}\right)^{n_{i}}}, \operatorname{Re}\left(p_{i}\right)>0
$$

It turns out that $F(s) G(s)=G(s)-G(s) / 2(s+1)+3 G(s) / 2(s-1)$ so that

$$
\Gamma_{F} G(s)=\Pi_{s}[F(s) G(s)]=-\frac{G(-1)}{2(s+1)}
$$

On the other hand

$$
F_{s}(s) G(s)=-\frac{G(-1)}{2(s+1)}+\sum_{i=1}^{k} \frac{\beta_{i}}{\left(s-p_{i}\right)^{n_{i}}}, \operatorname{Re}\left(p_{i}\right)>0
$$

where $\beta_{i}$ are suitable scalars. Hence

$$
\Gamma_{F_{s}} G(s)=\Pi_{s}\left[F_{s}(s) G(s)\right]=-\frac{G(-1)}{2(s+1)}
$$

The adjoint Laurent and Hankel operators are now defined in the two lemmas below according to Definition 2.22 .

Lemma 2.26 (The adjoint Laurent operator with symbol $\mathbf{F}$ ) The adjoint Laurent operator with symbol $F$ is the Laurent operator with symbol $F^{\sim}$, i.e.

$$
\Lambda_{F}^{*}=\Lambda_{F} \sim
$$

Proof Let $G_{1}(s)$ and $G_{2}(s)$ be two generic elements of $R L_{2}$ and $F(s) \in R L_{\infty}$. Then, by recalling Definition 2.22,

$$
\begin{aligned}
<G_{1}(s), \Lambda_{F}^{*} G_{2}(s)> & =<\Lambda_{F} G_{1}(s), G_{2}(s)> \\
& =\frac{1}{2 \pi} \int_{-\infty}^{\infty} \operatorname{trace}\left[G_{1}^{\sim}(j \omega) F^{\sim}(j \omega) G_{2}(j \omega) d \omega\right. \\
& =<G_{1}(s), F^{\sim}(s) G_{2}(s)> \\
& =<G_{1}(s), \Lambda_{F \sim} \sim G_{2}(s)>
\end{aligned}
$$

Being $G_{1}(s)$ arbitrary, from this expression it follows

$$
\Lambda_{F}^{*} G_{2}(s)=\Lambda_{F^{\sim}} G_{2}(s)
$$

so that $\Lambda_{F}^{*}=\Lambda_{F^{\sim}}$, since $G_{2}(s)$ is arbitrary too.
Lemma 2.27 (The adjoint Hankel operator with symbol F) The adjoint Hankel operator with symbol $F$ is the operator $\Gamma_{F}^{*}:=\Pi_{a} \Lambda_{F}^{*}=\Pi_{a} \Lambda_{F^{\sim}}$

Proof Let $G_{1}(s) \in R H_{2}^{\perp}, H(s) \in R H_{2}$ and $F(s) \in R L_{\infty}$. Preliminarily observe that

$$
\begin{aligned}
<G(s), \Gamma_{F}^{*} H(s)> & =<\Gamma_{F} G(s), H(s)> \\
& =<\Pi_{s} \Lambda_{F} G(s), H(s)> \\
& =<\Pi_{s} F(s) G(s), H(s)>=<F(s) G(s), H(s)>
\end{aligned}
$$

since $F(s) G(s)=\Pi_{s} F(s) G(s)+\Pi_{a} F(s) G(s)$ and being $H(s) \in R H_{2}$, it turns out that $<\Pi_{a} F(s) G(s), H(s)>=0$ (recall Theorem 2.11). Analogously, $F^{\sim}(s) H(s)=$ $\Pi_{a} F^{\sim}(s) H(s)+\Pi_{s} F^{\sim}(s) H(s)$, so that Theorem 2.11 implies $<G(s), \Pi_{s} F^{\sim}(s) H(s)>$ $=0$, since $G(s) \in R H_{2}^{\perp}$. Taking in mind this fact, it follows

$$
<F(s) G(s), H(s)>=<G(s), F^{\sim}(s) H(s)>=<G(s), \Pi_{a} F^{\sim}(s) H(s)>
$$

so that

$$
<G(s), \Gamma_{F}^{*} H(s)>=<G(s), \Pi_{a} F^{\sim}(s) H(s)>=<G(s), \Pi_{a} \Lambda_{F} \sim H(s)>
$$

Since $G(s)$ is arbitrary it follows that $\Gamma_{F}^{*} H(s)=\Pi_{a} \Lambda_{F^{\sim}} H(s)$. Finally, being $H(s)$ arbitrary, it results $\Gamma_{F}^{*}=\Pi_{a} \Lambda_{F} \sim=\Pi_{a} \Lambda_{F}^{*}$.

Interestingly, one can associate with a generic Hankel operator, defined in the frequency domain, a suitable function of time. Precisely, consider the operator $\Gamma_{F}$ and assume, without any loss of generality, (recall Remark 2.25), that

$$
F(s):=\left[\begin{array}{c|c}
\mathrm{A} & \mathrm{~B} \\
\hline \mathrm{C} & 0
\end{array}\right]
$$

with $A$ stable and $\Sigma(A, B, C, 0)$ in minimal form with order $n$. Then, the function

$$
f(t):= \begin{cases}0 & t \leq 0 \\ C e^{A t} B & t>0\end{cases}
$$

is the inverse Laplace transform of $F(s)$. The map

$$
\Gamma_{f}: R L_{2}(-\infty 0] \rightarrow R L_{2}[0 \infty)
$$

defined by

$$
\Gamma_{f}: u(t) \mapsto \Gamma_{f} u(t):=y(t):= \begin{cases}0 & t<0 \\ C e^{A t} \int_{-\infty}^{0} e^{-A \tau} B u(\tau) d \tau & t \geq 0\end{cases}
$$

is easily shown to be an operator, which is legitimated to be considered as the time domain counterpart of the operator $\Gamma_{F}$. Actually, the Laplace transform of $y$ equals the stable projection of the product of $F(s)$ with the Laplace transform of $u$, i.e.

$$
\begin{equation*}
y_{L}=\left(\Gamma_{f} u\right)_{L}=\Gamma_{F} u_{L} \tag{2.45}
\end{equation*}
$$

It is also useful to characterize further the operator $\Gamma_{f}$ as follows. Consider the system $\Sigma(A, B, C, 0)$ and assume that its initial state at $t=-\infty$ is zero, i.e. $x(-\infty)=0$. Now apply to the system an input $u(\cdot)$ which is different from zero only for nonpositive time instants. Hence,

$$
x(0)=\int_{-\infty}^{0} e^{-A \tau} B u d \tau
$$

Moreover, consider the free output $y$, defined for nonnegative time instants, due to the above defined initial state at $t=0$. It turns out that the operator $\Gamma_{f}$ maps the input $u(t \leq 0)$ to the output $y(t \geq 0)$, through $x(0)$. As such, $\Gamma_{f}$ can be viewed as the composition of two operators $\Psi_{r}: R L_{2}(-\infty 0] \rightarrow R^{n}$ and $\Psi_{o}: R^{n} \rightarrow R L_{2}[0 \infty)$, defined by

$$
\begin{aligned}
& \Psi_{r}: u(t) \mapsto \Psi_{r} u(t):=x:=\int_{-\infty}^{0} e^{-A \tau} B u d \tau \\
& \Psi_{o}: x \mapsto \Psi_{o} x:=y(t):= \begin{cases}0 & t<0 \\
C e^{A t} x & t \geq 0\end{cases}
\end{aligned}
$$

Actually, it results $\Gamma_{f}=\Psi_{o} \Psi_{r}$. The observability operator $\Psi_{o}$ is obviously injective (the unique element of $R^{n}$ which is mapped in zero is the zero element), thanks to the assumed observability condition. Slightly less obvious is to recognize that the operator $\Psi_{r}$ is surjective. In fact, thanks to the reachability condition, any element of $R^{n}$ is the result of the transformation of a suitable element of $R L_{2}(\infty 0]$. To see this, let $x \in R^{n}$ be a fixed element. The input

$$
u(t):= \begin{cases}0 & t>0 \\ B^{\prime} e^{-A^{\prime} t}\left[\int_{-\infty}^{0} e^{-A \sigma} B B^{\prime} e^{-A^{\prime} \sigma} d \sigma\right]^{-1} x, & t \leq 0\end{cases}
$$

is such that $\Psi_{r} u=x$. Actually, recalling the discussion in Remark 2.24, it follows

$$
\int_{-\infty}^{0} e^{-A \sigma} B B^{\prime} e^{-A^{\prime} \sigma} d \sigma=\int_{0}^{\infty} e^{A \sigma} B B^{\prime} e^{A^{\prime} \sigma} d \sigma=P_{r}
$$

where $P_{r}$ is the unique solution of the Lyapunov equation (in the unknown $P$ ) $P A^{\prime}+$ $A P+B B^{\prime}=0$. Moreover, such solution is positive definite (for this consider points $i$ ) and $i i$ ) of the "dual" version of Lemma C.1, i.e. when the pair $(A, C)$ is replaced by the pair $\left.\left(A^{\prime}, B^{\prime}\right)\right)$. Hence $u$ is well defined. Moreover, since $A$ is stable, $u \in R L_{2}(-\infty 0]$. Finally,

$$
\Psi_{r} u=\int_{-\infty}^{0} e^{-A \tau} B B^{\prime} e^{-A^{\prime} \tau} d \tau\left[\int_{-\infty}^{0} e^{-A \sigma} B B^{\prime} e^{-A^{\prime} \sigma} d \sigma\right]^{-1} x=x
$$

By resorting to Definition 2.22 it is easy to determine the adjoint operators of $\Psi_{r}$ and $\Psi_{o}$. It results that $\Psi_{r}^{*}: R^{n} \rightarrow R L_{2}(-\infty 0]$ and $\Psi_{o}^{*}: R L_{2}[0 \infty) \rightarrow R^{n}$ are defined by

$$
\begin{aligned}
& \Psi_{r}^{*}: x \mapsto \Psi_{r}^{*} x:=u(t):= \begin{cases}B^{\prime} e^{-A^{\prime} t} x & t \leq 0 \\
0 & t>0\end{cases} \\
& \Psi_{o}^{*}: y(t) \mapsto \Psi_{o}^{*} y(t):=x:=\int_{0}^{\infty} e^{A^{\prime} \sigma} C^{\prime} y d \sigma
\end{aligned}
$$

respectively. By exploiting the fact that $\Gamma_{f}=\Psi_{o} \Psi_{r}$ and taking into account the above expressions, it follows, in view of Theorem 2.9, that $\Gamma_{f}^{*}: R L_{2}[0 \infty) \rightarrow R L_{2}(-\infty 0]$ is defined by

$$
\Gamma_{f}^{*}: y(t) \mapsto \Gamma_{f}^{*} y(t):=u(t):= \begin{cases}B^{\prime} e^{A^{\prime} t} \int_{0}^{\infty} e^{A^{\prime} \tau} C^{\prime} y d \tau & t \leq 0 \\ 0 & t>0\end{cases}
$$

This operator is legitimated to be the time domain counterpart of the operator $\Gamma_{F}^{*}$. Actually, the Laplace transform of $u$ equals the antistable projection of the product of $F^{\sim}(s)$ with the Laplace transform of $y$, i.e.

$$
\begin{equation*}
u_{L}=\left(\Gamma_{f}^{*} y\right)_{L}=\Gamma_{F}^{*} y_{L} \tag{2.46}
\end{equation*}
$$

The operators $\Psi_{r}$ and $\Psi_{o}$ have rank equal to $n$, i.e. equal to the dimension of the minimal realization $\Sigma(A, B, C, 0)$. Moreover, since $\Psi_{r}$ is surjective and $\Psi_{o}$ is injective, it turns out that $\Gamma_{f}=\Psi_{o} \Psi_{r}$ and, obviously, $\Gamma_{F}$ have rank equal to $n$. Therefore the self adjoint operator $\Gamma_{F}^{*} \Gamma_{F}$ has rank $n$ and, thanks to Remark 2.14 it admits eigenvalues.

Remark 2.26 The operator $\Gamma_{F}^{*} \Gamma_{F}$ has a zero eigenvalue. Actually, write $F(s) \in R H_{2}$ as $F(s)=N(s) / \psi(s)$, where $N(s)$ is a polynomial matrix and $\psi(s)$ is the least common multiple of all denominators of the elements of $F(s)$. Moreover, let $H(s) \in R H_{2}^{\perp}$ be such that

$$
G(s):=\psi(s) H(s) \in R H_{2}^{\perp}
$$

Then,

$$
\begin{aligned}
\Gamma_{F} G(s) & =\Pi_{s} F(s) G(s) \\
& =\Pi_{s} \frac{N(s)}{\psi(s)} \psi(s) H(s) \\
& =\Pi_{s} N(s) H(s)=0
\end{aligned}
$$

since $N(s) H(s) \in R H_{2}^{\perp}$. Consequently, $\Gamma_{F}^{*} \Gamma_{F} G(s)=0=0 G(s)$ so that $\lambda=0$ is an eigenvalue.

Example 2.23 Let $F(s)=1 /(s+1)$. Taken $G(s)=1 /(s-1)$, it turns out that

$$
\Gamma_{F} G(s)=\Pi_{s} F(s) G(s)=\frac{G(-1)}{s+1}=-\frac{1}{2(s+1)}
$$

On the other hand,

$$
\Gamma_{F}^{*} \Gamma_{F} G(s)=\Gamma_{F}^{*}\left[-\frac{1}{2(s+1)}\right]=\Pi_{a} F^{\sim}(s)\left[-\frac{1}{2(s+1)}\right]=\frac{1}{4} G(s)
$$

so that $\lambda=1 / 4$ is an eigenvalue and $G(s)$ an eigenvector of $\Gamma_{F}^{*} \Gamma_{F}$.
The operator $\Gamma_{F}^{*} \Gamma_{F}$ enjoys the interesting properties stated in the following result, for the proof of which the reader is referred to specialized texts.

Theorem 2.17 Let $F(s) \in R L_{\infty}$. The eigenvalues of $\Gamma_{F}^{*} \Gamma_{F}$ are real and nonnegative. The greatest of them is $\left\|\Gamma_{F}\right\|^{2}$.

Example 2.24 Consider the function $F(s)=1 /(s+1)$ already introduced in Example 2.23. It is easy to check that the equation $\Gamma_{F}^{*} \Gamma_{F} G(s)=\lambda G(s)$, with $G(s) \in R H_{2}^{\perp}$ and $G(s) \neq 0$, is satisfied only for $\lambda=0$ or $\lambda=0.25$. Hence, $\Gamma_{F}^{*} \Gamma_{F}$ has two distinct eigenvalues, which are real and nonnegative. The norm of $\Gamma_{F}^{*} \Gamma_{F}$ is equal to 0.25 and $\left\|\Gamma_{F}\right\|=0.5$ (recall Theorem 2.9, point 2)).

The operators $\Psi_{r} \Psi_{r}^{*}$ and $\Psi_{o}^{*} \Psi_{o}$ map the space $R^{n}$ into itself and, as such, can be represented by suitable (real) $n$-dimensional matrices. Actually, as for the former, it follows

$$
\begin{aligned}
x_{r} & :=\Psi_{r} \Psi_{r}^{*} x_{0} \\
& =\int_{-\infty}^{0} e^{-A \tau} B B^{\prime} e^{-A^{\prime} \tau} d \tau x_{0} \\
& =\int_{0}^{\infty} e^{A \tau} B B^{\prime} e^{A^{\prime} \tau} d \tau x_{0}=P_{r} x_{0}
\end{aligned}
$$

where $P_{r}$ again represents the unique solution of the Lyapunov equation in the unknown $P$ (recall Lemma C.1)

$$
\begin{equation*}
0=A P+P A^{\prime}+B B^{\prime} \tag{2.47}
\end{equation*}
$$

Such a solution is also positive definite. Since the equation $\Psi_{r} \Psi_{r}^{*} x_{0}=P_{r} x_{0}$ holds for any $x_{0}$, it results

$$
\begin{equation*}
\Psi_{r} \Psi_{r}^{*}=P_{r} \tag{2.48}
\end{equation*}
$$

As for the operator $\Psi_{o}^{*} \Psi_{o}$ it follows that, for any $x_{0}$,

$$
x_{f}:=\Psi_{o}^{*} \Psi_{o} x_{0}=\int_{o}^{\infty} e^{A^{\prime} \tau} C^{\prime} C e^{A \tau} d \tau x_{0}=P_{o} x_{0}
$$

where, thanks to what has been shown in Remark 2.20, matrix $P_{o}$ is the unique solution of the Lyapunov equation in the unknown $P$

$$
\begin{equation*}
0=P A+A^{\prime} P+C^{\prime} C \tag{2.49}
\end{equation*}
$$

Such a solution is positive definite (recall points $i$ ) and $i i i$ ) of Lemma C.1). Hence, since $x_{0}$ is arbitrary,

$$
\begin{equation*}
\Psi_{o}^{*} \Psi_{o}=P_{o} \tag{2.50}
\end{equation*}
$$

The link previously expressed between $\Gamma_{F}$ and $\Gamma_{f}$ allows one to precisely state the following important result, which also provides a procedure for the computation of the eigenvalues of $\Gamma_{F}^{*} \Gamma_{F}$ and hence the norm of $\Gamma_{F}$.

Theorem 2.18 The operator $\Gamma_{F}^{*} \Gamma_{F}$ and the matrix $P_{r} P_{o}$ share the same nonzero eigenvalues.

Proof Let $\lambda \neq 0$ be an eigenvalue of $\Gamma_{F}^{*} \Gamma_{F}$. In view of eqs. (2.45),(2.46) it is easy to check that, if $0 \neq u_{L} \in R H_{2}^{\perp}$ is an associate eigenvector, then its inverse Laplace transform, $u \neq 0$, is an eigenvector of $\Gamma_{f}^{*} \Gamma_{f}$, associated with the same eigenvalue. Moreover, in view of Theorem 2.9, it is $\Gamma_{f}^{*}=\Psi_{r}^{*} \Psi_{o}^{*}$ so that

$$
\begin{equation*}
\Psi_{r}^{*} \Psi_{o}^{*} \Psi_{o} \Psi_{r} u=\lambda u \tag{2.51}
\end{equation*}
$$

Premultiplying both sides of this equation by $\Psi_{r}$ and defining $x_{0}:=\Psi_{r} u$, one obtains (recall eqs. (2.48),(2.50))

$$
\begin{equation*}
P_{r} P_{o} x_{0}=\lambda x_{0} \tag{2.52}
\end{equation*}
$$

Notice that $x_{0} \neq 0$. Actually, if not, then $\Psi_{r} u=0$ and, from eq. (2.51), $\lambda u=0$. Hence $\lambda \neq 0$ would imply $u=0$, a contradiction.

Conversely, if $\lambda \neq 0$ is an eigenvalue of $P_{r} P_{o}$, then there exists $x_{0} \neq 0$ such that eq. (2.52) holds. Premultiplying this expression by $\Psi_{r}^{*} P_{o}$, defining $u(t):=\Psi_{r}^{*} P_{o} x_{0}$, and recalling eqs. (2.48),(2.50), eq. (2.51) is derived. It is left to show that $u \neq 0$. If it were not so, eqs. $(2.52),(2.48)$ would imply $\lambda x_{0}=0$, a contradiction. Hence $u$ is an eigenvector of $\Gamma_{f}^{*} \Gamma_{f}$ associated with $\lambda$ and its Laplace transform is an eigenvector of $\Gamma_{F}^{*} \Gamma_{F}$, associated with the same eigenvalue.

Remark 2.27 (Computation of the norm of the Hankel operator) As already said, Theorem 2.18 provides a procedure for the computation of $\left\|\Gamma_{F}\right\|$. It can be summarized in the following steps, with reference to the case in which $F(s) \in R H_{2}$. Such situation can be always be matched by using the arguments discussed in Remark 2.20.

1) Determine a minimal realization $\Sigma(A, B, C, 0)$ of $F(s)$.
2) Solve the two Lyapunov equations (2.47),(2.49).
3) Compute the greatest eigenvalue $\lambda_{M}^{2}$ of $P_{r} P_{o}$
4) $\left\|\Gamma_{F}\right\|=\lambda_{M}$

Example 2.25 Take again the function $F(s)=1 /(s+1)$ considered in Examples 2.23, 2.24. Being $F(s)=\Sigma(-1,1,1,0)$, eqs. (2.47),(2.49) have solutions $P_{r}=P_{o}=0.5$, so that $P_{r} P_{o}=0.25$ and $\left\|\Gamma_{F}\right\|=0.5$.

At the end of the present section, it is presented an important result, the so called Nehari theorem, which sets at the basis of an operatorial technique for the solution of the control problem in the $R H_{\infty}$ context, treated in Chapter 5 . The proof of this result is not simple and therefore is not reported here.

Theorem 2.19 (Nehari theorem) Let $F(s) \in R L_{\infty}$. There exists $X^{o}(s) \in R H_{\infty}$ such that

$$
\inf _{X(s) \in R H_{\infty}}\|F(s)-X(s)\|_{\infty}=\left\|F(s)-X^{o}(s)\right\|_{\infty}=\left\|\Gamma_{F \sim}\right\|
$$

The Nehari theorem concerns the problem of finding a stable function $X(s)$ which approximate an assigned function of $R L_{\infty}$ by minimizing the distance between $X(s)$ and $F(s)$. Such a distance is defined as the $R L_{\infty}$ norm of the difference $F(s)-$ $X(s)$. The theorem comes out with two important conclusions: first, the existence of a function $X^{\circ}(s) \in R H_{\infty}$ which represents the best approximation of $F(s)$ and, second, that the minimal distance is given by the norm of the Hankel operator with symbol $F^{\sim}$. It should be apparent, consistently with such an interpretation, that if $F(s)$ is stable, the minimal distance is zero, since one can choose $X^{o}(s)=F(s)$. In general, the optimal approximation is not unique. A way to compute an optimal approximation will be presented in Chapter 5 , in the case where $F(s)$ is a scalar.

Example 2.26 Consider the function $F(s)=\gamma /(s-\beta)$ with $\beta>0$ and $\gamma>0$. Then, $F^{\sim}(s)=-\gamma /(s+\beta)$. Based on Remark 2.27 one finds that $\left\|\Gamma_{F \sim}\right\|=|\gamma / 2 \beta|$.

### 2.9 Notes and references

More about the material of this chapter can be found in many places. Restricting the attention to system and control theory point of view (which is the one mostly pursued in the present book), it is worth quoting the following sources with reference to the various sections. Sections 2.3 and 2.4: the books of Kailath [29], Maciejowski [43] and Vidyasagar [60]. Section 2.5: besides the books [29], [60], the papers by MacFarlane and Karcanias [42] and Kouvaritakis and MacFarlane [33], [34]. Section 2.6: the book of Lawson and Hanson [38]. Section 2.7: any text of functional analysis, for instance that of Rudin [55]. Section 2.8: the book of Francis [19], the paper by Boyd et al. [9] on the computation of the $H_{\infty}$ norm and the book of Power [52] for the Nehari theorem.

## Chapter 3

## Feedback Systems Stability

### 3.1 Introduction

One of the most significant problems of linear control theory is no doubt that of characterizing the set of all controllers which stabilize a given system. The present chapter is devoted to this problem, which is also conceptually linked with the results presented in the forthcoming Chapters 4 and 5 .

In order to avoid any possible misunderstanding, it is well advisable to place in the right context the problem covered in the present treatment. Therefore, consider the system shown in fig. 3.1 where $\Sigma_{1}$ is a given time invariant and finite dimensional linear system and $\Sigma_{2}$ is a controller which receives informations from the system through the output variable $y$ only and can drive the system through the control variable $u$. Its duty is rendering the feedback connected system in fig. 3.1 (asymptotically) stable. An obvious necessary condition for a stabilizing regulator to exist is stabilizability and detectability of system $\Sigma_{1}$. This condition is also sufficient if $\Sigma_{2}$ can be chosen in the class $\mathcal{C}_{L}$ of linear time invariant finite dimensional systems. This is precisely the class to which the present discussion will be limited.

The actual determination of an element of $\mathcal{C}_{L S}\left(\Sigma_{1}\right)$, the subset of $\mathcal{C}_{L}$ constituted by the stabilizing controllers for $\Sigma_{1}$, can be worked out by resorting to a number of classical techniques. Among them, it is worth recalling the optimal linear quadratic Gaussian control or the pole assignment technique. Greatly more difficult is rather the problem of individuating the whole set $\mathcal{C}_{L S}\left(\Sigma_{1}\right)$, in the sense of establishing a significant correspondence between an element of the set and a suitable free parameter.

In Section 3.4 it will be shown how the set $\mathcal{C}_{L S}\left(\Sigma_{1}\right)$ can be parametrized as a


Figure 3.1: The feedback connection of two systems


Figure 3.2: The set $\mathcal{C}_{L S}\left(\Sigma_{1}\right)$


Figure 3.3: The system $\Sigma$ - feedback connection of $G_{1}(s)$ and $G_{2}(s)$
function of a free parameter $Q(s)$. Specifically, it will be proved that: $i$ ) each element of $\mathcal{C}_{L S}\left(\Sigma_{1}\right)$ can be represented accordingly to the block scheme in fig. 3.2 , where $R_{f}$ is a system which is determined from the given system $\Sigma_{1}$ once for ever, whereas $Q(s)$ is a linear system belonging to a suitable subset of $\left.R H_{\infty} ; i i\right)$ chosen an element, one obtains an element of $\mathcal{C}_{L S}\left(\Sigma_{1}\right)$. In order to present the relevant results, some preliminary facts are now introduced. In particular, Section 3.2 is devoted to discuss the basic relationships between internal and $B I B O$ (external) stability of a feedback system, whereas in Section 3.3 the important concept of double coprime factorization is introduced.

### 3.2 Internal and external stability

Consider system $\Sigma$ depicted in fig. 3.3 where each one of the two blocks with transfer functions $G_{1}(s)$ and $G_{2}(s)$ is a detectable and stabilizable system (otherwise $\Sigma$ should never be stable). The aim of this section is presenting some simple results which establish precise connections between the (internal) stability of $\Sigma$ and the (external) stability of some transfer functions defined in the block-scheme of fig. 3.3. Notice that the (stable) dynamic matrix of the unreachable and/or unobservable parts of the two subsystems which constitute $\Sigma$ do not affect anyone of the transfer functions which can be defined in the scheme of fig. 3.3. For this reason reference will be made, without loss of generality, to minimal realizations of both $G_{1}(s)$ and $G_{2}(s)$, i.e.

$$
G_{i}(s):=C_{i}\left(s I-A_{i}\right)^{-1} B_{i}+D_{i}, \quad i=1,2
$$

with

$$
\begin{aligned}
\dot{x}_{i} & =A_{i} x_{i}+B_{i} u_{i}, \quad i=1,2 \\
y_{i} & =C_{i} x_{i}+D_{i} u_{i}, \quad i=1,2 \\
u_{1} & =v_{1}+y_{2}
\end{aligned}
$$

$$
u_{2}=v_{2}+y_{1}
$$

where, for $i=1,2$, the pair $\left(A_{i}, B_{i}\right)$ is reachable and the pair $\left(A_{i}, C_{i}\right)$ is observable. Obviously, the above equations make sense only provided that the matrix $\Delta_{21}:=$ $\left(I-D_{2} D_{1}\right)$ (or, equivalently, matrix $\left.\Delta_{12}:=\left(I-D_{1} D_{2}\right)\right)$ is nonsingular. Letting $x:=\left[x_{1}^{\prime} x_{2}^{\prime}\right]^{\prime} \in R^{n}, u:=\left[\begin{array}{ll}u_{1}^{\prime} & u_{2}^{\prime}\end{array}\right]^{\prime}, v:=\left[\begin{array}{ll}v_{1}^{\prime} & v_{2}^{\prime}\end{array}\right]^{\prime}$, a realization of $\Sigma$ is given by

$$
\begin{aligned}
\dot{x} & =A x+B v \\
u & =C x+D v
\end{aligned}
$$

with

$$
\begin{aligned}
& A:=\left[\begin{array}{cc}
A_{1}+B_{1} \Delta_{21}^{-1} D_{2} C_{1} & B_{1} \Delta_{21}^{-1} C_{2} \\
B_{2} \Delta_{12}^{-1} C_{1} & A_{2}+B_{2} \Delta_{12}^{-1} D_{1} C_{2}
\end{array}\right] \\
& B:=\left[\begin{array}{cc}
B_{1} \Delta_{21}^{-1} & B_{1} \Delta_{21}^{-1} D_{2} \\
B_{2} \Delta_{12}^{-1} D_{1} & B_{2} \Delta_{12}^{-1}
\end{array}\right] \\
& C:=\left[\begin{array}{cc}
\Delta_{21}^{-1} D_{2} C_{1} & \Delta_{21}^{-1} C_{2} \\
\Delta_{12}^{-1} C_{1} & \Delta_{12}^{-1} D_{1} C_{2}
\end{array}\right] \\
& D:=\left[\begin{array}{cc}
\Delta_{21}^{-1} & \Delta_{21}^{-1} D_{2} \\
\Delta_{12}^{-1} D_{1} & \Delta_{12}^{-1}
\end{array}\right]
\end{aligned}
$$

having exploited Lemma B.9. Letting $T(u, v ; s):=C(s I-A)^{-1} B+D$ denote the transfer function from the input $v$ to the output $u$, the following result can be proved.

Theorem 3.1 The system $\Sigma$ in fig. 3.3 is internally stable if and only if $T(u, v ; s) \in$ $R H_{\infty}$.

Proof If matrix $A$ is stable, then, obviously, $T(u, v ; s) \in R H_{\infty}$. In order to prove the converse statement it will be shown that all the eigenvalues of $A$ coincide with the poles of $T(u, v ; s)$, or, in other words that the four matrices $(A, B, C, D)$ constitute a minimal realization of $T(u, v ; s)$. The proof is worked out by exploiting the $P B H$ test (see Lemmas D. 1 and D.3).

Recall that the pair $(A, C)$ is observable if and only if

$$
\operatorname{rank}\left[P_{C}(\lambda)\right]:=\operatorname{rank}\left[\left[\begin{array}{c}
\lambda I-A \\
C
\end{array}\right]\right]=n, \quad \forall \lambda
$$

Moreover, the matrix

$$
T_{o}:=\left[\begin{array}{cccc}
I & 0 & B_{1} & 0 \\
0 & I & 0 & B_{2} \\
0 & 0 & I & -\Delta_{21}^{-1} D_{2} \Delta_{12} \\
0 & 0 & -\Delta_{12}^{-1} D_{1} \Delta_{21} & I
\end{array}\right]
$$

is nonsingular since the Schur formula for the computation of the determinant of a block matrix (see Lemma B.13) gives
$\operatorname{det}\left[\left[\begin{array}{cc}I & -\Delta_{21}^{-1} D_{2} \Delta_{12} \\ -\Delta_{12}^{-1} D_{1} \Delta_{21} & I\end{array}\right]\right]=\operatorname{det}\left[I-\Delta_{12}^{-1} D_{1} \Delta_{21} \Delta_{21}^{-1} D_{2} \Delta_{12}\right]=\operatorname{det}\left[\Delta_{12}\right]$

Thus,

$$
\begin{aligned}
\operatorname{rank}\left[P_{C}(\lambda)\right] & =\operatorname{rank}\left[T_{o} P_{C}(\lambda)\right] \\
& =\operatorname{rank}\left[\left[\left[\begin{array}{cc}
\lambda I-A_{1} & 0 \\
0 & \lambda I-A_{2} \\
0 & C_{2} \\
C_{1} & 0
\end{array}\right]\right]=n, \quad \forall \lambda\right.
\end{aligned}
$$

since, for $i=1,2$, the pair $\left(A_{i}, C_{i}\right)$ is observable. In a similar way the pair $(A, B)$ is proved to be reachable. To this aim, let

$$
T_{r}:=\left[\begin{array}{cccc}
I & 0 & 0 & 0 \\
0 & I & 0 & 0 \\
0 & -C_{2} & -\Delta_{21} D_{2} \Delta_{12}^{-1} & I \\
-C_{1} & 0 & I & -\Delta_{12} D_{1} \Delta_{21}^{-1}
\end{array}\right]
$$

Simple computations show that (recall also Lemma B.9)

$$
\begin{aligned}
P_{B}(\lambda) T_{r} & =\left[\begin{array}{lllc}
\lambda I-A & -B
\end{array}\right] T_{r} \\
& =\left[\begin{array}{cccc}
\lambda I-A_{1} & 0 & 0 & -B_{1} \\
0 & \lambda I-A_{2} & -B_{2} & 0
\end{array}\right]
\end{aligned}
$$

In view of the $P B H$ test the pair $(A, B)$ is reachable if and only if $\operatorname{rank}\left[P_{B}(\lambda)\right]=$ $\operatorname{rank}\left[T_{r} P_{B}(\lambda)\right]=n, \forall \lambda$. This condition holds thanks to the reachability of the pair $\left(A_{i}, B_{i}\right), i=1,2$.

Example 3.1 Let $G_{1}(s)=\bar{G}_{1}(s):=-1 /(s-1)$ and $G_{2}(s)=(s-1) /(s+1)$. The dynamic matrix of a particular realization of system $\Sigma$ is

$$
A=\left[\begin{array}{rr}
0 & -2 \\
-1 & -1
\end{array}\right]
$$

the eigenvalues of which are $\lambda_{1}=-2$ and $\lambda_{2}=\mathbf{1}$, so that $\Sigma$ is not internally stable. Consistently, the matrix

$$
T(u, v ; s)=\left[\begin{array}{cc}
(s+1) /(s+2) & (s-1) /(s+2) \\
-(s+1) /(s-1)(s+2) & (s+1) /(s+2)
\end{array}\right]
$$

does not belong to $R H_{\infty}$. On the contrary, if $G_{1}(s)=\widetilde{G}_{1}(s):=1 /(s+2)$, the dynamic matrix of a particular realization of $\Sigma$ is

$$
A=\left[\begin{array}{rr}
-1 & -2 \\
1 & -1
\end{array}\right]
$$

whose eigenvalues are $\lambda_{1,2}=-1 \pm j \sqrt{2}$, so that the system is internally stable. Consistently, the matrix

$$
T(u, v ; s)=\frac{1}{s^{2}+2 s+3}\left[\begin{array}{cc}
(s+1)(s+2) & (s-1)(s+2) \\
(s+1) & (s+1)(s+2)
\end{array}\right]
$$

is an element of $R H_{\infty}$.
This result can be specialized to the case where one of the two transfer functions $G_{1}(s)$ or $G_{2}(s)$ is stable, leading to the following theorem where $T_{i j}(u, v ; s), i=1,2$, $j=1,2$ denotes the transfer function from $v_{j}$ to $u_{i}$.

Theorem 3.2 Let $G_{2}(s)$ be stable. Then the system $\Sigma$ in fig. 3.3 is internally stable if and only if $T_{21}(u, v ; s) \in R H_{\infty}$.

Proof In view of Theorem 3.1, the present theorem is proved once it is shown that stability of $T_{21}(u, v ; s)$ implies stability of $T_{11}(u, v ; s), T_{12}(u, v ; s), T_{22}(u, v ; s)$. By exploiting Lemma B. 9 one has

$$
\begin{aligned}
I+G_{2}(s) T_{21}(u, v ; s) & =I+G_{2}(s) G_{1}(s)\left[I-G_{2}(s) G_{1}(s)\right]^{-1} \\
& =\left[I-G_{2}(s) G_{1}(s)\right]^{-1}=T_{11}(u, v ; s)
\end{aligned}
$$

so that, if $T_{21}(u, v ; s) \in R H_{\infty}$, also $T_{11}(u, v ; s) \in R H_{\infty}$. Similarly, by exploiting once more Lemma B.9,

$$
\begin{aligned}
I+T_{21}(u, v ; s) G_{2}(s) & =I+\left[I-G_{1}(s) G_{2}(s)\right]^{-1} G_{1}(s) G_{2}(s) \\
& =\left[I-G_{1}(s) G_{2}(s)\right]^{-1}=T_{22}(u, v ; s)
\end{aligned}
$$

so that, again, if $T_{21}(u, v ; s) \in R H_{\infty}, T_{22}(u, v ; s) \in R H_{\infty}$, as well. Finally, being $T_{12}(u, v ; s)=T_{11}(u, v ; s) G_{2}(s)$ then stability of $T_{11}(u, v ; s)$ entails stability of $T_{12}(u, v ; s)$.

Example 3.2 Again consider the functions $G_{1}(s)$ and $G_{2}(s)$ defined in Example 3.1. Being $G_{2}(s)$ stable, the internal stability of $\Sigma$ can be tested by checking the stability of $T_{21}(u, v ; s)$ only. When $G_{1}(s)=\bar{G}_{1}(s), T_{21}(u, v ; s)$ is not stable, whereas $T_{21}(u, v ; s) \in R H_{\infty}$ when $G_{1}(s)=\widetilde{G}_{1}(s)$, consistently with $\Sigma$ being stable only in the second case.

A further specialization of the above results can be found when both transfer functions $G_{1}(s)$ and $G_{2}(s)$ are stable.

Theorem 3.3 Assume that $G_{i}(s) \in R H_{\infty}, i=1,2$. Then the system $\Sigma$ depicted in fig. 3.3 is internally stable if and only if

$$
\operatorname{det}\left[I-G_{1}(s) G_{2}(s)\right] \neq 0, \quad \operatorname{Re}(s) \geq 0
$$

Proof Preliminarily notice that the transfer function of the system

$$
\Sigma_{a}:=\left[\begin{array}{cc|c}
A_{1} & B_{1} C_{2} & B_{1} D_{2} \\
0 & A_{2} & B_{2} \\
\hline-C_{1} & -D_{1} C_{2} & \Delta_{12}
\end{array}\right]
$$

is $I-G_{1}(s) G_{2}(s)$. By exploiting Lemma B. 9 and recalling Definition 2.14 it is easy to check that the dynamic matrix of $\Sigma_{a}^{-1}$ is precisely $A$. Moreover, any common root of the polynomial $\pi_{z t}(s)$ of the transmission zeros and $\pi_{p}(s)$ of the poles of system $I-G_{1}(s) G_{2}(s)$ must necessarily lie in the open left half plane, since such a system is stable.

Now suppose that $\operatorname{det}\left[I-G_{1}(s) G_{2}(s)\right] \neq 0, \operatorname{Re}(s) \geq 0$, so that, in view of Remark 2.7,

$$
\frac{\pi_{z t}(s)}{\pi_{p}(s)}=\operatorname{det}\left[I-G_{1}(s) G_{2}(s)\right] \neq 0, \quad \operatorname{Re}(s) \geq 0
$$

which implies that all transmission zeros of $\Sigma_{a}$ have negative real part. The stability of $\Sigma_{a}$ (notice that $A_{i}$ is stable, $i=1,2$, since $G_{i}(s) \in R H_{\infty}, i=1,2$, and the associated realizations are minimal) entails that any invariant zero of it which is not a transmission zero must lie in the open left half plane. Hence, thanks to Theorem
2.7, all the eigenvalues of $\Sigma_{a}^{-1}$ have negative real part, matrix $A$ is stable and in turn system $\Sigma$ is internally stable.

Conversely, if system $\Sigma$ is internally stable, then matrix $A$ is stable and system $\Sigma_{a}^{-1}$ is internally stable, so that, in view of Theorem 2.7 , all invariant zeros of system $\Sigma_{a}$ are in the open left half plane. Since the set of transmission zeros is contained in the set of invariant zeros (Theorem 2.4), all the transmission zeros of system $\Sigma_{a}$ have negative real part. Then, from Remark 2.7, it follows

$$
\operatorname{det}\left[I-G_{1}(s) G_{2}(s)\right] \neq 0, \quad \operatorname{Re}(s) \geq 0
$$

Example 3.3 Consider the functions $G_{1}(s)=\widetilde{G}_{1}(s)$ and $G_{2}(s)$ defined in Example 3.1. Both of them are stable and system $\Sigma$ is internally stable if and only if $\operatorname{det}\left[I-G_{1}(s) G_{2}(s)\right]$ is not zero in the closed right half plane. Actually, it is

$$
\operatorname{det}\left[I-\widetilde{G}_{1}(s) G_{2}(s)\right]=\frac{s^{2}+2 s+3}{(s+1)(s+2)}
$$

the zeros of which have negative real part, consistently with the (already established) internal stability of system $\Sigma$.

### 3.3 Double coprime factorizations

The definitions of right and left coprimeness for rational matrices which, together with some related properties, have been presented in Section 2.4, are now exploited in order to introduce an intermediate result on the way of presenting the so called parametrization of all stabilizing controllers in the next section.

Theorem 3.4 Let $G(s)$ be any proper rational matrix. Then there exist eight matrices $M(s), N(s), \hat{M}(s), \hat{N}(s), X(s), Y(s), \hat{X}(s), \hat{Y}(s)$, all belonging to $R H_{\infty}$, such that
a)

$$
G(s)=N(s) M^{-1}(s)=\hat{M}^{-1}(s) \hat{N}(s)
$$

b)

$$
\left[\begin{array}{rr}
\hat{X}(s) & -\hat{Y}(s) \\
-\hat{N}(s) & \hat{M}(s)
\end{array}\right]\left[\begin{array}{rr}
M(s) & Y(s) \\
N(s) & X(s)
\end{array}\right]=I
$$

The matrices $M(s), N(s), \hat{M}(s), \hat{N}(s)$ constitute a double coprime factorization of $G(s)$.

Proof First notice that, if the statement of the theorem is correct, the matrices $M(s)$ and $N(s)$ are right coprime in view of Lemma 2.5 , while $\hat{M}(s)$ and $\hat{N}(s)$ are left coprime thanks to the "left version" of the same lemma.

Now let $G(s)=\Sigma(A, B, C, D)$ and assume that the triple $(A, B, C)$ is stabilizable and detectable. Further, let $F$ and $H$ be any two matrices such that $(A+B F)$ and $(A+H C)$ are stable and set $\bar{C}:=C+D F, \bar{B}:=B+H D$. Define

$$
M(s):=\left[\begin{array}{c|c}
A+B F & B  \tag{3.1}\\
\hline F & I
\end{array}\right] \quad, \quad N(s):=\left[\begin{array}{c|c}
A+B F & B \\
\hline \bar{C} & D
\end{array}\right]
$$

$$
\begin{array}{ll}
X(s):=\left[\begin{array}{c|c}
A+B F & H \\
\hline-\bar{C} & I
\end{array}\right] \quad, \quad Y(s):=\left[\begin{array}{c|c}
A+B F & H \\
\hline-F & 0
\end{array}\right] \\
\hat{M}(s):=\left[\begin{array}{c|c}
A+H C & H \\
\hline C & I
\end{array}\right], \quad \hat{N}(s):=\left[\begin{array}{cc}
A+H C & \bar{B} \\
\hline C & D
\end{array}\right] \\
\hat{X}(s):=\left[\begin{array}{cc}
A+H C & \bar{B} \\
\hline-F & I
\end{array}\right] \quad, \quad \hat{Y}(s):=\left[\begin{array}{cc}
A+H C & H \\
\hline-F & 0
\end{array}\right] \tag{3.4}
\end{array}
$$

Point a) It is obvious that the four matrices $M(s), N(s), \hat{M}(s), \hat{N}(s)$ belong to $R H_{\infty}$. Moreover, $M(s)$ and $\hat{M}(s)$ have inverses since

$$
\lim _{s \rightarrow \infty} M(s)=I, \quad \lim _{s \rightarrow \infty} \hat{M}(s)=I
$$

Now consider system $\Sigma_{1}$ (with input $v$ and outputs $y$ and $u$ ) defined by

$$
\begin{aligned}
\dot{x} & =(A+B F) x+B v \\
y & =\bar{C} x+D v \\
u & =F x+v
\end{aligned}
$$

Thus, $\Sigma_{1}$ is nothing but system $G(s)$ after the control law $u=F x+v$ has been implemented. It follows

$$
\begin{aligned}
u_{L 0} & =M(s) v_{L} \\
y_{L 0} & =N(s) v_{L}
\end{aligned}
$$

so that $y_{L 0}=G(s) u_{L 0}=N(s) M^{-1}(s) u_{L 0}$. Hence

$$
\begin{equation*}
G(s)=N(s) M^{-1}(s) \tag{3.5}
\end{equation*}
$$

Consider system $\Sigma_{2}$ (with inputs $u$ and $y$ and output $\eta$ ) defined by

$$
\begin{aligned}
& \dot{\vartheta}=A \vartheta+B u+H \eta \\
& \eta=C \vartheta+D u-y
\end{aligned}
$$

Thus, $\Sigma_{2}$ is nothing but a state observer for the system with transfer function $G(s)$, so that $\eta$ is the output observation error, namely

$$
\eta=C(\vartheta-x)
$$

which is well known not to depend on $u$, since $\dot{\vartheta}-\dot{x}=(A+H C)(\vartheta-x)$. Therefore,

$$
\begin{aligned}
\eta_{L 0}= & {\left[C(s I-(A+H C))^{-1} \bar{B}+D\right] u_{L}-} \\
& -\left[C(s I-(A+H C))^{-1} H+I\right] y_{L 0} \\
= & \hat{N}(s) u_{L}-\hat{M}(s) y_{L 0} \\
= & {[\hat{N}(s)-\hat{M}(s) G(s)] u_{L}=0 }
\end{aligned}
$$

which implies

$$
\begin{equation*}
G(s)=\hat{M}^{-1}(s) \hat{N}(s) \tag{3.6}
\end{equation*}
$$

Point $a$ ) is then proved in view of eqs. (3.5),(3.6).

Point b) Consider system $\Sigma_{3}$ (with inputs $u$ and $y$ and output $\gamma$ ) defined by

$$
\begin{aligned}
\dot{\mu} & =(A+H C) \mu+\bar{B} u-H y \\
\gamma & =u-F \mu
\end{aligned}
$$

so that $\gamma_{L 0}=\hat{X}(s) u_{L}-\hat{Y}(s) y_{L}$. Then in the series connection $\Sigma_{3} \Sigma_{1}$ it is

$$
\gamma_{L 0}=\hat{X}(s) u_{L 0}-\hat{Y}(s) y_{L 0}=[\hat{X}(s) M(s)-\hat{Y}(s) N(s)] v_{L}
$$

However, $\dot{\mu}-\dot{x}=(A+H C)(\mu-x)$ and $\gamma=v-F(\mu-x)$ which implies that the transfer function from $v$ to $\gamma$ is $I$, that is

$$
\begin{equation*}
\hat{X}(s) M(s)-\hat{Y}(s) N(s)=I \tag{3.7}
\end{equation*}
$$

Now consider system $\Sigma_{4}$ (with input $w$ and outputs $y$ and $u$ ) defined by

$$
\begin{aligned}
\dot{z} & =(A+B F) z+H w \\
y & =\bar{C} z-w \\
u & =F z
\end{aligned}
$$

so that $y_{L 0}=-X(s) w_{L}$ and $u_{L 0}=-Y(s) w_{L}$. Then in the series connection $\Sigma_{2} \Sigma_{4}$ it is

$$
\eta_{L 0}=\hat{N}(s) u_{L 0}-\hat{M}(s) y_{L 0}=[-\hat{N}(s) Y(s)+\hat{M}(s) X(s)] w_{L}
$$

However, $\dot{z}-\dot{\vartheta}=(A+H C)(z-\vartheta)$ and $\eta=w-C(z-\vartheta)$ which implies that the transfer function from $w$ to $\eta$ is $I$, that is

$$
\begin{equation*}
-\hat{N}(s) Y(s)+\hat{M}(s) X(s)=I \tag{3.8}
\end{equation*}
$$

Finally, in the series connection $\Sigma_{3} \Sigma_{4}$ it is

$$
\gamma_{L 0}=\hat{X}(s) u_{L 0}-\hat{Y}(s) y_{L 0}=[-\hat{X}(s) Y(s)+\hat{Y}(s) X(s)] w_{L}
$$

However it is $\dot{z}-\dot{\mu}=(A+H C)(z-\mu)$ and $\gamma=F(z-\mu)$ which implies that the transfer function from $w$ to $\gamma$ is 0 ,

$$
\begin{equation*}
-\hat{X}(s) Y(s)+\hat{Y}(s) X(s)=0 \tag{3.9}
\end{equation*}
$$

From eqs. (3.5),(3.6) it follows that $\hat{M}(s) N(s)-\hat{N}(s) M(s)=0$ which, together with eqs. (3.7)-(3.9), proves point $b$ ).

Remark 3.1 The proof of Theorem 3.4 outlines the fact that, corresponding to a given function $G(s)$, there exists an infinite number of distinct double coprime factorizations in $R H_{\infty}$ (recall that matrices $F$ and $H$ can be chosen in an almost arbitrary way).

Example 3.4 Consider the rational matrix

$$
G(s)=\left[\begin{array}{cc}
\left(s^{2}+s-1\right) /\left(s^{3}-s\right) & s /\left(s^{2}-1\right) \\
\left(2 s^{2}-1\right) /\left(s^{3}-s\right) & s^{2} /\left(s^{2}-1\right)
\end{array}\right]
$$

which admits the quadruple

$$
\begin{array}{ll}
A=\left[\begin{array}{lll}
0 & 1 & 0 \\
1 & 0 & 0 \\
0 & 0 & 0
\end{array}\right], \quad B=\left[\begin{array}{ll}
0 & 1 \\
1 & 0 \\
1 & 0
\end{array}\right] \\
C=\left[\begin{array}{lll}
1 & 0 & 1 \\
0 & 1 & 1
\end{array}\right], \quad D=\left[\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right]
\end{array}
$$

as a minimal realization. Chosen

$$
F:=\left[\begin{array}{rrr}
0 & 0 & -1 \\
-1 & -2 & 0
\end{array}\right], \quad H:=\left[\begin{array}{rr}
-4 & 0 \\
-4 & 0 \\
1 & 0
\end{array}\right]
$$

it is easy to check the stability of matrices $A+B F$ and $A+H C$ and then computing the eight matrices referred to in Theorem 3.4. It results

$$
\begin{aligned}
& M(s)=\left[\begin{array}{cc}
s /(s+1) & 0 \\
-s(2 s+1) /(s+1)\left(s^{2}+s+1\right) & \left(s^{2}-1\right) /\left(s^{2}+s+1\right)
\end{array}\right] \\
& N(s)=\left[\begin{array}{cc}
\left(s^{2}+1\right) /(s+1)\left(s^{2}+s+1\right) & s /\left(s^{2}+s+1\right) \\
1 /\left(s^{2}+s+1\right) & s^{2} /\left(s^{2}+s+1\right)
\end{array}\right] \\
& X(s)=\left[\begin{array}{cc}
\left(s^{3}+5 s^{2}+s-5\right) /(s+1)\left(s^{2}+s+1\right) & 0 \\
-(9 s+5) /\left(s^{2}+s+1\right) & 1
\end{array}\right] \\
& Y(s)=\left[\begin{array}{cc}
1 /(s+1) & 0 \\
-\left(12 s^{2}+26 s+13\right) /(s+1)\left(s^{2}+s+1\right) & 0
\end{array}\right] \\
& \hat{M}(s)=\left[\begin{array}{cc}
s(s-1) /(s+1)^{2} & 0 \\
-(3 s+1) /(s+1)^{2} & 1
\end{array}\right] \\
& \hat{N}(s)=\left[\begin{array}{cc}
\left(s^{2}+s-1\right) /(s+1)^{3} & s^{2} /(s+1)^{3} \\
\left(2 s^{2}+3 s\right) /(s+1)^{3} & \left(s^{3}+3 s^{2}+s\right) /(s+1)^{3}
\end{array}\right] \\
& \hat{X}(s)=\left[\begin{array}{cc}
\left(s^{3}+4 s^{2}+7 s+5\right) /(s+1)^{3} & s /(s+1)^{3} \\
\left(2 s^{2}-5 s-13\right) /(s+1)^{3} & \left(s^{3}+4 s^{2}-4 s-1\right) /(s+1)^{3}
\end{array}\right] \\
& \hat{Y}(s)=\left[\begin{array}{cc}
(s-1) /(s+1)^{2} & 0 \\
-12 s /(s+1)^{2} & 0
\end{array}\right]
\end{aligned}
$$

These matrices verify Theorem 3.4.
Theorem 3.4 allows one stating an equivalence condition between internal and external stability for a feedback system as the one depicted in fig. 3.3 which is different from the condition presented in Theorem 3.1. As done in the preceding section, it is assumed that a minimal realization of the two systems $G_{1}(s)$ and $G_{2}(s)$ is considered and their feedback connection is well defined (the matrix ( $I-D_{1} D_{2}$ ) is nonsingular, $\left.D_{i}:=\lim _{s \rightarrow \infty} G_{i}(s), i=1,2\right)$. For $i=1,2$, let $N_{i}(s)$ and $M_{i}(s)$ be elements of $R H_{\infty}$ such that $G_{i}(s)=N_{i}(s) M_{i}^{-1}(s)$ with $N_{i}(s)$ and $M_{i}(s)$ right coprime. Then the system in fig. 3.3 can be represented as in fig. 3.4 as well. With reference to this figure let $z:=\left[\begin{array}{ll}z_{1}^{\prime} & z_{2}^{\prime}\end{array}\right]^{\prime}, u:=\left[\begin{array}{ll}u_{1}^{\prime} & u_{2}^{\prime}\end{array}\right]^{\prime}$ and $v:=\left[\begin{array}{ll}v_{1}^{\prime} & v_{2}^{\prime}\end{array}\right]^{\prime}$. Then the following result holds.

Theorem 3.5 The system in fig. 3.4 is internally stable if and only if the transfer function $T(z, v ; s)$ from the input $v$ to the output $z$ is stable.

Proof In view of Theorem 3.1, the proof consists in showing that the matrix $T(z, v ; s)$ is stable if and only if the transfer function $T(u, v ; s)$ from the input $v$ to the output $u$ is stable.

Sufficiency With reference to fig. 3.4 it is

$$
u_{L 0}=T(z, v ; s) v_{L}
$$



Figure 3.4: Feedback connection of two systems

$$
\begin{aligned}
& =\left[\begin{array}{cc}
0 & N_{2}(s) \\
N_{1}(s) & 0
\end{array}\right] z_{L 0}+v_{L} \\
& =\left[\left[\begin{array}{cc}
0 & N_{2}(s) \\
N_{1}(s) & 0
\end{array}\right] T(z, v ; s)+I\right] v_{L}
\end{aligned}
$$

Therefore, if $T(z, v ; s) \in R H_{\infty}$ also $T(u, v ; s) \in R H_{\infty}$, since, for $i=1,2, N_{i}(s) \in$ $R H_{\infty}$.

Necessity Being $N_{i}(s)$ and $M_{i}(s), i=1,2$, right coprime, there exist matrices $X_{i}(s)$ and $Y_{i}(s), i=1,2$, which are elements of $R H_{\infty}$ and such that

$$
\begin{aligned}
& {\left[\begin{array}{cc}
X_{1}(s) & 0 \\
0 & X_{2}(s)
\end{array}\right]\left[\begin{array}{cc}
M_{1}(s) & 0 \\
0 & M_{2}(s)
\end{array}\right]+} \\
& +\left[\begin{array}{cc}
Y_{1}(s) & 0 \\
0 & Y_{2}(s)
\end{array}\right]\left[\begin{array}{cc}
N_{1}(s) & 0 \\
0 & N_{2}(s)
\end{array}\right]=I
\end{aligned}
$$

so that

$$
\begin{aligned}
& {\left[\begin{array}{cc}
X_{1}(s) & 0 \\
0 & X_{2}(s)
\end{array}\right]\left[\begin{array}{cc}
M_{1}(s) & 0 \\
0 & M_{2}(s)
\end{array}\right] z_{L 0}+} \\
& +\left[\begin{array}{cc}
Y_{1}(s) & 0 \\
0 & Y_{2}(s)
\end{array}\right]\left[\begin{array}{cc}
N_{1}(s) & 0 \\
0 & N_{2}(s)
\end{array}\right] z_{L 0}=z_{L 0}
\end{aligned}
$$

from which it follows

$$
\begin{aligned}
& {\left[\begin{array}{cc}
X_{1}(s) & 0 \\
0 & X_{2}(s)
\end{array}\right] u_{L 0}+\left[\begin{array}{cc}
0 & Y_{1}(s) \\
Y_{2}(s) & 0
\end{array}\right]\left(u_{L 0}-v_{L}\right)=} \\
& =\left[\left[\begin{array}{cc}
X_{1}(s) & Y_{1}(s) \\
Y_{2}(s) & X_{2}(s)
\end{array}\right] T(u, v ; s)-\left[\begin{array}{cc}
0 & Y_{1}(s) \\
Y_{2}(s) & 0
\end{array}\right]\right] v_{L} \\
& =T(z, v ; s) v_{L}
\end{aligned}
$$

Being $X_{i}(s), Y_{i}(s), i=1,2$, and $T(u, v ; s)$ all elements of $R H_{\infty}$, also $H(s)$ must be such.

Thanks to Theorem 3.5 it is possible to prove a result useful for the discussion in the next section. It clarifies the circumstances under which a system with transfer function $G_{1}(s)$ is internally stabilized by a system with transfer function $G_{2}(s)$ when they are feedback connected as shown in fig. 3.4, where, for $i=1,2$, matrices $N_{i}(s)$ and $M_{i}(s)$ possess the above mentioned properties.

Theorem 3.6 With reference to fig. 3.4, system $G_{2}(s)$ internally stabilizes system $G_{1}(s)$, that is, the resulting system is internally stable, if and only if $G^{-1}(s) \in R H_{\infty}$, where

$$
G(s):=\left[\begin{array}{ll}
M_{1}(s) & N_{2}(s) \\
N_{1}(s) & M_{2}(s)
\end{array}\right]
$$

Proof Preliminarily observe that matrix $G(s)$ is nonsingular, so that the statement makes sense. Indeed, one can write

$$
\left[\begin{array}{cc}
M_{1}(s) & N_{2}(s) \\
N_{1}(s) & M_{2}(s)
\end{array}\right]=\left[\begin{array}{cc}
I & G_{2}(s) \\
G_{1}(s) & I
\end{array}\right]\left[\begin{array}{cc}
M_{1}(s) & 0 \\
0 & M_{2}(s)
\end{array}\right]
$$

For the first matrix on the right hand side of this equation it is

$$
\begin{aligned}
\lim _{s \rightarrow \infty} \operatorname{det}\left[\left[\begin{array}{cc}
I & G_{2}(s) \\
G_{1}(s) & I
\end{array}\right]\right] & =\lim _{s \rightarrow \infty} \operatorname{det}\left[I-G_{1}(s) G_{2}(s)\right] \\
& =\operatorname{det}\left[I-D_{1} D_{2}\right]
\end{aligned}
$$

having exploited the Schur formula for the determinant of a block matrix (see Lemma B.13). Thus, matrix $G(s)$ is nonsingular since it is the product of two nonsingular matrices (recall that, for $i=1,2, M_{i}^{-1}(s)$ exists).

By taking into account fig. 3.4 one has

$$
\begin{aligned}
{\left[\begin{array}{rr}
M_{1}(s) & -N_{2}(s) \\
-N_{1}(s) & M_{2}(s)
\end{array}\right] z_{L 0} } & =\left[\begin{array}{c}
M_{1}(s) z_{1 L 0}-N_{2}(s) z_{2 L 0} \\
-N_{1}(s) z_{1 L 0}+M_{2}(s) z_{2 L 0}
\end{array}\right] \\
& =\left[\begin{array}{c}
u_{1 L 0}+v_{1 L}-u_{1 L 0} \\
v_{2 L}-u_{2 L 0}+u_{2 L 0}
\end{array}\right]=v_{L}
\end{aligned}
$$

so that, thanks to Theorem 3.5, $G_{2}(s)$ internally stabilizes $G_{1}(s)$ if and only if

$$
T(z, v ; s):=\left[\begin{array}{rr}
M_{1}(s) & -N_{2}(s) \\
-N_{1}(s) & M_{2}(s)
\end{array}\right]^{-1} \in R H_{\infty}
$$

Finally, $T(z, v ; s) \in R H_{\infty}$ if and only if $G^{-1}(s) \in R H_{\infty}$. In fact,

$$
T(z, v ; s)=\left[\begin{array}{rr}
I & 0 \\
0 & -I
\end{array}\right] G^{-1}(s)\left[\begin{array}{rr}
I & 0 \\
0 & -I
\end{array}\right]
$$

Example 3.5 Consider the functions $G_{1}(s)$ and $G_{2}(s)$ defined in Example 3.1. Being $G_{2}(s)$ stable one can set (see the proof of Theorem 3.4)

$$
N_{2}(s)=\frac{s-1}{s+1}, \quad M_{2}(s)=1
$$

whereas for $G_{1}(s)=\bar{G}_{1}(s)$, the choice $F=-3$ entails (see the proof of Theorem 3.4)

$$
N_{1}(s)=\bar{N}_{1}(s)=-\frac{1}{s+2}, \quad M_{1}(s)=\bar{M}_{1}(s)=\frac{s-1}{s+2}
$$

Finally, being $\widetilde{G}_{1}(s)$ stable, for $G_{1}(s)=\widetilde{G}_{1}(s)$, one can set

$$
N_{1}(s)=\widetilde{N}_{1}(s)=\frac{1}{s+2}, \quad M_{1}(s)=\widetilde{M}_{1}(s)=1
$$



Figure 3.5: Feedback connection of two systems

Thus,

$$
\begin{aligned}
\bar{G}^{-1}(s) & =\left[\begin{array}{cc}
(s-1) /(s+2) & (s-1) /(s+1) \\
-1 /(s+2) & 1
\end{array}\right]^{-1} \\
& =\left[\begin{array}{cc}
(s+1) /(s-1) & -1 \\
(s+1) /\left(s^{2}+s-2\right) & (s+1) /(s+2)
\end{array}\right]
\end{aligned}
$$

which does not belong to $R H_{\infty}$, consistently with $\bar{G}_{1}(s)$ not being stabilized by $G_{2}(s)$, whereas

$$
\begin{aligned}
\widetilde{G}^{-1}(s) & =\left[\begin{array}{cc}
1 & (s-1) /(s+1) \\
1 /(s+2) & 1
\end{array}\right]^{-1} \\
& =\left[\begin{array}{cc}
\left(s^{2}+3 s+2\right) /\left(s^{2}+2 s+3\right) & -\left(s^{2}+s-2\right) /\left(s^{2}+2 s+3\right) \\
-(s+1) /\left(s^{2}+2 s+3\right) & \left(s^{2}+3 s+2\right) /\left(s^{2}+2 s+3\right)
\end{array}\right]
\end{aligned}
$$

belongs to $R H_{\infty}$, consistently with $\widetilde{G}_{1}(s)$ being stabilized by $G_{2}(s)$.

### 3.4 The set of stabilizing controllers

The so called parametrization of all stabilizing controllers is presented in this section. This result is particularly significant as it allows to completely characterize the set of all linear, finite dimensional and time invariant controllers $K(s)$ which stabilize a given system with transfer function $G(s)$, if connected to it according to the blockscheme of fig. 3.5. It will be shown that the system in fig. 3.5 is internally stable if and only if the controller transfer function can be given a very simple form in terms of a free parameter which, besides ensuring the well-posedness of the feedback connection, has to comply with a single constraint only, namely being an element of the space $R H_{\infty}$.

Obviously, there exist stabilizing controllers for system $G(s):=C(s I-A)^{-1} B+D$ only provided that the pair $(A, B)$ is stabilizable and the pair $(A, C)$ is detectable. Therefore, such an unavoidable assumptions will be done throughout the subsequent discussion.

Consistently with Theorem 3.4, let the eight matrices $N(s), M(s), \hat{N}(s), \hat{M}(s)$, $X(s), Y(s), \hat{X}(s), \hat{Y}(s)$ specify a double coprime factorization in $R H_{\infty}$ of the transfer function $G(s)$. These matrices are derived according to eqs. (3.1)-(3.4) and are such that

$$
G(s)=N(s) M^{-1}(s)=\hat{M}^{-1}(s) \hat{N}(s)
$$

and

$$
\left[\begin{array}{rr}
\hat{X}(s) & -\hat{Y}(s)  \tag{3.10}\\
-\hat{N}(s) & \hat{M}(s)
\end{array}\right]\left[\begin{array}{cc}
M(s) & Y(s) \\
N(s) & X(s)
\end{array}\right]=I
$$

It is now possible to state the following result.
Theorem 3.7 (Youla parametrization) The set of all rational and proper transfer functions $K(s)$ which internally stabilize $G(s)$ is defined by controllers of the form

$$
\begin{align*}
K(s) & =[Y(s)-M(s) Q(s)][X(s)-N(s) Q(s)]^{-1} \\
& =[\hat{X}(s)-Q(s) \hat{N}(s)]^{-1}[\hat{Y}(s)-Q(s) \hat{M}(s)] \tag{3.11}
\end{align*}
$$

where the matrix $Q(s)$ is any element of $R H_{\infty}$ such that

$$
\begin{equation*}
\operatorname{det}\left[I-D \lim _{s \rightarrow \infty} Q(s)\right] \neq 0 \tag{3.12}
\end{equation*}
$$

Proof Preliminarily notice that the existence of the inverse matrices in eq. (3.11) is ensured by eqs. (3.1)-(3.4) and (3.12) (also recall Lemma B.8).

For any matrix $Q(s)$, eq. (3.10) implies that

$$
\begin{align*}
I= & {\left[\begin{array}{cc}
I & Q(s) \\
0 & I
\end{array}\right]\left[\begin{array}{cc}
I & -Q(s) \\
0 & I
\end{array}\right] } \\
= & {\left[\begin{array}{cc}
I & Q(s) \\
0 & I
\end{array}\right]\left[\begin{array}{cc}
\hat{X}(s) & -\hat{Y}(s) \\
-\hat{N}(s) & \hat{M}(s)
\end{array}\right] } \\
& \cdot\left[\begin{array}{cc}
M(s) & Y(s) \\
N(s) & X(s)
\end{array}\right]\left[\begin{array}{cc}
I & -Q(s) \\
0 & I
\end{array}\right] \\
= & {\left[\begin{array}{cc}
\hat{X}(s)-Q(s) \hat{N}(s) & -[\hat{Y}(s)-Q(s) \hat{M}(s)] \\
-\hat{N}(s) & \hat{M}(s)
\end{array}\right] } \\
& \cdot\left[\begin{array}{cc}
M(s) & -M(s) Q(s)+Y(s) \\
N(s) & -N(s) Q(s)+X(s)
\end{array}\right] \tag{3.13}
\end{align*}
$$

The $(1,2)$ block of the product of the two last matrices in eq. (3.13) is zero, i.e.

$$
\begin{align*}
{[\hat{X}(s)-Q(s) \hat{N}(s)][-M(s) Q(s)+Y(s)] } & - \\
- & {[\hat{Y}(s)-Q(s) \hat{M}(s)][-N(s) Q(s)+X(s)] } \tag{3.14}
\end{align*}=0
$$

so that the second equality sign in eq. (3.11) is proved.
Now the system with transfer function

$$
K(s):=[Y(s)-M(s) Q(s)][X(s)-N(s) Q(s)]^{-1}
$$

where $Q(s) \in R H_{\infty}$ satisfies condition (3.12), is proved to internally stabilize $G(s)$. To this aim, let $U(s):=Y(s)-M(s) Q(s), V(s):=X(s)-N(s) Q(s), \hat{U}(s):=$ $\hat{Y}(s)-Q(s) \hat{M}(s), \hat{V}(s):=\hat{X}(s)-Q(s) \hat{N}(s)$, so that $K(s)=U(s) V^{-1}(s)$. From eq. (3.10) it follows

$$
\left[\begin{array}{rr}
\hat{V}(s) & -\hat{U}(s) \\
-\hat{N}(s) & \hat{M}(s)
\end{array}\right]\left[\begin{array}{cc}
M(s) & U(s) \\
N(s) & V(s)
\end{array}\right]=I
$$

that is

$$
\left[\begin{array}{cc}
M(s) & U(s) \\
N(s) & V(s)
\end{array}\right]^{-1}=\left[\begin{array}{rr}
\hat{V}(s) & -\hat{U}(s) \\
-\hat{N}(s) & \hat{M}(s)
\end{array}\right]
$$

The matrix on the right hand side of this last equation belongs to $R H_{\infty}$ (each block in it is an element of $R H_{\infty}$ ), so it does the inverse on the left hand side. Hence, thanks to Theorem 3.6, the system in fig. 3.5 is internally stable.

Conversely, assume that $K(s)$ stabilizes $G(s)$. In order that the feedback system in fig. 3.5 is well defined, matrix $K(s)$ must satisfy the condition

$$
\begin{equation*}
\operatorname{det}[I-D \bar{K}] \neq 0 \tag{3.15}
\end{equation*}
$$

where $\bar{K}:=\lim _{s \rightarrow \infty} K(s)$. Let $U(s)$ and $V(s)$ be two right coprime elements of $R H_{\infty}$ such that $K(s)=U(s) V^{-1}(s)$ and consider the identity

$$
\begin{gather*}
{\left[\begin{array}{rr}
\hat{X}(s) & -\hat{Y}(s) \\
-\hat{N}(s) & \hat{M}(s)
\end{array}\right]\left[\begin{array}{rr}
M(s) & U(s) \\
N(s) & V(s)
\end{array}\right]=} \\
=\left[\begin{array}{rr}
I & \hat{X}(s) U(s)-\hat{Y}(s) V(s) \\
0 & -\hat{N}(s) U(s)+\hat{M}(s) V(s)
\end{array}\right] \tag{3.16}
\end{gather*}
$$

which follows from eq. (3.10). The two matrices on the left hand side of eq. (3.16) have inverses in $R H_{\infty}$, the first one thanks to eq. (3.10), the second one thanks to Theorem 3.6. Hence also the matrix on the right hand side has inverse in $R H_{\infty}$, so that $P^{-1}(s):=[-\hat{N}(s) U(s)+\hat{M}(s) V(s)]^{-1} \in R H_{\infty}$. Therefore the matrix $Q(s):=$ $-[\hat{X}(s) U(s)-\hat{Y}(s) V(s)] P^{-1}(s)$ belongs to $R H_{\infty}$, too. Now notice that from eqs. (3.1)-(3.4) it follows

$$
\begin{aligned}
I-D \lim _{s \rightarrow \infty} Q(s)= & I+D \lim _{s \rightarrow \infty}[[\hat{X}(s) U(s)-\hat{Y}(s) V(s)] \\
& \left.\cdot[-\hat{N}(s) U(s)+\hat{M}(s) V(s)]^{-1}\right] \\
= & I+D \lim _{s \rightarrow \infty}\left[\left[\hat{X}(s) U(s) V^{-1}(s)-\hat{Y}(s)\right] V(s) V^{-1}(s)\right. \\
& \left.\cdot\left[-\hat{N}(s) U(s) V^{-1}(s)+\hat{M}(s)\right]^{-1}\right] \\
= & I+D \bar{K}(I-D \bar{K})^{-1} \\
= & (I-D \bar{K})^{-1}
\end{aligned}
$$

where the last equality sign is due to Lemma B.9. Hence $Q(s)$ satisfies condition (3.12) in view of eq. (3.15).

From eq. (3.16) one obtains

$$
\begin{gathered}
{\left[\begin{array}{cc}
M(s) & Y(s) \\
N(s) & X(s)
\end{array}\right]\left[\begin{array}{rr}
\hat{X}(s) & -\hat{Y}(s) \\
-\hat{N}(s) & \hat{M}(s)
\end{array}\right]\left[\begin{array}{cc}
M(s) & U(s) \\
N(s) & V(s)
\end{array}\right]=} \\
\quad=\left[\begin{array}{cc}
M(s) & Y(s) \\
N(s) & X(s)
\end{array}\right]\left[\begin{array}{cc}
I & -Q(s) P(s) \\
0 & P(s)
\end{array}\right]
\end{gathered}
$$

which, in view of eq. (3.10), becomes

$$
\left[\begin{array}{cc}
M(s) & U(s) \\
N(s) & V(s)
\end{array}\right]=\left[\begin{array}{cc}
M(s) & Y(s) \\
N(s) & X(s)
\end{array}\right]\left[\begin{array}{cc}
I & -Q(s) P(s) \\
0 & P(s)
\end{array}\right]
$$

This last equation implies

$$
\begin{aligned}
& U(s)=[Y(s)-M(s) Q(s)] P(s) \\
& V(s)=[X(s)-N(s) Q(s)] P(s)
\end{aligned}
$$

so that

$$
\begin{aligned}
K(s) & =U(s) V^{-1}(s) \\
& =[Y(s)-M(s) Q(s)][X(s)-N(s) Q(s)]^{-1}
\end{aligned}
$$

according to eq. (3.11).
The controller $K_{0}(s)$ corresponding to the somehow most natural choice of $Q(s)$, namely $Q(s)=0$, is usually referred to as the central controller. In view of eq. (3.11) such a controller is given by

$$
K_{0}(s)=Y(s) X^{-1}(s)=\hat{X}^{-1}(s) \hat{Y}(s)
$$

and therefore only depends on the particular double coprime factorization of $G(s)$ which has been selected. Further, in spite of corresponding to the simplest matrix $Q(s) \in R H_{\infty}$ which satisfies condition (3.12), the central controller is not, in general, the lowest order stabilizing controller.

Example 3.6 Let $G(s)=1 / s$ and choose $F=H=-1$, thus obtaining $N(s)=\hat{N}(s)=$ $1 /(s+1), M(s)=\hat{M}(s)=s /(s+1), X(s)=\hat{X}(s)=(s+2) /(s+1), Y(s)=\hat{Y}(s)=-1 /(s+1)$. From eq. (3.11) it follows

$$
K(s)=-\frac{1+s Q(s)}{s+2-Q(s)}
$$

so that, corresponding to the choice $Q(s)=0$ it is $K_{0}(s)=-1 /(s+2)$, while, if $Q(s)=1$, it is $K(s)=-1$ (note that the first controller, the central one, is a dynamic system of order 1 , while the second one is purely algebraic). Vice-versa, by making use of the expressions for $Q(s)$ given in the proof of Theorem 3.7, one gets

$$
Q(s)=-\frac{V(s)+(s+2) U(s)}{s V(s)-U(s)}
$$

where $U(s)$ and $V(s)$ are such that $K(s)=U(s) V^{-1}(s)$. As an example, if $K(s)=-1 /(s+1)$, then $Q(s)=1 /\left(s^{2}+s+1\right)$.

The set of all stabilizing controllers can be further analyzed in order to enlighten interesting connections with other results of linear control theory. A preliminary result is presented in the forthcoming lemma, where reference is made to a double coprime factorization of $G(s)$ satisfying eq. (3.10).

Lemma 3.1 The set of proper rational transfer functions which internally stabilize $G(s)$ is given by

$$
K(s)=K_{0}(s)-\hat{X}^{-1}(s) Q(s)\left[I-X^{-1}(s) N(s) Q(s)\right]^{-1} X^{-1}(s)
$$

where $K_{0}(s)=Y(s) X^{-1}(s)=\hat{X}^{-1}(s) \hat{Y}(s)$ and $Q(s) \in R H_{\infty}$ is such that

$$
\operatorname{det}\left[I-D \lim _{s \rightarrow \infty} Q(s)\right] \neq 0
$$

Proof By exploiting the results in Lemma B. 9 and the expression of a generic stabilizing controller (Theorem 3.7), one obtains

$$
\begin{aligned}
K(s) & =[Y(s)-M(s) Q(s)][X(s)-N(s) Q(s)]^{-1} \\
& =[Y(s)-M(s) Q(s)]\left\{\left[I-N(s) Q(s) X^{-1}(s)\right] X(s)\right\}^{-1}
\end{aligned}
$$



Figure 3.6: The structure of the stabilizing controllers

$$
\begin{aligned}
= & {[Y(s)-M(s) Q(s)] X^{-1}(s)\left[I-N(s) Q(s) X^{-1}(s)\right]^{-1} } \\
= & {\left[K_{0}(s)-M(s) Q(s) X^{-1}(s)\right]\{I+N(s) Q(s) \cdot} \\
& \left.\cdot\left[I-X^{-1}(s) N(s) Q(s)\right]^{-1} X^{-1}(s)\right\} \\
= & K_{0}(s)+K_{0}(s) N(s) Q(s)\left[I-X^{-1}(s) N(s) Q(s)\right]^{-1} . \\
& \cdot X^{-1}(s)-M(s) Q(s) X^{-1}(s) \cdot \\
= & \cdot\left\{I+N(s) Q(s)\left[I-X^{-1}(s) N(s) Q(s)\right]^{-1} X^{-1}(s)\right\} \\
& \cdot X^{-1}(s)-M(s) Q(s) K_{0}(s) N(s) Q(s)\left[I-X^{-1}(s)[I-N(s) N(s) Q(s)]^{-1} .\right. \\
= & K_{0}(s)+K_{0}(s) N(s) Q(s)\left[I-X^{-1}(s) N(s) Q(s)\right]^{-1} . \\
& \cdot X^{-1}(s)-M(s) Q(s)\left[I-X^{-1}(s) N(s) Q(s)\right]^{-1} X^{-1}(s) \\
= & K_{0}(s)+\left[K_{0}(s) N(s)-M(s)\right] Q(s) . \\
& \cdot\left[I-X^{-1}(s) N(s) Q(s)\right]^{-1} X^{-1}(s)
\end{aligned}
$$

Thanks to eq. (3.10) it is

$$
\begin{aligned}
K_{0}(s) N(s)-M(s) & =\hat{X}^{-1}(s) \hat{Y}(s) N(s)-M(s) \\
& =\hat{X}^{-1}(s)[\hat{Y}(s) N(s)-\hat{X}(s) M(s)] \\
& =\hat{X}^{-1}(s)
\end{aligned}
$$

so that $K(s)$ is given by

$$
K(s)=K_{0}(s)-\hat{X}^{-1}(s) Q(s)\left[I-X^{-1}(s) N(s) Q(s)\right]^{-1} X^{-1}(s)
$$

and the lemma is proved.
The form of $K(s)$ as given by Lemma 3.1 shows that the generic stabilizing controller is constituted by two subsystems in feedback connection. The first subsystem, denoted with $R_{f}(s)$ in fig. 3.6, only depends on the particular double coprime factorization which has been selected for $G(s)$, while the second one is simply constituted by the system with transfer function $Q(s)$. In fact, it is straightforward to check that the transfer function from $y$ to $u$ of the controller depicted in fig. 3.6 is precisely $K(s)$. In view of this fact it is particularly meaningful to look for a realization of the
subsystem $R_{f}(s)$, since the generic stabilizing controller must result from the feedback connection of it with a stable system $Q(s)$ which is only constrained to actually produce a well defined overall system (condition (3.12)).

With reference to fig. 3.6, the transfer function of system $R_{f}(s)$ from the two inputs $y$ and $y_{1}$ to the two outputs $u$ and $u_{1}$ is

$$
R_{f}(s)=\left[\begin{array}{cc}
K_{0}(s) & -\hat{X}^{-1}(s) \\
X^{-1}(s) & X^{-1}(s) N(s)
\end{array}\right]
$$

Preliminarily a realization is presented for each one of the transfer functions appearing in $R_{f}(s)$. The proof of Theorem 3.4 (as for the form of $\left.N(s), X(s), Y(s), \hat{X}(s)\right)$ and Section 2.5 (as for the form of the inverse system) are expedient to such an operation. Letting $\bar{C}:=C+D F, \bar{B}:=B+H D$ and $\bar{A}:=A+B F+H \bar{C}$, one obtains

$$
X^{-1}(s) N(s)=\left[\begin{array}{cc|c}
A+B F & 0 & B \\
H \bar{C} & \bar{A} & H D \\
\hline \bar{C} & \bar{C} & D
\end{array}\right]
$$

It is easy to ascertain that a lower order realization can be found by performing a change of variables which put into evidence an unobservable part and is defined by the matrix

$$
T_{1}:=\left[\begin{array}{cc}
I & I \\
0 & I
\end{array}\right]
$$

namely,

$$
X^{-1}(s) N(s)=\left[\begin{array}{c|c}
\bar{A} & \bar{B} \\
\hline \bar{C} & D
\end{array}\right]
$$

In a similar way one gets

$$
K_{0}(s)=Y(s) X^{-1}(s)=\left[\begin{array}{cc|c}
\bar{A} & 0 & H \\
H \bar{C} & A+B F & H \\
\hline 0 & -F & 0
\end{array}\right]
$$

simplified to

$$
K_{0}(s)=\left[\begin{array}{r|r}
\bar{A} & H \\
\hline-F & 0
\end{array}\right]
$$

through the change of variables defined by the matrix

$$
T_{2}:=\left[\begin{array}{rr}
-I & I \\
0 & I
\end{array}\right]
$$

which put into evidence an unreachable part. Finally, one has

$$
X^{-1}(s)=\left[\begin{array}{c|c}
\bar{A} & H \\
\hline \bar{C} & I
\end{array}\right], \quad \hat{X}^{-1}(s)=\left[\begin{array}{c|c}
\bar{A} & \bar{B} \\
\hline F & I
\end{array}\right]
$$

Thus a realization for system $R_{f}(s)$ is

$$
R_{f}(s)=\left[\begin{array}{r|rr}
\bar{A} & H & \bar{B} \\
\hline-F & 0 & -I \\
\bar{C} & I & D
\end{array}\right]=\left[\begin{array}{r|rr}
\bar{A} & -H & -\bar{B} \\
\hline F & 0 & -I \\
-\bar{C} & I & D
\end{array}\right]
$$



Figure 3.7: A realization of $R_{f}(s)$

This particular realization corresponds to the block-scheme depicted in fig. 3.7. This scheme clearly shows that when $y_{1}=0$, that is when $Q(s)=0$ or, in other words, when the central controller has been adopted, the stabilizing controller generates the control variable $u$ through the law $u=F \xi, \xi$ being the state of an observer with dynamic matrix $A+H C$. Thus the central controller coincides with the controller designed via the well known pole placement technique and any stabilizing controller for $G(s)$ can be said to be "built around" a state observer.

Finally, the eigenvalues of the resulting control system, that is of the system in fig. 3.5 with $K(s)$ given by the block-scheme of fig. 3.6 where $R_{f}(s)$ is specified by the block-diagram of fig. 3.7, are those of matrices $A+B F$ and $A+H C$ together with those of the dynamic matrix of $Q(s):=\Sigma\left(A_{q}, B_{q}, C_{q}, D_{q}\right)$. In fact, denoting by $x$ and $x_{q}$ the state variables of $G(s)$ and $Q(s)$, respectively, the choice $x_{t}:=\left[x^{\prime}-\xi^{\prime} x_{q}^{\prime} x^{\prime}\right]^{\prime}$ as state variables of resulting system, yields the (closed loop) dynamic matrix

$$
A_{t}=\left[\begin{array}{ccc}
A+H C & 0 & 0 \\
B_{q} C & A_{q} & 0 \\
-B\left(F+D_{q} C\right) & -B C_{q} & A+B F
\end{array}\right]
$$

### 3.5 Notes and references

The line of reasoning pursued in this chapter and the derivation of some results presented herein are inspired by the books of Francis [19] and Maciejowski [43]. The result on the parametrization of stabilizing controllers can be found in the pioneering paper by Youla et al. [63].

## Chapter 4

## $\mathrm{RH}_{2}$ Control

### 4.1 Introduction

This chapter presents the most significant results concerning the control problem in the $R H_{2}$ context: it simply consists in minimizing the $R H_{2}$ norm of a transfer function. As it will be apparent in the sequel, a number of connections can be established between this problem (more precisely, between the subproblems which actually constitute its frame) and the most celebrated set of results in optimal filtering and control problems, the well known Linear Quadratic Gaussian theory. Approaching these problems within the $\mathrm{RH}_{2}$ setting gives a somehow more complete picture of the structure of the results (see the forthcoming Remarks 4.3, 4.10, 4.19).

Throughout this chapter reference will be made to a controlled system described by the following equations

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{4.1}\\
z & =C_{1} x+D_{11} w+D_{12} u  \tag{4.2}\\
y & =C_{2} x+D_{21} w+D_{22} u \tag{4.3}
\end{align*}
$$

The measured output variable $y$ is the input to the controller, which is constrained to be a finite dimensional, linear, time invariant system, while the control variable $u$ is its output. Therefore, the controller takes on the form

$$
\begin{align*}
\dot{\xi} & =F \xi+G y  \tag{4.4}\\
u & =H \xi+E y \tag{4.5}
\end{align*}
$$



Figure 4.1: The standard control system

The situation is summarized in the block-scheme of fig. 4.1, where $P(s)$ and $K(s)$ denote the transfer functions of the controlled system and the controller, respectively. More precisely,

$$
\begin{gathered}
P(s):=\left[\begin{array}{c|cc}
A & B_{1} & B_{2} \\
\hline C_{1} & D_{11} & D_{12} \\
C_{2} & D_{21} & D_{22}
\end{array}\right] \\
K(s):=\left[\begin{array}{c|c}
F & G \\
\hline H & E
\end{array}\right]
\end{gathered}
$$

It is apparent that the feedback connection of system (4.1)-(4.3) with system (4.4),(4.5) is well defined only if the controller is such that

$$
\begin{equation*}
\operatorname{det}\left[I-E D_{22}\right] \neq 0 \tag{4.6}
\end{equation*}
$$

that is only if the algebraic loop deriving from its implementation can be solved.
The structure of system (4.1)-(4.3) is typical of the control problems stated within the $R H_{2}$ and $R H_{\infty}$ context. Indeed, it is sufficiently general so as to encompass the most meaningful situations. A fairly comprehensive discussion on this aspect is presented at the beginning of Chapter 5 to adequately motivate the significance of the results there presented for the $R H_{\infty}$ control problems.

The input variable $w$ collects all exogenous signals which can be viewed as disturbances acting on the control system, while the output performance variable $z$ is expedient to specify an index of the performances to be attained by the control system.

With reference to the control system depicted in fig. 4.1, the main objective is minimizing (with respect to $K(s)$ ) the $R H_{2}$ norm of the transfer function $T(z, w ; s)$ from the input $w$ to the output $z$.

If $T(z, w ; s)$ has to belong to $R H_{2}$, then it must necessarily be strictly proper and stable. The first requirement is equivalent to the condition

$$
\begin{equation*}
D_{11}+D_{12}\left(I-E D_{22}\right)^{-1} E D_{21}=0 \tag{4.7}
\end{equation*}
$$

As for the second requirement, it is no doubt satisfied if the internal stability of the control system is ensured, that is if the controller (4.4),(4.5) internally stabilizes the controlled system (4.1)-(4.3), which simply amounts to asking for the stability of the dynamic matrix of the resulting system, namely

$$
\begin{equation*}
\operatorname{Re}\left(\lambda_{i}\left(A_{F}\right)\right)<0, \forall i \tag{4.8}
\end{equation*}
$$

where

$$
A_{F}:=\left[\begin{array}{cc}
A+B_{2}\left(I-E D_{22}\right)^{-1} E C_{2} & B_{2}\left(I-E D_{22}\right)^{-1} H  \tag{4.9}\\
G\left[I+D_{22}\left(I-E D_{22}\right)^{-1} E\right] C_{2} & F+G D_{22}\left(I-E D_{22}\right)^{-1} H
\end{array}\right]
$$

In view of the discussion above the notion of admissible controller is introduced in the following definition.

Definition 4.1 (Admissible controller in $R H_{2}$ ) A controller $K(s)$ is said to be admissible in $\mathrm{RH}_{2}$ for $P(s)$ if conditions (4.6)-(4.9) are verified.

Three problems will be discussed in the forthcoming Sections 4.2-4.4: (i) the full information problem; (ii) the output estimation problem; (iii) the partial information problem. Each of these problems refers to a particular structure of the controlled system $P(s)$ and exhibits deep and significant mutual connections. Indeed, the last one presents an interesting separation property and will be dealt with by exploiting the solutions relevant to the first two, which, in turn, are characterized by strong duality properties.

The main result concerning each one of the above problems provides the answer to three important questions: (i) the existence of the optimal controller; (ii) the actual form of such a controller; (iii) the parametrization of suboptimal controllers. The material will be presented according to the scheme formally stated in Problem 4.1 which refers to the feedback connected system shown in fig. 4.1 and to the set $\mathcal{F}_{2 \gamma}$ of controllers which are admissible in $R H_{2}$ for $P(s)$ and such that $\|T(z, w ; s)\|_{2}$ is bounded by a given positive scalar $\gamma$.

## Problem 4.1 (Standard problem in $\mathrm{RH}_{2}$ ) Find

a) The minimum value (if any) of $\|T(z, w ; s)\|_{2}$ attainable by a controller $K(s)$ which is admissible in $R H_{2}$ for $P(s)$.
b) An admissible controller which minimizes $\|T(z, w ; s)\|_{2}$.
c) A set of controllers $\mathcal{F}_{2 \gamma r} \subseteq \mathcal{F}_{2 \gamma}$ whose elements generate the whole set of functions $T(z, w ; s)$ which are generated by the elements of $\mathcal{F}_{2 \gamma}$.

Remark 4.1 An obvious necessary condition for the existence of a stabilizing controller (and therefore for the existence of an admissible controller in $R H_{2}$ for $P(s)$ ) is the stabilizability of the pair $\left(A, B_{2}\right)$ and the detectability of the pair $\left(A, C_{2}\right)$. The statement of Problem 4.1 makes sense only if both properties actually hold true.

In the forthcoming sections the parametrization of the controllers in the family $\mathcal{F}_{2 \gamma}$ will be presented in specific remarks which follow the main theorems concerning the solution of Problem 4.1. Hence, the issues relative to the family $\mathcal{F}_{2 \gamma}$ and those relative to Problem 4.1 will be treated separately.

Problem 4.1 will be tackled in the subsequent sections by assuming $D_{11}=0$ in eq. (4.2). Actually, this assumption does not cause any loss of generality. In fact, assume that the set of matrices $E$ which satisfy eqs. (4.6),(4.7) is not empty (otherwise the problem would not admit any solution in $R H_{2}$ ) and let $\hat{E}$ be one of such matrices. It is straightforward to check that the situation $D_{11}=0$ is recovered after the output feedback

$$
u=\hat{E} y+v
$$

has been implemented.

### 4.2 The full information problem

In this section Problem 4.1 is considered by assuming that the output signal $y$ of system (4.1)-(4.3) is constituted by the state vector $x$ and the disturbance vector $w$. Moreover, no direct influence of the input signals $w$ and $u$ exists on the outputs $z$ and $y$, respectively. As a consequence, the controlled plant, whose transfer function from


Figure 4.2: The full information problem
$\left[\begin{array}{ll}w^{\prime} & u^{\prime}\end{array}\right]^{\prime}$ to $\left[\begin{array}{ll}z^{\prime} & y^{\prime}\end{array}\right]^{\prime}$ will be denoted by $P_{F I}(s)$, is described by

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{4.10}\\
z & =C_{1} x+D_{12} u  \tag{4.11}\\
y & =\left[y_{1}^{\prime} y_{2}^{\prime}\right]^{\prime}  \tag{4.12}\\
y_{1} & =x  \tag{4.13}\\
y_{2} & =w \tag{4.14}
\end{align*}
$$

Furthermore, the following assumptions will be made.
Assumption 4.1 The pair $\left(A, B_{2}\right)$ is stabilizable and no eigenvalue of the unobservable part of the pair $\left[\left(A-B_{2} D_{12}^{\prime} C_{1}\right),\left(I-D_{12} D_{12}^{\prime}\right) C_{1}\right]$ lies on the imaginary axis.

Assumption 4.2 $D_{12}^{\prime} D_{12}=I$.
In the theorem below reference is made to the block-scheme of fig. 4.2 where $K_{F I}(s)$ is a generic controller admissible in $R H_{2}$ for $P_{F I}(s)$. In this diagram $T(z, w ; s)$ is the transfer function from $w$ to $z$.

Theorem 4.1 (Full information) Consider Problem 4.1 relative to system (4.10)(4.14). Then, under Assumptions 4.1, 4.2, it has the solution
a)

$$
\min \|T(z, w ; s)\|_{2}=\left\|P_{c}(s) B_{1}\right\|_{2}=\sqrt{\operatorname{trace}\left[B_{1}^{\prime} P_{2} B_{1}\right]}
$$

b)

$$
K_{F I}^{o}(s)=\left[\begin{array}{c|cc}
\emptyset & \emptyset & \emptyset \\
\hline \emptyset & F_{2} & 0
\end{array}\right]
$$

c) The set $\mathcal{F}_{2 \gamma r}$ of the controllers $K_{\text {FIr }}(s)$ is defined by the diagram of fig. 4.3, where $Q(s):=\Sigma\left(A_{q}, B_{q}, C_{q}, 0\right)$ with $A_{q}$ stable and $\|Q(s)\|_{2}^{2}<\gamma^{2}-\left\|P_{c}(s) B_{1}\right\|_{2}^{2}$.

In the three points above, $\gamma$ is a positive scalar such that $\gamma>\left\|P_{c}(s) B_{1}\right\|_{2}$ and

$$
\begin{gather*}
F_{2}:=-B_{2}^{\prime} P_{2}-D_{12}^{\prime} C_{1}  \tag{4.15}\\
P_{c}(s):=\left[\begin{array}{c|c}
A_{c}-B_{2} B_{2}^{\prime} P_{2} & I \\
\hline C_{1 c}-D_{12} B_{2}^{\prime} P_{2} & 0
\end{array}\right]  \tag{4.16}\\
A_{c}:=A-B_{2} D_{12}^{\prime} C_{1}, C_{1 c}:=\left(I-D_{12} D_{12}^{\prime}\right) C_{1} \tag{4.17}
\end{gather*}
$$



Figure 4.3: The set $\mathcal{F}_{2 \gamma r}$ of the controllers $K_{\text {FIr }}(s)$
where $P_{2}$ is the symmetric, positive semidefinite and stabilizing solution of the Riccati equation (in the unknown $P$ )

$$
\begin{equation*}
0=P A_{c}+A_{c}^{\prime} P-P B_{2} B_{2}^{\prime} P+C_{1 c}^{\prime} C_{1 c} \tag{4.18}
\end{equation*}
$$

i.e. such that matrix $A_{c c}$ defined by

$$
\begin{equation*}
A_{c c}:=A_{c}-B_{2} B_{2}^{\prime} P_{2}=A+B_{2} F_{2} \tag{4.19}
\end{equation*}
$$

is stable.
Proof First observe that the necessary condition for the problem at hand to make sense (recall Remark 4.1) is satisfied. Indeed, being measurable the state of the system, detectability of the pair $\left(A, C_{2}\right)$ trivially holds, while, on the other hand, stabilizability of the pair $\left(A, B_{2}\right)$ is guaranteed by Assumption 4.1.

Points a) and b) Notice that the pair $\left(A, B_{2}\right)$ is stabilizable if and only if the pair $\left(A_{c}, B_{2}\right)$ is such (recall that state feedback does not modify the stabilizability property). Therefore, Assumptions 4.1 and 4.2 together with Lemma C. 4 guarantee the existence of the symmetric, positive semidefinite and stabilizing solution $P_{2}$ of eq. (4.18), so that the matrix $A_{c c}$ defined by eqs. (4.19),(4.15) is stable.

Furthermore, let

$$
\begin{equation*}
v:=u-F_{2} x \tag{4.20}
\end{equation*}
$$

Equation (4.20) apparently defines the control law

$$
\begin{equation*}
u:=v+F_{2} x \tag{4.21}
\end{equation*}
$$

From eqs. (4.10),(4.11),(4.21) it follows

$$
\begin{equation*}
z_{L 0}=P_{c}(s) B_{1} w_{L}+U(s) v_{L} \tag{4.22}
\end{equation*}
$$

where

$$
\begin{equation*}
U(s):=\Sigma\left(A_{c}-B_{2} B_{2}^{\prime} P_{2}, B_{2}, C_{1 c}-D_{12} B_{2}^{\prime} P_{2}, D_{12}\right) \tag{4.23}
\end{equation*}
$$

The variable $v$ defined by eq. (4.20) is one of the output of the system

$$
P_{v}(s):=\left\{\begin{array}{l}
\dot{x}=A x+B_{1} w+B_{2} u \\
v=-F_{2} x+u \\
y=\left[y_{1}^{\prime} y_{2}^{\prime}\right]^{\prime}, \quad y_{1}=x, \quad y_{2}=w
\end{array}\right.
$$



Figure 4.4: The equivalent full information problem

In order to evaluate the transfer function $T(z, w ; s)$ (relevant to the scheme in fig. 4.2 ) by means of eq. (4.22), it is convenient to make reference to the block-diagram in fig. 4.4 and to the related transfer function $T(v, w ; s)$ from the input $w$ to the output $v$. From eqs. (4.10)-(4.14) and the definition of $P_{v}(s)$ it follows that the systems in fig. 4.2 and 4.4 are well defined for all $K_{F I}(s)$, since condition (4.6) is satisfied in both cases. Therefore, it is possible to write

$$
\begin{equation*}
T(z, w ; s)=P_{c}(s) B_{1}+U(s) T(v, w ; s) \tag{4.24}
\end{equation*}
$$

From eqs. (4.10)-(4.14) and the definition of $P_{v}(s)$ it also follows that the system in fig. 4.4 is stable if and only if the system in fig. 4.2 is such. Moreover, letting $K_{F I}(\infty):=\left[E_{1} E_{2}\right]$, it is $T(z, w ; \infty)=D_{12} E_{2}$ and $T(v, w ; \infty)=E_{2}$. Thus, $D_{12} E_{2}=0$ if and only if $E_{2}=0$ since, thanks to Assumption 4.2, the matrix $D_{12}$ has rank equal to the number of its columns. Therefore, the function $T(z, w ; s)$ is strictly proper if and only if the function $T(v, w ; s)$ is such. In conclusion, it can be claimed that the controller $K_{F I}(s)$ is admissible in $R H_{2}$ for $P_{F I}(s)$ if and only if it is such also for $P_{v}(s)$.

From Assumption 4.2 it follows that $\left(I-D_{12} D_{12}^{\prime}\right) D_{12}=0$, so that $C_{1 c}^{\prime} D_{12}=0$. This fact, together with Lemma C. 5 implies that the function $U(s)$ defined by eq. (4.23) is inner and $U^{\sim}(s) P_{c}(s) \in R H_{2}^{\perp}$.

In view of eq. (4.24), it follows that

$$
\begin{align*}
\|T(z, w ; s)\|_{2}^{2}= & \left\|P_{c}(s) B_{1}\right\|_{2}^{2}+\|U(s) T(v, w ; s)\|_{2}^{2}+ \\
& +2<U(s) T(v, w ; s), P_{c}(s) B_{1}> \tag{4.25}
\end{align*}
$$

for any controller $K_{F I}(s)$ admissible in $R H_{2}$ for $P_{F I}(s)$. Being $U(s)$ an inner function, one gets

$$
\begin{aligned}
\|U(s) T(v, w ; s)\|_{2}^{2} & =<U(s) T(v, w ; s), U(s) T(v, w ; s)> \\
& =<T(v, w ; s), T(v, w ; s)>=\|T(v, w ; s)\|_{2}^{2}
\end{aligned}
$$

while, being $U^{\sim}(s) P_{c}(s) \in R H_{2}^{\perp}$, it is

$$
<U(s) T(v, w ; s), P_{c}(s) B_{1}>=<T(v, w ; s), U^{\sim}(s) P_{c}(s) B_{1}>=0
$$

since $T(v, w ; s) \in R H_{2}$ (recall Theorem 2.11). Then, eq. (4.25) becomes

$$
\begin{equation*}
\|T(z, w ; s)\|_{2}^{2}=\left\|P_{c}(s) B_{1}\right\|_{2}^{2}+\|T(v, w ; s)\|_{2}^{2} \tag{4.26}
\end{equation*}
$$

so that

$$
\min _{K_{F I}(s)}\|T(z, w ; s)\|_{2}^{2}=\left\|P_{c}(s) B_{1}\right\|_{2}^{2}+\min _{K_{F I}(s)}\|T(v, w ; s)\|_{2}^{2}
$$

The minimum is by sure attained if $\|T(v, w ; s)\|_{2}^{2}=0$, that is if $v=0$. This is the case if

$$
\begin{equation*}
u=F_{2} x \tag{4.27}
\end{equation*}
$$

Now observe that $P_{2}$, being a solution of eq. (4.18), is also a positive semidefinite solution of the Lyapunov equation (in the unknown $P$ )

$$
0=P A_{c c}+A_{c c}^{\prime} P+P_{2} B_{2} B_{2}^{\prime} P_{2}+C_{1 c}^{\prime} C_{1 c}
$$

where $A_{c c}$ is defined by eq. (4.19). Matrix $A_{c c}$ is the dynamic matrix of the system (4.10)-(4.14),(4.27), namely

$$
\begin{align*}
& \dot{x}=A_{c c} x+B_{1} w  \tag{4.28}\\
& z=\left(C_{1 c}-D_{12} B_{2}^{\prime} P_{2}\right) x \tag{4.29}
\end{align*}
$$

The norm of the (optimal) transfer function of system (4.28), (4.29) may be computed by exploiting what has been presented in Remark 2.20. To this aim, notice that from Assumption 4.2 and $C_{1 c}^{\prime} D_{12}=0$ it follows

$$
\left(C_{1 c}-D_{12} B_{2}^{\prime} P_{2}\right)^{\prime}\left(C_{1 c}-D_{12} B_{2}^{\prime} P_{2}\right)=P_{2} B_{2} B_{2}^{\prime} P_{2}+C_{1 c}^{\prime} C_{1 c}
$$

and recall that $P_{2}$ solves the above Lyapunov equation. Therefore, $\|T(z, w ; s)\|_{2}^{2}=$ trace $\left[B_{1}^{\prime} P_{2} B_{1}\right]$. Points $\left.a\right)$ and $b$ ) are thus proved.

Point c) Let a generic controller $K_{F I}(s) \in \mathcal{F}_{2 \gamma}$ be described by the equations

$$
\begin{aligned}
\dot{\xi} & =L \xi+M_{1} x+M_{2} w \\
u & =N \xi+O_{1} x+O_{2} w
\end{aligned}
$$

If $v$ is given by eq. (4.20), then the same arguments exploited for proving points $a$ ) and $b$ ) lead to eq. (4.26), so that

$$
\|T(v, w ; s)\|_{2}^{2}<\gamma^{2}-\left\|P_{c}(s) B_{1}\right\|_{2}^{2}
$$

Recall that if $K_{F I}(s)$ is admissible in $R H_{2}$ for $P_{F I}(s)$, then it is also admissible in $R H_{2}$ for $P_{v}(s)$. Therefore, $T(v, w ; s)$ can be written as

$$
T(v, w ; s):=Q(s):=\Sigma\left(A_{q}, B_{q}, C_{q}, 0\right)
$$

with $A_{q}$ stable. A realization of $Q(s)$ can be easily derived by recalling the definition of $P_{v}(s)$ and the above given expression for $K_{F I}(s)$. It results

$$
\begin{aligned}
\dot{\vartheta} & =L \vartheta+M_{1} \sigma+M_{2} w \\
\dot{\sigma} & =B_{2} N \vartheta+\left(A+B_{2} O_{1}\right) \sigma+\left(B_{1}+B_{2} O_{2}\right) w \\
v & =N \vartheta+\left(O_{1}-F_{2}\right) \sigma+O_{2} w
\end{aligned}
$$

The controller $K_{\text {FIr }}(s)$ defined by these equations and $u=F_{2} x+v$ is now shown to belong to the set $\mathcal{F}_{2 \gamma r}$. Obviously, it possesses the structure of the system in fig. 4.3. Moreover, thanks to eqs. (4.10), (4.11), the equations of the system resulting from the feedback connection of $P_{F I}(s)$ with $K_{F I r}(s)$ are, letting $\varepsilon:=\sigma-x$,

$$
\begin{aligned}
\dot{x} & =A x+B_{1} w+B_{2} u \\
\dot{\vartheta} & =L \vartheta+M_{1} x+M_{2} w+M_{1} \varepsilon \\
\dot{\varepsilon} & =\left(A+B_{2} F_{2}\right) \varepsilon \\
u & =N \vartheta+O_{1} x+\left(O_{1}-F_{2}\right) \varepsilon+O_{2} w \\
z & =C_{1} x+D_{12} u
\end{aligned}
$$

The comparison of these equations with those relevant to the feedback connection of system $P_{F I}(s)$ with the controller $K_{F I}(s)$ allows one to derive the following two conclusions: first, the transfer functions from $w$ to $z$ are equal in the two cases; second, the system having $K_{\text {FIr }}(s)$ as a controller is stable because the system having $K_{F I}(s)$ as a controller is such and matrix $A+B_{2} F_{2}$ is stable.

Vice versa, if $K_{F I r}(s)$ belongs to the set described in fig. 4.3, then it is $u=F_{2} x+v$ where $v$ is the output of the system $T(v, w ; s)=\Sigma\left(A_{q}, B_{q}, C_{q}, 0\right)$, with $A_{q}$ stable and $\|T(v, w ; s)\|_{2}^{2}<\gamma^{2}-\left\|P_{c}(s) B_{1}\right\|_{2}^{2}$. The same arguments exploited in proving points $a$ ) and $b$ ) imply that $K_{F I r}(s)$ is admissible in $R H_{2}$ for $P_{F I}(s)$ and, thanks to eq. (4.26), $\|T(z, w ; s)\|_{2}^{2}<\gamma^{2}$.

Example 4.1 Consider system (4.10)-(4.14) with

$$
\begin{gathered}
A=\left[\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right], \quad B_{1}=\left[\begin{array}{l}
0 \\
2
\end{array}\right], \quad B_{2}=\left[\begin{array}{l}
0 \\
1
\end{array}\right] \\
C_{1}=\left[\begin{array}{ll}
\alpha & 1
\end{array}\right], \quad D_{12}=1
\end{gathered}
$$

where $\alpha \neq 0$ and let $u=u_{w}+u_{x}$. The classic way of designing a controller which reduces the effect of the disturbance $w$ on the output $z$ is trying to zeroing the transfer function from $w$ to $z$ by means of a suitable controller $K_{D C}(s)$ with input $w$ and output $u_{w}$ which performs a direct compensation. In the problem at hand such a transfer function vanishes (namely, $z$ does not depend on $w$ ) if

$$
K_{D C}(s)=-\frac{2(s+\alpha)}{s^{2}+s+\alpha}
$$

However, it is apparent that the resulting control system can not be stabilized (by means of a control law $u=u_{w}+u_{x}$ which makes $u_{x}$ to depend on $x$ in a suitable manner) if $K_{D C}(s)$ is not stable by itself. Hence, if $\alpha>0$, a perfect direct compensation of the disturbance can be achieved. On the contrary, if $\alpha<0$, this is no more possible. By resorting to Theorem 4.1 one obtains, correspondingly to $Q(s)=0, u=-\alpha x_{1}-x_{2}$ and $T(z, w ; s)=0$, when $\alpha>0$, while, when $\alpha<0$, one gets $u=\alpha x_{1}-\sqrt{1-4 \alpha} x_{2}$ and

$$
T(z, w ; s)=2 \frac{2 \alpha+(1-\sqrt{1-4 \alpha}) s}{s^{2}+\sqrt{1-4 \alpha} s-\alpha}
$$

Remark 4.2 The structure of the controllers $K_{F I r}(s)$ which are admissible in $R H_{2}$ for $P_{F I}(s)$ and defined by the block-scheme in fig. 4.3 allows one to easily conclude that the eigenvalues of the resulting control system are those of matrix $A+B_{2} F_{2}$ and those of matrix $A_{q}$.

Remark 4.3 (Parametrization of the set $\mathcal{F}_{2 \gamma}$ ) Observe that the set $\mathcal{F}_{2 \gamma r}$ is a proper subset of the set $\mathcal{F}_{2 \gamma}$. In fact, consider a generic controller $K_{F I}(s)$ admissible in $R H_{2}$ for $P_{F I}(s)$ which is purely algebraic and makes the control variables to depend on the state variables only. Namely, let

$$
\begin{equation*}
u=\Lambda x \tag{4.30}
\end{equation*}
$$

with $\Lambda \neq F_{2}$. Such a controller certainly exists because of a continuity argument, since the controller $K_{F I}^{o}(s)$ actually has the form given in eq. (4.30).

Within the set $\mathcal{F}_{2 \gamma r}$ the only element with the form given in eq. (4.30) is $K_{F I}^{o}(s)$, so that it is necessary to resort to a dynamic controller in the set $\mathcal{F}_{2 \gamma r}$ in order to generate the
same transfer function $T(z, w ; s)$. As an example, for the system $P_{F I}(s)$ given by

$$
\begin{aligned}
& \dot{x}=w+u \\
& z=\left[\begin{array}{l}
1 \\
0
\end{array}\right] x+\left[\begin{array}{l}
0 \\
1
\end{array}\right] u \\
& y=\left[\begin{array}{l}
x \\
w
\end{array}\right]
\end{aligned}
$$

the controller (admissible in $R H_{2}$ for such a $P_{F I}(s)$ ) described by

$$
\begin{equation*}
u=-3 x \tag{4.31}
\end{equation*}
$$

has a "corresponding element" in the set $\mathcal{F}_{2 \gamma_{r}}$ given by (recall what has been shown in the proof of Theorem 4.1, point $c$ ))

$$
\begin{align*}
\dot{\sigma} & =-3 \sigma+w  \tag{4.32}\\
u & =-x-2 \sigma \tag{4.33}
\end{align*}
$$

Indeed, by exploiting eq. (4.31), it follows

$$
T(z, w ; s)=\frac{1}{s+3}\left[\begin{array}{r}
1 \\
-3
\end{array}\right]
$$

which coincides with the expression deriving from eqs. (4.32),(4.33).
A parametrization of the set $\mathcal{F}_{2 \gamma}$ is now presented. Consider the system $P_{F}(s)$ which is obtained from system (4.10)-(4.14) after the control law (4.21) has been implemented, namely the system

$$
\begin{aligned}
P_{F}(s): & =\left[\begin{array}{cc}
P_{F 11}(s) & P_{F 12}(s) \\
P_{F 21}(s) & P_{F 22}(s)
\end{array}\right] \\
& =\left[\begin{array}{c|cc}
A_{c c} & B_{1} & B_{2} \\
\hline C_{1 F} & 0 & D_{12} \\
{\left[\begin{array}{c}
I \\
0
\end{array}\right]} & {\left[\begin{array}{cc}
0 & 0 \\
I & 0
\end{array}\right]}
\end{array}\right]
\end{aligned}
$$

where eqs. (4.19) and (4.15) have been taken into account and

$$
C_{1 F}:=C_{1 c}-D_{12} B_{2}^{\prime} P_{2}
$$

The set of controllers $K_{F}(s)$ which stabilize $P_{F}(s)$ apparently coincides with the set of controllers which stabilize $P_{F 22}(s)$, the latter being the system

$$
P_{F 22}(s)=\left[\begin{array}{c|c}
A_{c c} & B_{2} \\
\hline I & 0 \\
0 & 0
\end{array}\right]
$$

consistently to what has been previously defined. In view of Theorem 3.7, such a set can be expressed as

$$
K_{F}(s)=[Y(s)-M(s) \Theta(s)][X(s)-N(s) \Theta(s)]^{-1}
$$

with $\Theta(s) \in R H_{\infty}$. Notice that condition (3.12) is no doubt satisfied in the present context, since $P_{F 22}(s)$ is strictly proper. From the proof of Theorem 3.4 it follows that, being the matrix $A_{c c}$ stable, the choice

$$
Y(s)=0, \quad M(s)=I, \quad X(s)=I, \quad N(s)=P_{F 22}(s)
$$

is admissible, so that

$$
\begin{equation*}
K_{F}(s)=-\Theta(s)\left[I-P_{F 22}(s) \Theta(s)\right]^{-1} \tag{4.34}
\end{equation*}
$$

In the closed loop situation it is

$$
y_{L 0}=P_{F 21}(s) w_{L}-P_{F 22}(s) \Theta(s)\left[I-P_{F 22}(s) \Theta(s)\right]^{-1} y_{L 0}
$$

By exploiting Lemma B.9, this equation can be written as

$$
y_{L 0}=\left[I-P_{F 22}(s) \Theta(s)\right] P_{F 21}(s) w_{L}
$$

which, in turn, can be utilized to obtain

$$
\begin{aligned}
z_{L 0} & =P_{F 11}(s) w_{L}-P_{F 12}(s) \Theta(s)\left[I-P_{F 22}(s) \Theta(s)\right]^{-1} y_{L 0} \\
& =\left[P_{F 11}(s)-P_{F 12}(s) \Theta(s) P_{F 21}(s)\right] w_{L}
\end{aligned}
$$

so that, when the controller described by

$$
u_{L 0}=F_{2} x_{L 0}+K_{F}(s) y_{L 0}
$$

is adopted, the closed loop transfer function from $w$ to $z$ is

$$
T_{\Theta}(z, w ; s)=P_{F 11}(s)-P_{F 12}(s) \Theta(s) P_{F 21}(s)
$$

On the other hand, if a controller $K_{F I r}(s)$ of the set $\mathcal{F}_{2 \gamma r}$ is adopted, namely, a controller described by

$$
u_{L 0}=F_{2} x_{L 0}+Q(s) w_{L}
$$

with (recall point (c) in the statement of Theorem 4.1) $Q(s) \in R H_{2}$, and $\|Q(s)\|_{2}^{2}<\gamma^{2}-$ $\left\|P_{c}(s) B_{1}\right\|_{2}^{2}=\gamma^{2}-\left\|P_{F 11}(s)\right\|_{2}^{2}$, one gets

$$
z_{L 0}=\left[P_{F 11}(s)+P_{F 12}(s) Q(s)\right] w_{L}
$$

Therefore, the set of the $K_{F I r}(s)$ generates the set of transfer functions from $w$ to $z$

$$
T_{Q}(z, w ; s)=P_{F 11}(s)+P_{F 12}(s) Q(s)
$$

By equating the transfer functions $T_{Q}(z, w ; s)$ and $T_{\ominus}(z, w ; s)$ it is possible to characterize the set of functions $\Theta(s)$ associated with controllers $K_{F}(s)$ admissible in $R H_{2}$ for $P_{F}(s)$ and such that $\left\|T_{\Theta}(z, w ; s)\right\|_{2}<\gamma$. In so doing one obtains

$$
P_{F 12}(s)\left[\Theta(s) P_{F 21}(s)+Q(s)\right]=0
$$

Letting $\Phi(s):=\left(s I-A_{c c}\right)^{-1}$, from Assumption 4.2 it follows that the rank of $P_{F 12}(s)=$ $C_{1 F} \Phi(s) B_{2}+D_{12}$ equals the number of its columns, so that the above written equation is equivalent to

$$
\begin{equation*}
\Theta(s) P_{F 21}(s)=-Q(s) \tag{4.35}
\end{equation*}
$$

A particular solution of this equation is

$$
\bar{\Theta}(s)=Q(s)[0-I]
$$

since

$$
P_{F 21}(s)=\left[\begin{array}{c}
\Phi(s) B_{1} \\
I
\end{array}\right]
$$

Thus, the general solution of eq. (4.35) is

$$
\Theta_{Q}(s)=\bar{\Theta}(s)+\hat{\Theta}(s)
$$

where $\hat{\Theta}(s)$ denotes any solution in $R H_{\infty}$ of the homogeneous equation

$$
\Theta(s) P_{F 21}(s)=0
$$



Figure 4.5: The generic admissible controller for $P_{F I}(s)$

Letting $\Theta(s):=\left[\Theta_{1}(s) \Theta_{2}(s)\right]$, from the last equation it follows that $\Theta_{2}(s)=-\Theta_{1}(s) \Phi(s) B_{1}$, so that

$$
\hat{\Theta}(s)=\Theta_{1}(s)\left[I-\Phi(s) B_{1}\right]
$$

belongs to $R H_{\infty}$ if and only if $\Theta_{1}(s) \in R H_{\infty}$, since $\Phi(s) \in R H_{\infty}$, being $A_{c c}$ a stable matrix.
Therefore, the set of functions $\Theta_{Q}(s)$ which generate controllers $K_{F}(s)$ admissible in $R H_{2}$ for $P_{F}(s)$ and such that $\left\|T_{\Theta}(z, w ; s)\right\|_{2}<\gamma$ is defined by

$$
\begin{gathered}
\Theta_{Q}(s)=\left[\Theta_{1}(s)-\left[Q(s)+\Theta_{1}(s) \Phi(s) B_{1}\right]\right] \\
\Theta_{1}(s) \in R H_{\infty}, Q(s) \in R H_{2},\|Q(s)\|_{2}^{2}<\gamma^{2}-\left\|P_{F 11}(s)\right\|_{2}^{2}
\end{gathered}
$$

By exploiting Lemma B.9, from eq. (4.34) it follows

$$
\begin{align*}
K_{F}(s)= & -\left[I-\Theta_{Q}(s) P_{F 22}(s)\right]^{-1} \Theta_{Q}(s) \\
= & {\left[I-\Theta_{1}(s) \Phi(s) B_{2}\right]^{-1} } \\
& \cdot\left[-\Theta_{1}(s)\left[Q(s)+\Theta_{1}(s) \Phi(s) B_{1}\right]\right] \tag{4.36}
\end{align*}
$$

since $\Theta_{Q}(s) P_{F 22}(s)=\Theta_{1}(s) \Phi(s) B_{2}$. In conclusion, by recalling eq. (4.21), the generic controller $K_{F I}(s)$ in the set $\mathcal{F}_{2 \gamma}$ is described by

$$
\begin{aligned}
& K_{F I}(s)= {\left[I-\Theta_{1}(s) \Phi(s) B_{2}\right]^{-1} } \\
& \cdot\left[-\Theta_{1}(s)\left[Q(s)+\Theta_{1}(s) \Phi(s) B_{1}\right]\right]+\left[F_{2} 0\right] \\
& \Theta_{1}(s) \in R H_{\infty}, Q(s) \in R H_{2},\|Q(s)\|_{2}^{2}<\gamma^{2}-\left\|P_{F 11}(s)\right\|_{2}^{2}
\end{aligned}
$$

Such a set is depicted in fig. 4.5. In the controller shown in this figure it is

$$
\begin{aligned}
v_{L 0} & =\Theta_{1}(s) \Phi(s)\left[-\Phi^{-1}(s) x_{L}+B_{2} v_{L 0}+B_{1} w_{L}\right]+Q(s) w_{L} \\
& =-\Theta_{1}(s) \Phi(s) x(0)+Q(s) w_{L}
\end{aligned}
$$

having taken into account eqs. (4.10) and (4.21). Therefore, the effect of the parameter $\Theta_{1}(s)$ (which is responsible of the difference between the controllers in $\mathcal{F}_{2 \gamma}$ and those in $\mathcal{F}_{2 \gamma r}$ ) amounts to a term which only depends on the initial state of system (4.10)-(4.14).

With reference to the example previously considered it is easy to verify that the algebraic controller characterized by $u=-3 x$ corresponds to the choice

$$
Q(s)=-\frac{2}{s+3}, \quad \Theta_{1}(s)=2 \frac{s+\mathbf{1}}{s+3}
$$

Remark 4.4 (Optimal control problems) The control problem to which Theorem 4.1 can be applied is strictly related to the linear quadratic deterministic ( $L Q$ ) or stochastic (LQS) control problem with measurable state. This connection can be made explicit by exploiting what has been presented in Remark 2.20.
$\underline{\text { LQ Problem Consider the } n \text {-th order system }}$

$$
\begin{gather*}
\dot{x}=A x+B \bar{u}  \tag{4.37}\\
x(0)=x_{0} \tag{4.38}
\end{gather*}
$$

and the cost functional

$$
\begin{equation*}
J_{1}:=\int_{0}^{\infty}\left[x^{\prime}(t) Q x(t)+2 x^{\prime}(t) S \bar{u}(t)+\bar{u}^{\prime}(t) R \bar{u}(t)\right] d t \tag{4.39}
\end{equation*}
$$

where

$$
\left[\begin{array}{ll}
Q & S  \tag{4.40}\\
S^{\prime} & R
\end{array}\right]:=L=L^{\prime} \geq 0, \quad R>0
$$

Observe that $L \geq 0$ and $R>0$ imply that

$$
\hat{Q}:=Q-S R^{-1} S^{\prime} \geq 0
$$

since $\hat{Q}=Z^{\prime} L Z$ with $Z^{\prime}:=\left[I-S R^{-1}\right]$. Let $C_{11} \in R^{n \times n}$ be a factorization of $\hat{Q}$, so that

$$
\begin{equation*}
C_{11}^{\prime} C_{11}=\hat{Q} \tag{4.41}
\end{equation*}
$$

and define

$$
\begin{gather*}
C_{1}:=\left[\begin{array}{c}
C_{11} \\
R^{-1 / 2} S^{\prime}
\end{array}\right], \quad D_{12}:=\left[\begin{array}{c}
0 \\
I
\end{array}\right]  \tag{4.42}\\
u:=R^{1 / 2} \bar{u}, z:=C_{1} x+D_{12} u \tag{4.43}
\end{gather*}
$$

It is easy to verify that it is

$$
J_{1}=\int_{0}^{\infty} z^{\prime}(t) z(t) d t
$$

On the other hand, the state free motion can always be interpreted as the forced motion caused by an impulsive input acting on the system through a suitable input matrix. Therefore, if eqs. (4.41)-(4.43) are taken into account, system (4.37),(4.38) can be described by

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{4.44}\\
z & =C_{1} x+D_{12} u \tag{4.45}
\end{align*}
$$

with $w:=\delta(t), B_{1}:=x_{0}, B_{2}:=B R^{-1 / 2}$ and $x(0)=0$.
The optimal control problem at hand ( $L Q$ problem) consists in finding a controller of the form

$$
\begin{align*}
& \dot{\xi}=F \xi+G x  \tag{4.46}\\
& u=H \xi+N x \tag{4.47}
\end{align*}
$$

such that the system (4.44)-(4.47) is stable and the performance index $J_{1}$ is minimized. Notice that the feedback connection of any controller of the form (4.46),(4.47) to the system (4.44),(4.45) is well defined and the relevant transfer function $T(z, w ; s)$ from $w$ to $z$ is strictly proper. Therefore, if system (4.44)-(4.47) is stable, then the controller (4.46),(4.47) is admissible in $\mathrm{RH}_{2}$ for system (4.44),(4.45). Further, in view of Remark 2.20, any controller (4.46),(4.47) which is admissible in $R H_{2}$ for system (4.44), (4.45) is such that $J_{1}=$ $\|T(z, w ; s)\|_{2}^{2}$.

If the pair $(A, B)$ is stabilizable and no eigenvalue of the unobservable part of the pair $\left(A-B R^{-1} S^{\prime}, C_{11}\right)$ lies on the imaginary axis, Assumptions 4.1 and 4.2 are satisfied. In fact, from one hand, Assumption 4.2 is readily ascertained to hold in view of eqs. (4.41)-(4.43). On the other hand, by performing the required substitutions, it is still straightforward to check that the unobservable part of the pair $\left[\left(A-B_{2} D_{12}^{\prime} C_{1}\right),\left(I-D_{12} D_{12}^{\prime}\right) C_{1}\right]$ coincides with the unobservable part of the pair $\left[\left(A-B_{2} R^{-1 / 2} S^{\prime}\right), C_{11}\right]$. Finally, observe that the pair $\left(A, B_{2}\right)$ is stabilizable if and only if the pair $(A, B)$ is such (recall the definition of $\left.B_{2}\right)$.

If stated in terms of the system (4.44),(4.45) with the additional output equation $y=$ $\left[x^{\prime} w^{\prime}\right]^{\prime}$, the control problem addressed to by Theorem 4.1 is solved by a controller of the form (4.46),(4.47). Therefore, the controller defined at point (b) of Theorem 4.1 constitutes the solution of the $L Q$ control problem, too.

However, observe that within the context of the classical optimal control theory the linear quadratic problem is stated without requiring the stability of the resulting control system. Thus the assumptions which are necessary to guarantee the existence of the solution in that context (stability of the observable and unreachable part of system $\Sigma\left(A-B R^{-1} S^{\prime}, B, C_{11}, 0\right)$ ) are weaker than those required within the $R H_{2}$ context. As an example, consider system (4.37) and the performance index (4.39) with

$$
A=\left[\begin{array}{cc}
\alpha & 1 \\
0 & 0
\end{array}\right], \quad B=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad Q=\left[\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right], \quad R=1, \quad S=0
$$

where $\alpha=1$ or $\alpha=-1$. Within the framework of the classical optimal control theory the solution of this problem is given by the control law $u_{C L}^{o}(x)=-\left[\begin{array}{ll}0 & 1\end{array}\right] x$ to which there corresponds the value $J_{1 C L}^{o}(x(0))=x_{2}^{2}(0)$. The resulting control system is stable when $\alpha=-1$ and unstable when $\alpha=1$. In the $\mathrm{RH}_{2}$ context the optimal control law is $u_{R H_{2}}^{o}(x)=-\left[\begin{array}{ll}4 & 3\end{array}\right] x$ when $\alpha=1$. Correspondingly, the value of the performance index is $J_{1 R H_{2}}^{o}(x(0))=8 x_{1}^{2}(0)+8 x_{1}(0) x_{2}(0)+3 x_{2}^{2}(0)$. When $\alpha=-1$ the same control law as in the classical setting is found. In the $\mathrm{RH}_{2}$ framework the resulting control system is stable in both cases. Finally, notice that $J_{1 C L}^{o}(x(0)) \leq J_{1 R H_{2}}^{o}(x(0)), \forall x(0)$.

LQS Problem Assume that the controlled system is

$$
\dot{x}=A x+B_{1} w+B \bar{u}
$$

where $w$ is a zero mean white noise with identity intensity. Let the pair $(A, B)$ be stabilizable and consider either the cost functional

$$
J_{2}:=\lim _{t \rightarrow \infty} E\left[x^{\prime}(t) Q x(t)+2 x^{\prime}(t) S \bar{u}(t)+\bar{u}^{\prime}(t) R \bar{u}(t)\right]
$$

or the cost functional

$$
J_{3}:=\lim _{T \rightarrow \infty} E\left[\frac{1}{T} \int_{0}^{T}\left[x^{\prime}(t) Q x(t)+2 x^{\prime}(t) S \bar{u}(t)+\bar{u}^{\prime}(t) R \bar{u}(t)\right] d t\right]
$$

where matrices $Q, R, S$ satisfy eq. (4.40). From eqs. (4.41)-(4.43) it follows

$$
\begin{gathered}
J_{2}=\lim _{t \rightarrow \infty} E\left[z^{\prime}(t) z(t)\right] \\
J_{3}=\lim _{T \rightarrow \infty} E\left[\frac{1}{T} \int_{0}^{T} z^{\prime}(t) z(t) d t\right]
\end{gathered}
$$

and the controlled system is described by

$$
\begin{align*}
& \dot{x}=A x+B_{1} w+B_{2} u  \tag{4.48}\\
& z=C_{1} x+D_{12} u \tag{4.49}
\end{align*}
$$



Figure 4.6: The disturbance feedforward problem


Figure 4.7: The auxiliary full information problem

The optimal control problem at hand ( $L Q S$ problem) consists in finding a controller of the form given in eqs. (4.46),(4.47) such that stability of the resulting system is ensured and $J_{2}$ or $J_{3}$ is minimized when it is connected to system (4.48),(4.49). The same arguments exploited for the $L Q$ problem lead to $J_{2}=J_{3}=\|T(z, w ; s)\|_{2}^{2}$. If system (4.48),(4.49) verifies the same assumptions as system (4.44),(4.45), then the solution of the $L Q S$ problem is again the one specified under point (b) of Theorem 4.1.

Remark 4.5 (Disturbance feedforward) Here reference is made to the block-scheme of fig. 4.6, where $P_{D F}(s)$ is described by the equations

$$
\begin{align*}
& \dot{x}=A x+B_{1} w+B_{2} u  \tag{4.50}\\
& z=C_{1} x+D_{12} u  \tag{4.51}\\
& y=C_{2} x+w \tag{4.52}
\end{align*}
$$

while $K_{D F}(s)$ is a generic controller admissible in $R H_{2}$ for $P_{D F}(s)$.
It is assumed that the pair $\left(A, B_{2}\right)$ is stabilizable and no eigenvalue of the unobservable part of the pair $\left[\left(A-B_{2} D_{12}^{\prime} C_{1}\right),\left(I-D_{12} D_{12}^{\prime}\right) C_{1}\right]$ lies on the imaginary axis. Moreover the matrix $A-B_{1} C_{2}$ is supposed to be stable and $D_{12}^{\prime} D_{12}=I$. Observe that stability of $A-B_{1} C_{2}$ implies detectability of the pair ( $A, C_{2}$ ) which, together with the stabilizability of the pair ( $A, B_{2}$ ), guarantees the fulfillment of the necessary condition reported in Remark 4.1.

Now, notice that system $P_{D F}(s)$ is equal to system $\hat{P}(s)$ defined in Lemma E.2. Moreover, let

$$
\bar{P}_{F I}(s):=\left[\begin{array}{c|cc}
A & B_{1} & B_{2} \\
\hline C_{1} & 0 & D_{12} \\
I & 0 & 0 \\
0 & I & 0
\end{array}\right]
$$

Observe that: ( $i$ ) System $\bar{P}_{F I}(s)$ is equal to system $\bar{P}(s)$ defined in Lemma E.2; (ii) System $\bar{P}_{F I}(s)$ is equal to system $P_{F I}(s)$ relative to which the full control problem has been stated.

Lemma E. 2 (which can be exploited, thanks to stability of matrix $A-B_{1} C_{2}$ ) ensures that the controller $\bar{K}_{F I}(s)$ connected to system $\bar{P}_{F I}(s)$ as shown in fig. 4.7 , stabilizes system


Figure 4.8: The structure of the controller $\hat{K}(s)$ in terms of $\bar{K}_{F I}(s)$
$\bar{P}_{F I}(s)$ if and only if the controller $\hat{K}(s)$, defined in the block-scheme of fig. 4.8, where

$$
K_{v}(s):=\left[\begin{array}{c|cc}
A-B_{1} C_{2} & B_{1} & B_{2} \\
\hline 0 & 0 & I \\
I & 0 & 0 \\
-C_{2} & I & 0
\end{array}\right]
$$

stabilizes $\hat{P}(s)$ to which it is connected according to the scheme shown in fig. 4.9, this latter scheme being identical to the one depicted in fig. 4.6. Moreover, the transfer function from $\bar{w}$ to $\bar{z}$ in fig. 4.7 and the transfer function from $\hat{w}$ to $\hat{z}$ in fig. 4.9 are equal. Therefore, the solution of Problem 4.1 relative to system $P_{D F}(s)$ (fig. 4.6) can be found by solving the same problem relative to system $\bar{P}_{F I}(s)$ (fig. 4.7) via Theorem 4.1. Notice that such a theorem can be exploited because the assumptions made for system $P_{D F}(s)$ imply the fulfillment of Assumptions 4.1 and 4.2. Thus
a) $\min \|T(z, w ; s)\|_{2}=\left\|P_{c}(s) B_{1}\right\|_{2}=\sqrt{\operatorname{trace}\left[B_{1}^{\prime} P_{2} B_{1}\right]}$;
b) The optimal controller is given by

$$
K_{D F}^{o}(s)=\left[\begin{array}{c|c}
A-B_{1} C_{2}+B_{2} F_{2} & B_{1} \\
\hline F_{2} & 0
\end{array}\right]
$$

c) The set $\mathcal{F}_{2 \gamma r}$ of the controllers $K_{D F r}(s)$ is defined by the block-scheme of fig. 4.10 where

$$
N_{2}(s):=\left[\begin{array}{c|cc}
A-B_{1} C_{2}+B_{2} F_{2} & B_{1} & B_{2} \\
\hline F_{2} & 0 & I \\
-C_{2} & I & 0
\end{array}\right]
$$

$Q(s):=\Sigma\left(A_{q}, B_{q}, C_{q}, 0\right),\|Q(s)\|_{2}^{2}<\gamma^{2}-\left\|P_{c}(s) B_{1}\right\|_{2}^{2}$ and $A_{q}$ is a stable matrix.
In the three points above $\gamma$ is a positive scalar such that $\gamma>\left\|P_{c}(s) B_{1}\right\|_{2}$ and reference has been made to eqs. (4.15)-(4.20). The problem at hand is referred to as the disturbance feedforward problem in view of the following discussion.

Preliminarily, observe that, if in eq. (4.52) $C_{2}=0$, then $y=w$ so that the disturbance $w$ can be measured (direct compensation). However, in such a case the controller $K_{D F}(s)$ is connected to the controlled system $P_{D F}(s)$ in an open rather than closed loop configuration (see also fig. 4.6), so that the stability assumption of the matrix $A-B_{1} C_{2}$ simply reduces to the stability assumption of the controlled system, which is obviously necessary to ensure the stability of the resulting control system.

If, on the contrary, $C_{2} \neq 0$, it is still possible to get the disturbance $w$ in terms of the variables $u$ and $y$. In fact, from eqs. (4.50),(4.52) it follows

$$
\begin{equation*}
y_{L 0}=P_{21}(s) w_{L}+P_{22}(s) u_{L} \tag{4.53}
\end{equation*}
$$



Figure 4.9: The equivalent disturbance feedforward problem


Figure 4.10: The set $\mathcal{F}_{2 \gamma r}$ of the controllers $K_{D F r}(s)$
where

$$
P_{21}(s):=\left[\begin{array}{c|c}
A & B_{1} \\
\hline C_{2} & I
\end{array}\right], \quad P_{22}:=\left[\begin{array}{c|c}
A & B_{2} \\
\hline C_{2} & 0
\end{array}\right]
$$

System $P_{21}(s)$ is invertible (recall what has been presented in Section 2.5) so that from eq. (4.53) one gets

$$
w_{L}=P_{21}^{-1}(s)\left[y_{L 0}-P_{22}(s) u_{L}\right]
$$

with

$$
P_{21}^{-1}(s)=\left[\begin{array}{c|c}
A-B_{1} C_{2} & B_{1} \\
\hline-C_{2} & I
\end{array}\right]
$$

Now define a precompensator $K_{F I D F}(s)$ by means of the equations

$$
\begin{align*}
\dot{\xi} & =\left(A-B_{1} C_{2}\right) \xi+B_{2} u+B_{1} y  \tag{4.54}\\
\hat{w} & =-C_{2} \xi+y \tag{4.55}
\end{align*}
$$

Comparing eqs. (4.50),(4.52) with eq. (4.54) leads to the conclusion that, letting $\varepsilon:=\xi-x$,

$$
\begin{equation*}
\dot{\varepsilon}=\left(A-B_{1} C_{2}\right) \varepsilon \tag{4.56}
\end{equation*}
$$

Thus, $x_{L 0}=\xi_{L 0}$ and, from eqs. (4.52),(4.55), also $w_{L}=\hat{w}_{L 0}$, so that, also when $C_{2} \neq 0$, it is still possible to think to $w$ as being measurable by resorting to a suitable dynamical system (4.54),(4.55) (indirect compensation).

Therefore, one can consider Problem 4.1 relative to system (4.50)-(4.52), rather than to system (4.50),(4.51) with the additional equation

$$
\hat{y}:=\left[\begin{array}{l}
x  \tag{4.57}\\
w
\end{array}\right]
$$

The new problem has the very structure of a full information problem. Thus, the controller $K_{F I}(s)$ which solves Problem 4.1 relative to system (4.50),(4.51),(4.57) must be connected


Figure 4.11: $K_{D F}(s)$ in terms of the precompensator $K_{F I D F}(s)$
to system (4.54),(4.55) in order to obtain the controller $K_{D F}(s)$ (see also fig. 4.11). The resulting system includes an unreachable part (recall eq. (4.56)) which is nevertheless stable, thanks to the stated assumptions. Moreover, it is easy to check that whenever the controller solving Problem 4.1 relative to system (4.50),(4.51),(4.57) belongs to the set of controllers $K_{F I r}(s)$, then the controller which is obtained by substituting $K_{F I}(s)$ with $K_{F I r}(s)$ in fig. 4.11 belongs to the set of controllers $K_{D F r}(s)$. In other words, the procedure just now presented for solving Problem 4.1 leads to the results previously found.

Finally, notice that $\mathcal{F}_{2 \gamma r}=\mathcal{F}_{2 \gamma}$. In fact, by recalling what has been said at the end of Remark 4.3, resorting to the controller $K_{F I}(s)$ rather than to the (simpler) controller $K_{F I r}(s)$ implies adding the term $\Theta_{1}(s) \Phi(s) \xi(0)$ to the control variable $u_{L 0}$. Such a term depends on the initial state of the precompensator $K_{F I D F}(s)$ only, so that the transfer function from $y$ to $u$ is independent of the parameter $\Theta_{1}(s)$. Alternatively, from eqs. (4.21),(4.54),(4.55) one gets

$$
\xi_{L 0}=\Phi(s)\left(B_{2} v_{L 0}+B_{1} \hat{w}_{L 0}\right)
$$

where $\Phi(s):=\left(s I-A-B_{2} F_{2}\right)^{-1}$. Hence, by recalling eq. (4.36),

$$
\begin{aligned}
v_{L 0}= & {\left[I-\Theta_{1}(s) \Phi(s) B_{2}\right]^{-1}\left\{-\Theta_{1}(s) \Phi(s)\left(B_{2} v_{L 0}+B_{1} \hat{w}_{L 0}\right)+\right.} \\
& \left.+\left[Q(s)+\Theta_{1}(s) \Phi(s) B_{1}\right] \hat{w}_{L 0}\right\} \\
= & {\left[I-\Theta_{1}(s) \Phi(s) B_{2}\right]^{-1}\left[Q(s) \hat{w}_{L 0}-\Theta_{1}(s) \Phi(s) B_{2} v_{L 0}\right] }
\end{aligned}
$$

Therefore, it follows

$$
v_{L 0}=Q(s) \hat{w}_{L 0}
$$

so that, in view of eq. (4.21),

$$
u_{L 0}=F_{2} \xi_{L 0}+Q(s) \hat{w}_{L 0}
$$

which, thanks to eqs. (4.54),(4.55), implies that the transfer function of the generic controller $K_{D F}(s)$ is independent of $\Theta_{1}(s)$.

Example 4.2 Consider system (4.50)-(4.52) with

$$
\begin{aligned}
& A=\left[\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right], \quad B_{1}=\left[\begin{array}{l}
0 \\
2
\end{array}\right], \quad B_{2}=\left[\begin{array}{l}
0 \\
1
\end{array}\right] \\
& C_{1}=\left[\begin{array}{ll}
1 & \alpha
\end{array}\right], \quad C_{2}=\left[\begin{array}{ll}
1 & 1
\end{array}\right], \quad D_{12}=1
\end{aligned}
$$

where $\alpha \neq 0$. The classical synthesis procedure for a controller which reduces the influence of the disturbance $w$ on the output $z$ consist in trying to make zero the transfer function from $w$ to $z$ by introducing a controller $K_{I C}(s)$ with input $y$ and output $u$ (indirect compensation). In the problem at hand such a transfer function is zero (hence $z$ does not depend on $w$ ) if

$$
K_{I C}(s)=-\frac{2(\alpha s+1)}{s^{2}+(\alpha+2) s+3}
$$

However, it is easy to verify that the resulting system is stable when $\alpha>0$, while it is not stable when $\alpha<0$. Thus, in this second case it is no more possible to perform a (perfect) indirect compensation of the disturbance $w$. On the contrary, by resorting to Remark 4.5, one gets $K_{D F}^{o}(s)=K_{I C}(s)$ when $\alpha>0$ and

$$
K_{D F}^{o}(s)=-2 \frac{1-\alpha s}{s^{2}+(2-\alpha) s+3}
$$

when $\alpha<0$.
Remark 4.6 Assumption 4.2 is restrictive though standard in control problems. It implies that matrix $D_{12}$ has rank equal to the number of its columns. Whenever such an assumption is not verified one has to tackle the so called singular problem the solution of which (if it exists) calls for a theoretical development far beyond the scope of this book.

On the contrary, notice that the condition $D_{12}^{\prime} D_{12}=I$ can be replaced by $D_{12}^{\prime} D_{12}=R$, with $R>0$. Indeed, the here presented case can be easily recovered by a suitable redefinition of the control variable, according to eq. (4.43).

Finally, it is often set $D_{12}^{\prime} C_{1}=0$. This simplifying orthogonality assumption, besides greatly reducing the notational burden, corresponds to the absence of the cross term $x^{\prime}(t) S \bar{u}(t)$ in the performance index of the optimal control problems dealt with in Remark 4.4.

Remark 4.7 The Riccati equation (4.18) coincides with the one encountered within the (classical) Optimal Regulator problem defined in terms of the system

$$
\dot{x}=A_{c} x+B_{2} u
$$

and the performance index

$$
J:=\int_{0}^{\infty}\left[x^{\prime} C_{1 c}^{\prime} C_{1 c} x+u^{\prime} u\right] d t
$$

Remark 4.8 Under Assumption 4.2, the fact that the no eigenvalue of the unobservable part of the pair [ $\left.\left(A-B_{2} D_{12}^{\prime} C_{1}\right),\left(I-D_{12} D_{12}^{\prime}\right) C_{1}\right]$ lies on the imaginary axis (Assumption 4.1) is equivalent to the subsystem of $P_{F I}(s)$ corresponding to the transfer function $P_{F I 12}(s)$ from the input $u$ and the output $z$ (namely system $\Sigma\left(A, B_{2}, C_{1}, D_{12}\right)$ ) not to have invariant zeros with $\operatorname{Re}(s)=0$. In fact, in view of the material in Section 2.5 and recalling that the dimension of $u$ is not greater than the dimension of $z$, if it exists $\lambda$ with $\operatorname{Re}(\lambda)=0$ such that

$$
\begin{gathered}
(\lambda I-A) x-B_{2} u=0 \\
C_{1} x+D_{12} u=0
\end{gathered}
$$

with $\left[x^{\prime} u^{\prime}\right]^{\prime} \neq 0$, then from the identity $D_{12}^{\prime} D_{12}=I$ it follows

$$
u=-D_{12}^{\prime} C_{1} x
$$

and hence also

$$
\begin{gathered}
\left(A-B_{2} D_{12}^{\prime} C_{1}\right) x=\lambda x \\
\left(I-D_{12} D_{12}^{\prime}\right) C_{1} x=0
\end{gathered}
$$

Being $x \neq 0$, since, otherwise, also $u=0$, this violates the assumption that no eigenvalue of the unobservable part of the pair $\left[\left(A-B_{2} D_{12}^{\prime} C_{1}\right),\left(I-D_{12} D_{12}^{\prime}\right) C_{1}\right]$ lies on the imaginary axis (recall Lemma D.1, part (a)). By going the other way on, if the last two equations hold corresponding to $x \neq 0$, the proof of the above mentioned equivalence can be carried out by letting $u:=-D_{12}^{\prime} C_{1} x$.


Figure 4.12: The output estimation problem

### 4.3 The output estimation problem

The problem of optimally observing a linear combination of the state variables is faced in this section by following an approach which relies on $\mathrm{RH}_{2}$ techniques.

Consider the system

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{4.58}\\
z & =C_{1} x+u  \tag{4.59}\\
y & =C_{2} x+D_{21} w \tag{4.60}
\end{align*}
$$

and let $P_{O E}(s)$ denote its transfer function. With reference to the block-scheme in fig. 4.12, suppose that $K_{O E}(s)$ is a $R H_{2}$ admissible controller for $P_{O E}(s)$ which makes "small" $\|T(z, w ; s)\|_{2}$ (the norm of the transfer function from $w$ to $z$ ). Then, the variable $u$ provides a "good" estimate of the linear combination $-C_{1} x$ of the state variables. Indeed, should $T(z, w ; s)=0$ (namely, $z_{L 0}=0$ ) then it would follows $u_{L 0}=-C_{1} x_{L 0}$. In such a case the signal $u$ apparently constitutes the best possible estimate of $-C_{1} x$ (when $\mathrm{x}(0)=0$ ). The following assumptions are now introduced.

Assumption 4.3 The pair $\left(A, C_{2}\right)$ is detectable and no eigenvalue of the unreachable part of the pair $\left[\left(A-B_{1} D_{21}^{\prime} C_{2}\right), B_{1}\left(I-D_{21}^{\prime} D_{21}\right)\right]$ lies on the imaginary axis.

Assumption 4.4 $D_{21} D_{21}^{\prime}=I$.
Assumption 4.5 $A-B_{2} C_{1}$ is stable.
Under these assumptions it is possible to state the following theorem.
Theorem 4.2 (Output estimation) Consider Problem 4.1 relative to system (4.58)(4.60). Then, under Assumptions 4.3-4.5, it has the solution
a)

$$
\min \|T(z, w ; s)\|_{2}=\left\|C_{1} P_{f}(s)\right\|_{2}=\sqrt{\operatorname{trace}\left[C_{1} \Pi_{2} C_{1}^{\prime}\right]}
$$

b)

$$
K_{O E}^{o}(s)=\left[\begin{array}{c|c}
A_{f}-B_{2} C_{1}-\Pi_{2} C_{2}^{\prime} C_{2} & L_{2} \\
\hline C_{1} & 0
\end{array}\right]
$$

c) The set $\mathcal{F}_{2 \gamma r}$ of controllers $K_{O E r}(s)$ is defined by the block-scheme in fig. 4.13, where

$$
M_{2}(s):=\left[\begin{array}{c|cc}
A_{f}-B_{2} C_{1}-\Pi_{2} C_{2}^{\prime} C_{2} & L_{2} & -B_{2} \\
\hline C_{1} & 0 & I \\
C_{2} & I & 0
\end{array}\right]
$$



Figure 4.13: The set $\mathcal{F}_{2 \gamma r}$ of the controllers $K_{O E r}(s)$
with $Q(s):=\Sigma\left(A_{q}, B_{q}, C_{q}, 0\right),\|Q(s)\|_{2}^{2}<\gamma^{2}-\left\|C_{1} P_{f}(s)\right\|_{2}^{2}$ and $A_{q}$ stable.
In the three preceding points $\gamma$ is a positive scalar such that $\gamma>\left\|C_{1} P_{f}(s)\right\|_{2}$, while

$$
\begin{gather*}
L_{2}:=-\Pi_{2} C_{2}^{\prime}-B_{1} D_{21}^{\prime}  \tag{4.61}\\
P_{f}(s):=\left[\begin{array}{c|c}
A_{f}-\Pi_{2} C_{2}^{\prime} C_{2} & B_{1 f}-\Pi_{2} C_{2}^{\prime} D_{21} \\
\hline I & 0
\end{array}\right]  \tag{4.62}\\
A_{f}:=A-B_{1} D_{21}^{\prime} C_{2}, \quad B_{1 f}:=B_{1}\left(I-D_{21}^{\prime} D_{21}\right) \tag{4.63}
\end{gather*}
$$

where $\Pi_{2}$ is the symmetric, positive definite and stabilizing solution of the Riccati equation (in the unknown $\Pi$ )

$$
\begin{equation*}
0=\Pi A_{f}^{\prime}+A_{f} \Pi-\Pi C_{2}^{\prime} C_{2} \Pi+B_{1 f} B_{1 f}^{\prime} \tag{4.64}
\end{equation*}
$$

that is such that matrix $A_{f c}$ defined by

$$
\begin{equation*}
A_{f c}:=A_{f}-\Pi_{2} C_{2}^{\prime} C_{2}=A+L_{2} C_{2} \tag{4.65}
\end{equation*}
$$

is stable.
Proof Preliminarily, notice that Assumption 4.5 implies that the pair $\left(A, B_{2}\right)$ is stabilizable. This fact, together with the assumed detectability of the pair $\left(A, C_{2}\right)$, makes the necessary condition in Remark 4.1 satisfied.

Now consider system $\hat{P}_{O E}(s)$ obtained from system $P_{O E}(s)$ by transposition. From eqs. (4.58)-(4.60) it follows that it is described by

$$
\begin{aligned}
\dot{\xi} & =F \xi+G_{1} \zeta+G_{2} \eta \\
\omega & =H_{1} \xi+E \eta \\
\varphi & =H_{2} \xi+\zeta
\end{aligned}
$$

where

$$
\begin{equation*}
F:=A^{\prime}, \quad E:=D_{21}^{\prime} \tag{4.66}
\end{equation*}
$$

and, for $i=1,2$,

$$
\begin{equation*}
G_{i}:=C_{i}^{\prime}, \quad H_{i}:=B_{i}^{\prime} \tag{4.67}
\end{equation*}
$$

Therefore, system $\hat{P}_{O E}(s)$ possesses the very same structure of system $P_{D F}(s)$ considered in Remark 4.5. Thus, Problem 4.1 relative to system $\hat{P}_{O E}(s)$ is solved as it is mentioned in such a remark (relatively to $P_{O E}(s)$ ). Indeed, the assumption required for system $P_{O E}(s)$ are precisely Assumptions 4.3-4.5. Thanks to Lemma E.1, the results concerning system $P_{O E}(s)$ can be derived by transposition of those relevant to system $\hat{P}_{O E}(s)$, provided that eqs. (4.66),(4.67) are taken into account.

Example 4.3 Consider system (4.58)-(4.60) with

$$
\begin{aligned}
A & =\left[\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right], \quad B_{1}=\left[\begin{array}{ll}
1 & 0 \\
1 & 0
\end{array}\right], \quad B_{2}=\left[\begin{array}{l}
0 \\
1
\end{array}\right] \\
C_{1} & =\left[\begin{array}{ll}
1 & 1
\end{array}\right], \quad C_{2}=\left[\begin{array}{ll}
1 & 0
\end{array}\right], \quad D_{21}=\left[\begin{array}{ll}
0 & 1
\end{array}\right]
\end{aligned}
$$

One obtains $L_{2}=[-\sqrt{3}-1]^{\prime}$ and $\left\|C_{1} P_{f}(s)\right\|_{2}^{2}=4.46$. Taken $\gamma^{2}=9$ and

$$
Q(s)=\left[\begin{array}{rr|r}
-1 & 0 & 1 \\
0 & -2 & 1 \\
\hline 1 & 1 & 0
\end{array}\right]
$$

it is $\|T(z, w ; s)\|_{2}^{2}=5.88<\gamma^{2}$, consistently with $1.42=\|Q(s)\|_{2}^{2}<\gamma^{2}-4.46$.
Remark 4.9 The structure of a generic controller $K_{O E}(s)$ admissible in $\mathrm{RH}_{2}$ for $P_{O E}(s)$ as defined by the block-scheme of fig. 4.13 allows one to easily verify that the eigenvalues of the resulting control system are those of matrices $A+L_{2} C_{2}, A-B_{2} C_{1}$ and $A_{q}$. In fact, letting $x_{m}$ and $x_{q}$ be the state variables of systems $M_{2}(s)$ and $Q(s)$, respectively, and choosing the state vector of the resulting control system as $x_{t}:=\left[x^{\prime}+x_{m}^{\prime} x_{q}^{\prime} x^{\prime}\right]^{\prime}$, its dynamic matrix is

$$
A_{t}=\left[\begin{array}{ccc}
A+L_{2} C_{2} & 0 & 0 \\
B_{q} C_{2} & A_{q} & 0 \\
B_{2} C_{1} & B_{2} C_{q} & A-B_{2} C_{1}
\end{array}\right]
$$

the eigenvalues of which are precisely those above mentioned.
Observe that the order of the resulting system is $2 n+n_{q}$, where $n$ is the order of system (4.58)-(4.60) and $n_{q}$ is the order of system $Q(s)$. Moreover, when $Q(s)=0$ it results $T(z, w ; s)=C_{1}\left(s I-A_{f c}\right)^{-1}\left(B_{1}+L_{2} D_{21}\right)$, so that the transfer function from $w$ to $z$ does not depend on $B_{2}$ (recall that, in view of eqs. (4.61),(4.63),(4.64), $L_{2}$ is independent of $B_{2}$ ).

Remark 4.10 (Parametrization of the set $\mathcal{F}_{2 \gamma}$ ) Notice that $\mathcal{F}_{2 \gamma}=\mathcal{F}_{2 \gamma r}$. In fact, as shown in the proof of Theorem 4.2, the $R H_{2}$ admissible controllers for $P_{O E}(s)$ can be obtained by transposing those which are admissible for $\hat{P}_{O E}(s)$, this last system possessing the structure of system $P_{D F}(s)$ considered in Remark 4.5. Having proved that $\mathcal{F}_{2 \gamma}=\mathcal{F}_{2 \gamma r}$ relative to $P_{D F}(s)$, the same conclusion must hold for $P_{O E}(s)$.

Remark 4.11 (Optimal filtering) The control problem relative to system $P_{O E}(s)$ can be interpreted as an optimal filtering problem relative to the $n$-th order system

$$
\begin{align*}
\dot{x} & =A x+\zeta_{1}+B_{2} u  \tag{4.68}\\
\bar{y} & =C x+\zeta_{2} \tag{4.69}
\end{align*}
$$

where $\zeta:=\left[\zeta_{1}^{\prime} \zeta_{2}^{\prime}\right]^{\prime}$ is a zero mean, Gaussian white noise with intensity

$$
W:=\left[\begin{array}{ll}
W_{11} & W_{12}  \tag{4.70}\\
W_{12}^{\prime} & W_{22}
\end{array}\right], \quad W_{22}>0
$$

If an asymptotic estimate of a generic linear combination $S x$ of the system state has to be found, then one of the two functionals

$$
J_{4}:=\lim _{t \rightarrow \infty} E\left[[S x(t)-u(t)]^{\prime}[S x(t)-u(t)]\right]
$$

and

$$
J_{5}:=\lim _{T \rightarrow \infty} E\left[\frac{1}{T} \int_{0}^{T}[S x(t)-u(t)]^{\prime}[S x(t)-u(t)] d t\right]
$$

may conveniently be associated with the system. Notice that $\hat{W}:=W_{11}-W_{12} W_{22}^{-1} W_{12}^{\prime} \geq 0$, since $W \geq 0$ and $W_{22}>0$. Indeed, it is $\hat{W}=Z^{\prime} W Z$, with

$$
Z:=\left[\begin{array}{c}
I \\
-W_{22}^{-1} W_{12}^{\prime}
\end{array}\right]
$$

Let the $n \times n$ matrix $B_{11}$ be a factorization of $\hat{W}$, so that

$$
\begin{equation*}
B_{11} B_{11}^{\prime}=\hat{W} \tag{4.71}
\end{equation*}
$$

and define

$$
\begin{align*}
D_{21} & :=\left[\begin{array}{ll}
0 & I
\end{array}\right], \quad B_{1}:=\left[\begin{array}{ll}
B_{11} & W_{12} W_{22}^{-1 / 2}
\end{array}\right]  \tag{4.72}\\
y & :=W_{22}^{-1 / 2} \bar{y}, \quad C_{2}:=W_{22}^{-1 / 2} C  \tag{4.73}\\
z & :=C_{1} x+u, \quad C_{1}:=-S \tag{4.74}
\end{align*}
$$

Then system (4.68)-(4.70) can be rewritten as

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{4.75}\\
z & =C_{1} x+u  \tag{4.76}\\
y & =C_{2} x+D_{21} w \tag{4.77}
\end{align*}
$$

where $w$ is the zero mean Gaussian white noise with identity intensity which satisfies the equation

$$
\left[\begin{array}{l}
\zeta_{1}  \tag{4.78}\\
\zeta_{2}
\end{array}\right]=\left[\begin{array}{cc}
B_{11} & W_{12} W_{22}^{-1 / 2} \\
0 & W_{22}^{1 / 2}
\end{array}\right] w
$$

Observe that

$$
\begin{equation*}
B_{1} B_{1}^{\prime}=W_{11}, \quad B_{1} D_{21}^{\prime}=W_{12} W_{22}^{-1 / 2}, \quad D_{21} D_{21}^{\prime}=I \tag{4.79}
\end{equation*}
$$

It is easy to verify that

$$
J_{4}=\lim _{t \rightarrow \infty} E\left[z^{\prime}(t) z(t)\right]
$$

and

$$
J_{5}=\lim _{T \rightarrow \infty} E\left[\frac{1}{T} \int_{0}^{T} z^{\prime}(t) z(t) d t\right]
$$

The problem at hand consists in finding a controller of the form

$$
\begin{align*}
\dot{\xi} & =F \xi+G y  \tag{4.80}\\
u & =H \xi+N y \tag{4.81}
\end{align*}
$$

such that the control system (4.75)-(4.81) is stable and either the criterion $J_{4}$ or $J_{5}$ is minimized. Observe that the feedback connection of any controller of the above form with system (4.75)-(4.79) is always well defined, while the transfer function $T(z, w ; s)$ from $w$ to $z$ relevant to such a connection is strictly proper if and only if $N D_{21}=0$, that is, in view of eq. (4.72), if and only if $N=0$. Therefore, if system (4.75)-(4.81) is stable and $N=0$, then the controller (4.80), (4.81) is $R H_{2}$ admissible for system (4.75)-(4.79). Further, from Remark 2.20, any $\mathrm{RH}_{2}$ admissible controller for system (4.75)-(4.79) is such that

$$
J_{4}=J_{5}=\|T(z, w ; s)\|_{2}^{2}
$$

If the pair $(A, C)$ is detectable and no eigenvalue of the unreachable part of the pair [ $(A-$ $\left.W_{12} W_{22}^{-1} C_{1}\right), B_{11}$ ] lies on the imaginary axis, then Assumptions 4.3 and 4.4 are verified. In fact, eqs. (4.71)-(4.74),(4.79) are readily seen to imply the fulfillment of Assumption 4.4, while, by performing the required substitutions, it is easy to ascertain that the unreachable part of the pair [ $\left.\left(A-W_{12} W_{22}^{-1} C_{1}\right), B_{11}\right]$ coincides with the unreachable part of the pair
$\left[\left(A-B_{1} D_{21}^{\prime} C_{2}\right), B_{1}\left(I-D_{21}^{\prime} D_{21}\right)\right]$. As for detectability of the pair $\left(A, C_{2}\right)$, it is equivalent to the detectability of the pair $(A, C)$, thanks to eq. (4.73).

Finally, if matrix $A+B_{2} S$ is stable, so that also Assumption 4.5 is satisfied, then Theorem 4.2 (point $(b)$ ) ensures that the controller which solves the underlying problem relative to system (4.75)-(4.79) has the form (4.80), (4.81) with $N=0$. Precisely, the controller is

$$
\begin{align*}
\dot{\xi} & =\left(A-B_{2} C_{1}+L_{2} C_{2}\right) \xi+L_{2} y  \tag{4.82}\\
u & =C_{1} \xi \tag{4.83}
\end{align*}
$$

and provides a solution of the considered filtering problem. Here $L_{2}$ is given by eqs. (4.61), (4.63) and (4.64). However, the filtering problem for system (4.68)-(4.70) could have been more classically tackled via Kalman theory, yielding the filter

$$
\begin{align*}
\dot{\xi}_{K} & =\left(A+L_{K} C\right) \xi_{K}-L_{K} \bar{y}  \tag{4.84}\\
u_{K} & =S \xi_{K} \tag{4.85}
\end{align*}
$$

where $u_{K}$ is the optimal estimate of $S x$. In eq. (4.84) it is

$$
L_{K}=-\left(\Pi_{K} C^{\prime}+W_{12}\right) W_{22}^{-1}
$$

and $\Pi_{K}$ is the symmetric, positive definite and minimal solution of the Riccati equation (in the unknown $\Pi$ )

$$
0=\Pi\left(A-W_{12} W_{22}^{-1} C\right)^{\prime}+\left(A-W_{12} W_{22}^{-1} C\right) \Pi-\Pi C^{\prime} W_{22}^{-1} C \Pi+\hat{W}
$$

where, again, $\hat{W}=W_{11}-W_{12} W_{22}^{-1} W_{12}^{\prime}$. Notice that this equation coincides with eq. (4.64) once the substitutions (4.71)-(4.73) have been performed.

The assumption which guarantees the existence of $\Pi_{K}$ (namely, the stability of the unobservable but reachable part of system $\Sigma\left(A-W_{12} W_{22}^{-1} C, \hat{W}^{1 / 2}, C, 0\right)$ ), is weaker than those assuring the existence of the stabilizing (i.e. maximal) solution of eq. (4.64). Hence, in general, $\Pi_{K}$ might exist and $\Pi_{2}$ not exist; moreover it can also happen that $\Pi_{K} \neq \Pi_{2}$ so that the $\mathrm{RH}_{2}$ and Kalman filtering problems may substantially differ one from the other.

However, if the Kalman filter is required to be stable, then $\Pi_{2}=\Pi_{K}$ and $L_{2}=$ $L_{K} W_{22}$. Despite of being the stable Kalman filter a device different from the $R H_{2}$ controller, the transfer function from the input noise $\zeta$ to the estimation error $S\left(\xi_{K}-x\right)$ (system (4.68),(4.69),(4.84) and (4.85)) coincides with the transfer function from the input noise $w$ to the estimation error $z=C_{1}(x+\xi)$ (system (4.75)-(4.77), (4.82) and (4.83)). On the contrary, the two devices do coincide in the particular case where $B_{2}=0, W_{12}=0$ and $A$ stable.

Example 4.4 Consider system (4.68)-(4.70) with

$$
\begin{aligned}
A & =\left[\begin{array}{rr}
0 & 1 \\
-1 & 0
\end{array}\right], \quad B_{2}=\left[\begin{array}{r}
0 \\
-1
\end{array}\right], \quad C=\left[\begin{array}{ll}
1 & 0
\end{array}\right] \\
W_{11} & =\left[\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right], W_{12}=\left[\begin{array}{l}
0 \\
0
\end{array}\right], W_{22}=1
\end{aligned}
$$

The Kalman filter relative to the linear combination $\eta:=S x$ with $S=\left[\begin{array}{ll}1 & 1\end{array}\right]$ is described by the equations

$$
\begin{aligned}
& \dot{\xi}_{K}=\left(A+L_{K} C\right) \xi_{K}+B_{2} u-L_{K} \bar{y} \\
& u_{K}=S \xi
\end{aligned}
$$

where $L_{K}=-[1.350 .41]$ and $\left(A+L_{K} C\right)$ is stable since the relevant Riccati equation admits a unique symmetric and positive definite solution which is also stabilizing.


Figure 4.14: The full control problem

If, on the contrary, Theorem 4.2 has to be exploited according to Remark 4.11, the controller (corresponding to the choice $Q(s)=0$ ) is obtained

$$
\begin{aligned}
\dot{\xi} & =\left(A-B_{2} C_{1}+L_{2} C_{2}\right) \xi+L_{2} y \\
u & =C_{1} \xi
\end{aligned}
$$

where $C_{1}=-S, C_{2}=C, y=\bar{y}, L_{2}=L_{K}$. Thus the optimal estimate of $\eta$ is $-C_{1} \xi$.
Remark 4.12 (Full control) Here reference is made to the block-scheme of fig. 4.14 where the control vector $u$ is partitioned into two components $u_{1}$ and $u_{2}$. Assume that the first one of them acts in a direct way on the state derivative only, while the second one directly affects the performance output only. More precisely, the considered system $P_{F C}(s)$ is described by the equations

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+[I 0] u  \tag{4.86}\\
z & =C_{1} x+[0 I] u  \tag{4.87}\\
y & =C_{2} x+D_{21} w  \tag{4.88}\\
u & =\left[\begin{array}{lll}
u_{1}^{\prime} & u_{2}^{\prime}
\end{array}\right]^{\prime} \tag{4.89}
\end{align*}
$$

Further, let the pair ( $A, C_{2}$ ) be detectable, $D_{21} D_{21}^{\prime}=I$ and assume that no eigenvalue of the unreachable part of the pair $\left[\left(A-B_{1} D_{21}^{\prime} C_{2}\right), B_{1}\left(I-D_{21}^{\prime} D_{21}\right)\right]$ lies on the imaginary axis.

First, observe that it makes sense dealing with Problem 4.1 relative to the system above since the necessary condition in Remark 4.1 is verified. Indeed, stabilizability of the pair $\left(A, B_{2}\right)$ is guaranteed by the form of matrix $B_{2}$, while detectability of the pair $\left(A, C_{2}\right)$ holds by assumption. Now consider system $\hat{P}_{F C}(s):=P_{F C}^{\prime}(s)$. From eqs. (4.86)-(4.89) it follows that $\hat{P}_{F C}(s)$ is described by

$$
\begin{aligned}
\dot{\xi} & =F \xi+G_{1} \zeta+G_{2} \eta \\
\omega & =H_{1} \xi+E \eta \\
\varphi & =\left[\begin{array}{l}
I \\
0
\end{array}\right] \xi+\left[\begin{array}{l}
0 \\
I
\end{array}\right] \zeta
\end{aligned}
$$

where

$$
F:=A^{\prime}, \quad E:=D_{21}^{\prime}, \quad H_{1}:=B_{1}^{\prime}
$$

and, for $i=1,2$,

$$
G_{i}:=C_{i}^{\prime}
$$

Therefore, system $\hat{P}_{F C}(s)$ possesses the same structure as system $P_{F I}(s)$ which has been considered in Section 4.2. The assumptions on system $P_{F C}(s)$ make system $\hat{P}_{F C}(s)$ to satisfy Assumptions 4.1, 4.2, so that the results concerning the solution of Problem 4.1 relative to system $P_{F C}(s)$ can be derived, thanks to Lemma E.1, by transposing those concerning system $\hat{P}_{F C}(s)$ which, in turn, coincide with those supplied by Theorem 4.1 for the full information problem. Thus one obtains
a)

$$
\min \|T(z, w ; s)\|_{2}=\left\|C_{1} P_{f}(s)\right\|_{2}=\sqrt{\operatorname{trace}\left[C_{1} \Pi_{2} C_{1}^{\prime}\right]}
$$



Figure 4.15: The set $\mathcal{F}_{2 \gamma r}$ of the controllers $K_{F C r}(s)$


Figure 4.16: The generic admissible controller for $P_{F C}(s)$
b)

$$
K_{F C}^{o}(s)=\left[\begin{array}{c|c}
\emptyset & \emptyset \\
\hline \emptyset & L_{2} \\
\emptyset & 0
\end{array}\right]
$$

c) The set $\mathcal{F}_{2 \gamma r}$ of the controllers $K_{F C r}(s)$ is defined in the block-scheme of fig. 4.15, where $Q(s):=\Sigma\left(A_{q}, B_{q}, C_{q}, 0\right)$, with the matrix $A_{q}$ stable and $\|Q(s)\|_{2}^{2}<\gamma^{2}-$ $\left\|C_{1} P_{f}(s)\right\|_{2}^{2}$.
In the three points above $\gamma>\left\|C_{1} P_{f}(s)\right\|_{2}$ and reference has been made to eqs. (4.61)-(4.65).
Finally, the set $\mathcal{F}_{2 \gamma}$ of the $\mathrm{RH}_{2}$ admissible controllers for system $P_{F C}(s)$ can easily be found by exploiting (through transposition) the content of Remark 4.3 (which refers to system $P_{F I}(s)$ ). Therefore, this set is described by the block-scheme of fig. 4.16 (which has been obtained by "transposing" fig. 4.5), where $\Theta_{1}(s) \in R H_{\infty}$ and $\Psi(s):=\left(s I-A_{f c}\right)^{-1}$.

Remark 4.13 Assumption 4.4 is somehow restrictive though customary in estimation theory. It implies that matrix $D_{21}$ has rank equal to the number of its rows. Should this assumption not be verified one would have to face a singular problem the solution of which, if any, requires a discussion far beyond the scope of this book.

On the contrary, observe that the condition $D_{21} D_{21}^{\prime}=I$ can be substituted, without troubles, by $D_{21} D_{21}^{\prime}=R$ with $R>0$. Indeed, the present derivation can be exploited by redefining the output variable as shown in eq. (4.73).

Finally, the orthogonality assumption $D_{21} B_{1}^{\prime}=0$ is often made. Besides making simpler the notation, such an assumption implies that the noises $\zeta_{1}$ and $\zeta_{2}$, introduced in Remark 4.11 when dealing with filtering problems, are uncorrelated ( $W_{12}=0$ ).

Remark 4.14 The Riccati equation (4.64) is the one encountered in deriving the Kalman filter for the system $\dot{x}=A x+w_{1}, y=C_{2} x+w_{2}$, where the zero mean Gaussian white noise $w:=\left[\begin{array}{ll}w_{1}^{\prime} & w_{2}^{\prime}\end{array}\right]^{\prime}$ has intensity

$$
W:=\left[\begin{array}{cc}
B_{1} B_{1}^{\prime} & B_{1} D_{21}^{\prime} \\
D_{21} B_{1}^{\prime} & I
\end{array}\right]
$$

Remark 4.15 Under Assumption 4.4, the eigenvalues of the unreachable part of the pair [ $\left.\left(A-B_{1} D_{21}^{\prime} C_{2}\right), B_{1}\left(I-D_{21}^{\prime} D_{21}\right)\right]$ do not lie on the imaginary axis (Assumption 4.3) if and only if the invariant zeros of the subsystem of $P_{O E}(s)$ corresponding to the transfer function $P_{O E 21}(s)$ from the input $w$ to the output $y$, namely the system $\Sigma\left(A, B_{1}, C_{2}, D_{21}\right)$, all have real part different from zero. In fact, from Section 2.5 and by recalling that the number of components of the disturbance $w$ is not smaller than the number of components of the output $y$, if a scalar $\lambda$ with $\operatorname{Re}(\lambda)=0$ exists such that

$$
\begin{gathered}
\left(\lambda I-A^{\prime}\right) x-C_{2}^{\prime} y=0 \\
B_{1}^{\prime} x+D_{21}^{\prime} y=0
\end{gathered}
$$

with $\left[x^{\prime} y^{\prime}\right]^{\prime} \neq 0$, then, from the identity $D_{21} D_{21}^{\prime}=I$ it follows

$$
y=-D_{21} B_{1}^{\prime} x
$$

so that

$$
\begin{gathered}
\left(A^{\prime}-C_{2}^{\prime} D_{21} B_{1}^{\prime}\right) x=\lambda x \\
\left(I-D_{21}^{\prime} D_{21}\right) B_{1}^{\prime} x=0
\end{gathered}
$$

Being $x \neq 0$, since, otherwise, also $y=0$, these two equations would violate the assumption that no eigenvalue of the unreachable part of the pair $\left[\left(A-B_{1} D_{21}^{\prime} C_{2}\right), B_{1}\left(I_{D} 21^{\prime} D_{21}\right)\right]$ lies on the imaginary axis (recall Lemma D.3, point (a)). On the contrary, if the two last relations hold true for a certain $x \neq 0$, then letting $y:=-D_{21} B_{1}^{\prime} x$ and proceeding in the reverse way the conclusion straightforwardly follows.

Remark 4.16 The above results can be generalized to the fairly frequent case in which the output variable $y$ explicitly depends on the control variable $u$, that is when eq. (4.60) is substituted by

$$
\begin{equation*}
y=C_{2} x+D_{12} w+D_{22} u \tag{4.90}
\end{equation*}
$$

Indeed, letting

$$
\begin{equation*}
\bar{y}:=y-D_{22} u=C_{2} x+D_{12} w \tag{4.91}
\end{equation*}
$$

Problem 4.1 relative to system $\bar{P}_{O E}(s)$ (described by eqs. (4.58),(4.59), (4.91)) admits the solution presented in Theorem 4.2. If $\bar{K}_{O E}(s)$ is a $R H_{2}$ admissible controller for $\bar{P}_{O E}(s)$, then the controller $K_{O E}(s)$ defined in the block-scheme of fig. 4.17 is apparently a $\mathrm{RH}_{2}$ admissible controller for system $P_{O E}(s)$ (described by eqs. (4.58),(4.59),(4.90)), only provided that it is well defined. This is certainly the case since the set of controllers $\bar{K}_{O E}(s)$ is constituted by strictly proper systems (see Theorem 4.2).

### 4.4 The partial information problem

In the control problem considered in the present section only a partial information on the system state is available to the controller. Therefore, the controlled system is described by

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{4.92}\\
z & =C_{1} x+D_{12} u  \tag{4.93}\\
y & =C_{2} x+D_{21} w \tag{4.94}
\end{align*}
$$



Figure 4.17: The controller $K_{O E}(s)$ when $y$ directly depends on $u$


Figure 4.18: The partial information problem
and its transfer function is denoted by $P(s)$. Further, the following assumptions are done.

Assumption 4.6 The pair $\left(A, B_{2}\right)$ is stabilizable and the pair $\left(A, C_{2}\right)$ is detectable.
Assumption 4.7 $D_{12}^{\prime} D_{12}=I$.
Assumption 4.8 The eigenvalues of the unobservable and unreachable part of the pairs $\left[\left(A-B_{2} D_{12}^{\prime} C_{1}\right)\left(I-D_{12} D_{12}^{\prime}\right) C_{1}\right]$ and $\left[\left(A-B_{1} D_{21}^{\prime} C_{2}\right), B_{1}\left(I-D_{21}^{\prime} D_{21}\right)\right]$ respectively, do not lie on the imaginary axis.

Assumption 4.9 $D_{21} D_{21}^{\prime}=I$.
In the forthcoming theorem reference is made to the block-scheme of fig. 4.18 where $K(s)$ denotes a generic $R H_{2}$ admissible controller for $P(s)$ and $T(z, w ; s)$ is the transfer function from $w$ to $z$.

Theorem 4.3 (Partial information) Consider system (4.92)-(4.94). Then, under Assumptions 4.6-4.9, Problem 4.1 has the following solution.
a)

$$
\begin{aligned}
\min \|T(z, w ; s)\|_{2}^{2} & =\left\|P_{c}(s) B_{1}\right\|_{2}^{2}+\left\|F_{2} P_{f}(s)\right\|_{2}^{2} \\
& =\left\|P_{c}(s) L_{2}\right\|_{2}^{2}+\left\|C_{1} P_{f}(s)\right\|_{2}^{2} \\
& :=\left(\gamma^{o}\right)^{2}
\end{aligned}
$$

b)

$$
K^{o}(s):=\left[\begin{array}{c|c}
A+B_{2} F_{2}+L_{2} C_{2} & -L_{2} \\
\hline F_{2} & 0
\end{array}\right]
$$



Figure 4.19: The set $\mathcal{F}_{2 \gamma r}$ of the controllers $K_{\text {FIr }}(s)$
c) The set $\mathcal{F}_{2 \gamma r}$ of the controller $K_{r}(s)$ is defined by the block-scheme of fig. 4.19 where

$$
S_{2}(s):=\left[\begin{array}{c|cc}
A+B_{2} F_{2}+L_{2} C_{2} & -L_{2} & B_{2} \\
\hline F_{2} & 0 & I \\
-C_{2} & I & 0
\end{array}\right]
$$

$$
Q(s):=\Sigma\left(A_{q}, B_{q}, C_{q}, 0\right), \text { the matrix } A_{q} \text { is stable and }\|Q(s)\|_{2}^{2}<\gamma^{2}-\left(\gamma^{o}\right)^{2}
$$

In the three points above reference has been made to eqs. (4.15)-(4.18) and (4.61)(4.64), while $\gamma$ is a positive scalar such that $\gamma>\gamma^{\circ}$.

Proof First observe that Assumption 4.6 coincides with the necessary condition for the problem at hand to make sense (see Remark 4.1). Then notice that the above assumptions guarantee, thanks to Lemma C.3, the existence of the solutions $P_{2}$ and $\Pi_{2}$ of the Riccati equations (4.18) and (4.64) endowed with the relevant properties. Repeat now the first part of the proof of Theorem 4.1 (eqs. (4.20)-(4.26)) by making reference to fig. 4.18 and 4.19 rather than to fig. 4.2 and 4.3 , respectively, and defining the system $P_{v}(s)$ in fig. 4.4 as

$$
P_{v}(s):=\left[\begin{array}{c|cc}
A & B_{1} & B_{2} \\
\hline-F_{2} & 0 & I \\
C_{2} & D_{21} & 0
\end{array}\right]
$$

Then, one gets

$$
\|T(z, w ; s)\|_{2}^{2}=\left\|P_{c}(s) B_{1}\right\|_{2}^{2}+\|T(v, w ; s)\|_{2}^{2}
$$

On the other hand, system $P_{v}(s)$ has the same structure as system $P_{O E}(s)$ (the system considered in Section 4.3) and equals it if one let $-F_{2}=C_{1}$. It is easy to ascertain that under Assumptions 4.6-4.9, Assumptions 4.3-4.5 are verified for system $P_{v}(s)$. This is straightforward as for the first two of them, while Assumption 4.5, namely stability of the matrix $A-B_{2} C_{1}=A+B_{2} F_{2}$, follows from $P_{2}$ being the stabilizing solution of the Riccati equation (recall eq. (4.19)). Since the set of $\mathrm{RH}_{2}$ admissible controllers for $P(s)$ coincides with the set of $R H_{2}$ admissible controllers for $P_{v}(s)$, it is possible to minimize $\|T(v, w ; s)\|_{2}$ by resorting to Theorem 4.2. Thus the first equality sign in point $a$ ) and points $b$ ) and $c$ ) follow.

As for the second equality sign in point $a$ ), notice that system $\hat{P}(s):=P^{\prime}(s)$ has the same structure as $P(s)$. Therefore, the solution of Problem 4.1 relative to system $\hat{P}(s)$ is fully described by the statement of Theorem 4.3 apart, for the moment being,
from the second part of point $a$ ). The relevant results can be utilized also for system $P(s)$, thanks to Lemma E.1. In particular, if a cap sign marks the items concerning system $\hat{P}(s)$ which correspond to those introduced for system $P(s)$, in view of point (a) of Theorem 4.3 , first equality sign, one obtains

$$
\begin{aligned}
\min \|\hat{T}(\hat{z}, \hat{w} ; s)\|_{2}^{2} & =\left\|\hat{P}_{c}(s) \hat{B}_{1}\right\|_{2}^{2}+\left\|\hat{F}_{2} \hat{P}_{f}(s)\right\|_{2}^{2} \\
& =\left\|P_{c}(s) L_{2}\right\|_{2}^{2}+\left\|C_{1} P_{f}(s)\right\|_{2}^{2}
\end{aligned}
$$

since $\hat{P}_{c}^{\prime}(s)=P_{f}(s), \hat{B}_{1}^{\prime}=C_{1}, \hat{P}_{f}^{\prime}(s)=P_{c}(s)$ and $\hat{F}_{2}^{\prime}=L_{2}$, as it can be verified.
Example 4.5 Consider system (4.92)-(4.94) with

$$
\begin{gathered}
A=\left[\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right], \quad B_{1}=\left[\begin{array}{ll}
1 & 0 \\
1 & 0
\end{array}\right], \quad B_{2}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad D_{12}=\left[\begin{array}{l}
0 \\
1
\end{array}\right] \\
C_{1}=\left[\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right], \quad C_{2}=\left[\begin{array}{ll}
1 & 0
\end{array}\right], \quad D_{21}=\left[\begin{array}{ll}
0 & 1
\end{array}\right]
\end{gathered}
$$

One obtains $F_{2}=-\left[\begin{array}{ll}1 & \sqrt{2}\end{array}\right], L_{2}=-\left[\begin{array}{ll}\sqrt{3} & 1\end{array}\right]^{\prime}$ and

$$
\left(\gamma^{o}\right)^{2}:=\left\|P_{c}(s) B_{1}\right\|_{2}^{2}+\left\|F_{2} P_{f}(s)\right\|_{2}^{2}=10.85
$$

Taken $\gamma=16$ and

$$
Q(s)=\left[\begin{array}{rr|r}
-1 & 0 & 1 \\
0 & -3 & 1 \\
\hline 1 & 1 & 0
\end{array}\right]
$$

it is $\|T(z, w ; s)\|_{2}^{2}=12.02<\gamma^{2}$, consistently with $1.17=\|Q(s)\|_{2}^{2}<\gamma^{2}-\left(\gamma^{\circ}\right)^{2}$.
Example 4.6 Consider system (4.92)-(4.94) with

$$
\begin{gathered}
A=\left[\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right], \quad B_{1}=B_{2}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad D_{12}=D_{21}=1 \\
C_{1}=\left[\begin{array}{ll}
1 & \beta
\end{array}\right], \quad C_{2}=\left[\begin{array}{ll}
1 & \alpha
\end{array}\right]
\end{gathered}
$$

where $\alpha \neq 0$ and $\beta \neq 0$. It is easy to verify that the controller which makes the control variable $u$ to depend on the output variable $y$ through the transfer function

$$
K_{I C}(s):=-\frac{\beta s+1}{s^{2}+(\alpha+\beta) s+2}
$$

performs the indirect perfect compensation of the disturbance $w$, since the transfer function from $w$ to $z$ is zero. However, the resulting system is stable only for $\alpha \geq 0$ and $\beta \geq 0$. Therefore, this kind of solution is no more feasible for all other values of the pair $(\alpha, \beta)$.

On the contrary, by applying Theorem 4.3 one obtains (corresponding to the choice $Q(s)=0) K^{o}(s)=K_{I C}(s)$ when $\alpha>0$ and $\beta>0$, otherwise

$$
\begin{aligned}
& K^{o}(s)=\frac{(2 \alpha+\beta) s-1}{s^{2}-(\alpha+\beta) s+2 \alpha^{2}+2 \alpha \beta+2}, \alpha<0, \beta<0 \\
& K^{o}(s)=\frac{\beta s-1}{s^{2}+(\alpha-\beta) s+2}, \alpha>0, \quad \beta<0 \\
& K^{o}(s)=\frac{(2 \alpha-\beta) s-1}{s^{2}+(\beta-\alpha) s+2 \alpha^{2}-2 \alpha \beta+2}, \alpha<0, \quad \beta>0
\end{aligned}
$$

Remark 4.17 The structure of a generic $\mathrm{RH}_{2}$ admissible controller for $P(s)$ (as defined by the block-scheme of fig. 4.19) allows checking that the eigenvalues of the resulting control system are those of the matrices $A+B_{2} F_{2}, A+L_{2} C_{2}$ and $A_{q}$. In fact, letting $x_{s}$ and $x_{q}$ denote the state variables of systems $S_{2}(s)$ and $Q(s)$, respectively, the dynamic matrix of the resulting system (with state $x_{t}:=\left[x^{\prime}-x_{s}^{\prime} x_{q}^{\prime} x^{\prime}\right]^{\prime}$ ) is

$$
A_{t}=\left[\begin{array}{ccc}
A+L_{2} C_{2} & 0 & 0 \\
B_{q} C_{2} & A_{q} & 0 \\
B_{2} F_{2} & B_{2} C_{q} & A+B_{2} F_{2}
\end{array}\right]
$$

Remark 4.18 The optimal controller given in point $b$ ) of Theorem 4.3 may be interpreted as the result of a synthesis procedure made up of two independent steps. The first one consists in solving the full information problem (dealt with in Section 4.2) yielding matrix $F_{2}$. The second step tackles the output estimation problem (dealt with in Section 4.3) relative to the linear state combination $F_{2} x$, that is relative to system (4.58)-(4.60) with $C_{1}=-F_{2}$.

Remark 4.19 (Parametrization of the set $\mathcal{F}_{2 \gamma}$ ) Notice that $\mathcal{F}_{2 \gamma}=\mathcal{F}_{2 \gamma r}$. In fact, as it was done in the proof of Theorem 4.3, recall that the set of $\mathrm{RH}_{2}$ admissible controllers for $P(s)$ coincides with the set of $R H_{2}$ admissible controllers for $P_{v}(s)$. However, the latter system has the same structure as system $P_{O E}(s)$ (considered in Section 4.3), so that, in view of Remark 4.10, the above claim is correct.

Remark 4.20 The control problem addressed to in Theorem 4.3 can be viewed as an optimal linear quadratic stochastic problem with unmeasurable state. In fact, consider the $n$-th order system

$$
\begin{align*}
& \dot{x}=A x+B \bar{u}+\zeta_{1}  \tag{4.95}\\
& \bar{y}=C x+\zeta_{2} \tag{4.96}
\end{align*}
$$

where $\zeta:=\left[\zeta_{1}^{\prime} \zeta_{2}^{\prime}\right]^{\prime}$ is a zero mean white Gaussian noise with intensity

$$
W:=\left[\begin{array}{ll}
W_{11} & W_{12}  \tag{4.97}\\
W_{12}^{\prime} & W_{22}
\end{array}\right], \quad W_{22}>0
$$

Also consider the cost functionals

$$
J_{6}:=\lim _{t \rightarrow \infty} E\left[x^{\prime}(t) Q x(t)+2 x^{\prime}(t) S \bar{u}(t)+\bar{u}^{\prime}(t) R \bar{u}(t)\right]
$$

and

$$
J_{7}:=\lim _{T \rightarrow \infty} E\left[\frac{1}{T} \int_{0}^{T}\left[x^{\prime}(t) Q x(t)+2 x^{\prime}(t) S \bar{u}(t)+\bar{u}^{\prime}(t) R \bar{u}(t)\right] d t\right]
$$

where

$$
\left[\begin{array}{ll}
Q & S \\
S^{\prime} & R
\end{array}\right]:=L=L^{\prime} \geq 0, \quad R>0
$$

Letting

$$
\hat{Q}:=Q-S R^{-1} S^{\prime}, \quad \hat{W}:=W_{11}-W_{12} W_{22}^{-1} W_{12}^{\prime}
$$

which are positive semidefinite (recall Remarks 4.4 and 4.11) and

$$
\begin{aligned}
C_{11}^{\prime} C_{11} & :=\hat{Q}, C_{11} \in R^{n \times n}, \quad B_{11} B_{11}^{\prime}:=\hat{W}, B_{11} \in R^{n \times n} \\
B_{2} & :=B R^{-1 / 2}, \quad B_{1}:=\left[\begin{array}{ll}
B_{11} & W_{12} W_{22}^{-1 / 2}
\end{array}\right] \\
D_{21} & :=\left[\begin{array}{ll}
0 & I
\end{array}\right], C_{2}:=W_{22}^{-1 / 2} C \\
C_{1} & :=\left[\begin{array}{c}
C_{11} \\
R^{-1 / 2} S^{\prime}
\end{array}\right], \quad D_{12}:=\left[\begin{array}{c}
0 \\
I
\end{array}\right], u:=R^{1 / 2} \bar{u} \\
y & :=W_{22}^{-1 / 2} \bar{y}, \quad z:=C_{1} x+D_{12} u
\end{aligned}
$$

system (4.95)-(4.97) can be rewritten as

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{4.98}\\
z & =C_{1} x+D_{12}  \tag{4.99}\\
y & =C_{2} x+D_{21} \tag{4.100}
\end{align*}
$$

where $w$ is a zero mean white Gaussian noise with identity intensity and dimension $n+p$ which satisfies the equation

$$
\zeta=\left[\begin{array}{cc}
B_{11} & W_{12} W_{22}^{-1 / 2} \\
0 & W_{22}^{1 / 2}
\end{array}\right] w
$$

It is also easy to verify that

$$
J_{6}:=\lim _{t \rightarrow \infty} E\left[z^{\prime}(t) z(t)\right]
$$

and

$$
J_{7}:=\lim _{T \rightarrow \infty} E\left[\frac{1}{T} \int_{0}^{T} z^{\prime}(t) z(t) d t\right]
$$

The problem under consideration is finding a controller of the form

$$
\begin{align*}
\dot{\xi} & =F \xi+G y  \tag{4.101}\\
u & =H \xi+N y \tag{4.102}
\end{align*}
$$

such that system (4.98)-(4.102) is stable and $J_{6}$ or $J_{7}$ is minimized. Notice that the feedback connection of any controller described by eqs. (4.101),(4.102) with system (4.98)-(4.100) is well defined and the relevant transfer function from $w$ to $z$ is strictly proper if and only if $N D_{21}=0$, that is if and only if $N=0$ (recall the definition of $D_{21}$ ). Therefore, if system (4.98)-(4.102) is stable and $N=0$, then the controller (4.101),(4.102) is $R H_{2}$ admissible for system (4.98)-(4.100). In view of Remark 2.20 any $\mathrm{RH}_{2}$ admissible controller for system (4.98)-(4.100) is such that

$$
J_{6}=J_{7}=\|T(z, w ; s)\|_{2}^{2}
$$

Assume that neither the eigenvalues of the unobservable part of the pair [ $\left(A-B R^{-1} S^{\prime}\right), C_{11}$ ] nor those of the unreachable part of the pair $\left[\left(A-W_{12} W_{22}^{-1} C\right), B_{11}\right]$ lie on the imaginary axis. Moreover, assume that the pair $(A, C)$ is detectable and the pair $(A, B)$ is stabilizable. With the same kind of reasoning developed in Remark 4.11 it is easy to see that these assumptions are equivalent to Assumptions 4.6 and 4.8 , while Assumptions 4.7 and 4.9 are satisfied because of the definition of $D_{12}$ and $D_{21}$.

Therefore, Theorem 4.3 can be applied to system (4.98)-(4.100) and supplies the optimal controller which is described (with reference to system (4.95)-(4.97)) by

$$
\begin{align*}
& \dot{\xi}=A \xi+B \bar{u}+L(C \xi-\bar{y})  \tag{4.103}\\
& \bar{u}=\bar{H} \xi \tag{4.104}
\end{align*}
$$

if the relevant substitutions have been done. In eqs. (4.103),(4.104) it is $L:=-\left(\Pi_{2} C^{\prime}+\right.$ $\left.W_{12}\right) W_{22}^{-1}, \bar{H}:=-R^{-1}\left(B^{\prime} P_{2}+S^{\prime}\right), P_{2}$ and $\Pi_{2}$ being the symmetric, positive semidefinite and stabilizing solutions of the Riccati equations (in the unknown $P$ and $\Pi$, respectively)

$$
\begin{aligned}
& 0=P A_{c}+A_{c}^{\prime} P-P B R^{-1} B^{\prime} P+Q-S R^{-1} S^{\prime} \\
& 0=\Pi A_{f}^{\prime}+A_{f} \Pi-\Pi C^{\prime} W_{22}^{-1} C \Pi+W_{11}-W_{12} W_{22}^{-1} W_{12}^{\prime}
\end{aligned}
$$

with $A_{c}:=A-B R^{-1} S^{\prime}, A_{f}:=A-W_{12} W_{22}^{-1} C$. Notice that in the controller (4.103),(4.104) it is $N=0$.

Equation (4.103) is the equation of the stable (see Remark 4.11) Kalman filter for system (4.95)-(4.97). Thus, the controller (4.103),(4.104) can be considered as a Kalman filter on the state of which the control law has been implemented (eq. (4.104)) which is optimal and
stabilizing for the Linear Quadratic problem with infinite horizon defined on system (4.95) with $\zeta_{1} \equiv 0$ and performance index

$$
\bar{J}_{7}:=\int_{0}^{\infty}\left[x^{\prime}(t) Q x(t)+2 x^{\prime}(t) S \bar{u}(t)+\bar{u}^{\prime}(t) R \bar{u}(t)\right] d t
$$

The criterion $\bar{J}_{7}$ is the "deterministic version" of the functional $J_{7}$. This structural separation of the solution (optimal filtering and regulator problem) is outlined also by the spectrum of the dynamic matrix of the resulting system. In fact, letting $e:=\xi-x$, from eqs. (4.95)(4.97),(4.103),(4.104) it follows

$$
\begin{aligned}
& \dot{e}=(A+L C) e-\zeta_{1}-L \zeta_{2} \\
& \dot{x}=(A+B \bar{H}) x+B \bar{H} e+\zeta_{1}
\end{aligned}
$$

The problem of minimizing either the functional $J_{6}$ or the functional $J_{7}$ under the unique constraint expressed by eqs. (4.95)-(4.97) is known as the linear quadratic Gaussian ( $L Q G$ ) problem. The solution of such a problem (whenever it exists) is specified by eqs. (4.103),(4.104), where, however, the two matrices $\Pi_{2}$ and $P_{2}$ which determine $L$ and $\bar{H}$, respectively, may not be the stabilizing solutions of the relevant Riccati equations. Consequently, system (4.95)-(4.97), (4.103),(4.104) may be unstable. This outcome is consistent with the absence of any stability requirement put forth by the classical $L Q G$ theory. As a matter of fact, the assumptions required by such a theory (stability of the unreachable but observable part of system $\Sigma\left(A-B R^{-1} S^{\prime}, B, C_{11}, 0\right)$ and of the unobservable but reachable part of system $\left.\Sigma\left(A-W_{12} W_{22}^{-1} C, B_{11}, C, 0\right)\right)$ are weaker than those required within the $R H_{2}$ context.

Remark 4.21 The solution of the output estimation problem (see Section 4.2) can be derived as an application of Theorem 4.3 to a particular case. In fact, it suffices to set $D_{12}=I$ in Assumptions 4.6-4.9 in order to conclude that Assumptions 4.3-4.5 are satisfied. In particular, it results $C_{1 \mathrm{c}}=0$ so that the (unique) symmetric, positive semidefinite and stabilizing solution of the Riccati equation (4.18) is $P_{2}=0$. In such a context, the conclusions of Theorem 4.3 are immediately redrawn to those of Theorem 4.2. In a similar way, the solution of the disturbance feedforward problem dealt with in Remark 4.5 can be obtained by solving the partial information problem in the particular case $D_{21}=I$. Indeed, letting $D_{21}=I$ in Assumptions 4.6-4.9, it is easy to verify that the assumptions made in Remark 4.5 are satisfied. In particular, it results $B_{1 f}=0$ and hence $\Pi_{2}=0$, so that the conclusions of Theorem 4.3 coincide with those illustrated in Remark 4.5.

Remark 4.22 The contents of Remark 4.6 (as for matrix $D_{12}$ ) and of Remark 4.13 (as for matrix $D_{21}$ ) apply with no changes to the problem dealt with in the present section.

Remark 4.23 In view of Remarks 4.8 and 4.15 it can be said that, under Assumptions 4.7 and 4.9 , Assumption 4.8 amounts to requiring that the two subsystems of $P(s)$ having transfer functions $P_{12}(s)$ (that is, system $\Sigma\left(A, B_{2}, C_{1}, D_{12}\right)$ with input $u$ and output $z$ ) and $P_{21}(s)$ (that is, system $\Sigma\left(A, B_{1}, C_{2}, D_{21}\right)$ with input $w$ and output $\left.y\right)$, respectively, do not have zeros on the imaginary axis.

Remark 4.24 The discussion in Remark 4.16 can be applied with no changes to the present context. Therefore, the results presented in Theorem 4.3 can be extended with no difficulty to encompass the case in which eq. (4.94) is replaced by

$$
y=C_{2} x+D_{21} w+D_{22} u
$$

### 4.5 Notes and references

The material of this section mostly relies on the paper by Doyle et al. [17]. However, the assumptions under which Theorems 4.1, 4.2, 4.3 have been stated are more general than those adopted in such a paper. Moreover the results concerning the parametrization of the admissible controllers have been modified so as to encompass the remark put forward by Mita et al. [44]. Further insight on the connections existing between the $L Q G$ and $R H_{2}$ problems was given by Kucera [35].

## Chapter 5

## $R H_{\infty}$ Control

### 5.1 Introduction

The control problem for linear time-invariant systems has been classically tackled in the frequency domain. A typical, though not completely general, context is the one shown in fig. 5.1, where $G(s)$ and $K(s)$ are the transfer functions of the process to be controlled (possibly including actuators and sensors) and of the controller to be synthesized. The disturbances $d_{c}, d_{u}, d_{r}$ act on the controlled variable $c$, on the process input variable $u_{p}$ and on the feedback path, respectively. Finally, $c^{o}$ is the opposite of the set point.

As well known, the aim in designing $K(s)$ is, loosely speaking, guaranteeing the stability of the control system and achieving satisfactory performances. Usually, such performances are evaluated in terms of the behavior of suitable variables of interest to be specified according to the problem at hand and must be attained in spite of the disturbances acting on the system and inaccurate knowledge of the process model.

In general, the philosophy underlying the adopted synthesis procedure strongly affects the result: for instance, having either ignored or taken into account the inaccurate knowledge of the process model makes the controller quite different. Moreover, the design procedure significantly depends on the adopted description of the uncertainty: thus, subdividing the relevant discussion (and consequently this section) into three parts, appears to be a fairly natural strategy. First, the design problem is faced in nominal conditions, that is the process is supposed perfectly known. Second, some


Figure 5.1: A typical control system
typical situations in which the process model is not precisely known are shown to be conveniently described by inserting a suitable perturbation $\Delta(s)$ into the block scheme of fig. 5.1. The topology of this scheme does therefore reflects the particular nature of the considered uncertainty. Third, the design problem is tackled in a robust way, namely it is stated within various uncertainty scenarios, modeled in accordance with the previous discussion.

### 5.1.1 Nominal design

The process under control is assumed to be perfectly described by its nominal transfer function $G_{n}(s)$ so that in the block scheme of fig. 5.1 $G(s)=G_{n}(s)$. The transfer functions which are suited to evaluate the effects of the disturbances and hence can be exploited to express meaningful performance requirements are the following:

$$
\begin{align*}
& S_{n}(s)=\left[I-G_{n}(s) K(s)\right]^{-1}  \tag{5.1}\\
& T_{n}(s)=G_{n}(s) K(s)\left[I-G_{n}(s) K(s)\right]^{-1}  \tag{5.2}\\
& V_{n}(s)=K(s)\left[I-G_{n}(s) K(s)\right]^{-1} \tag{5.3}
\end{align*}
$$

The function $S_{n}(s)$, usually referred to as sensitivity function, describes the effect of the disturbance $d_{c}$ on the controlled variable $c$ and the effect of the signal $c^{o}$ on the controller input variable $y$. The function $T_{n}(s)$, which is readily recognized to equal $S_{n}(s)-I$, is, for such a reason, referred to as complementary sensitivity function and accounts for the effect of $c^{o}$ or $d_{r}$ on the controlled variable. Finally, the function $V_{n}(s)$ is responsible of the effect of $c^{o}, d_{r}$ or $d_{c}$ on the process input: thus, it will be referred to as input sensitivity function.

A good solution to the design problem spontaneously calls for making small, in some suitable sense to be specified, the effects of the disturbances on the variables of interest. Thus, the desire of making small the above introduced transfer functions naturally arises.

In the scalar case (that is the case where all the relevant variables are scalar), the desire of making small a transfer function $\varphi(s)$ is consistent with the request that the absolute value of $\varphi(j \omega)$ be, for each frequency $\omega$, smaller than a given (possibly frequency dependent) quantity, namely $|\varphi(j \omega)|<\vartheta(\omega), \forall \omega$. For this to make sense, it is necessary that no poles of $\varphi(s)$ lie on the imaginary axis: in the here considered framework this requirement is naturally fulfilled since, in view of the unavoidable stability constraint, all the transfer functions of interest must belong to $R H_{\infty}$. The natural extension of this philosophy to the multivariable case leads to asking that an inequality of the above type be verified by the maximum singular value of $\varphi(j \omega)$, namely

$$
\begin{equation*}
\bar{\sigma}[\varphi(j \omega)]<\vartheta(\omega), \quad \forall \omega \tag{5.4}
\end{equation*}
$$

Let $W(s) \in R H_{\infty}$ be any matrix such that $W^{-1}(s) \in R H_{\infty}$ and $\bar{\sigma}\left[W^{-1}(j \omega)\right]=$ $\vartheta(\omega)$. Then, the inequality (5.4) is no doubt satisfied if

$$
\bar{\sigma}[W(j \omega) \varphi(j \omega)]<1, \quad \forall \omega
$$

that is if

$$
\begin{equation*}
\|W(s) \varphi(s)\|_{\infty}<1 \tag{5.5}
\end{equation*}
$$

since, in view of Lemma 2.21,

$$
\begin{aligned}
\bar{\sigma}[\varphi(j \omega)] & =\bar{\sigma}\left[W^{-1}(j \omega) W(j \omega) \varphi(j \omega)\right] \\
& \leq \bar{\sigma}[W(j \omega) \varphi(j \omega)] \bar{\sigma}\left[W^{-1}(j \omega)\right]
\end{aligned}
$$

The function $W(s)$ is usually referred to as shaping function.
For the control system of fig. 5.1 the more classical request is the sensitivity performance one, which, consistently with eq. (5.5), can be expressed as

$$
\begin{equation*}
\left\|W_{1}(s) S_{n}(s)\right\|_{\infty}<1 \tag{5.6}
\end{equation*}
$$

Not at all less meaningful are the requests related to the complementary sensitivity performance and input sensitivity performance which can be expressed as

$$
\begin{equation*}
\left\|W_{2}(s) T_{n}(s)\right\|_{\infty}<1 \tag{5.7}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\|W_{3}(s) V_{n}(s)\right\|_{\infty}<1 \tag{5.8}
\end{equation*}
$$

respectively.
Equally of interest are the requests calling for the simultaneous satisfaction of more than one of the inequalities (5.6)-(5.8). As an example, a common request involves both the sensitivity and the complementary sensitivity performance by asking for the fulfillment of

$$
\left\|\left[\begin{array}{l}
W_{1}(s) S_{n}(s)  \tag{5.9}\\
W_{2}(s) T_{n}(s)
\end{array}\right]\right\|_{\infty}<1
$$

since (recall Definition 2.24 and Lemmas 2.17 and 2.23) if (5.9) holds, then both (5.6) and (5.7) hold. It should be apparent that a multiple goal as the one expressed by eq. (5.9) must not ignore the intrinsic constraints existing among the involved functions. A wise selection of the shaping functions $W_{i}(s)$ to be associated with the performance functions $S_{n}(s), T_{n}(s), V_{n}(s)$ is therefore mandatory.

With reference to the particular case of eq. (5.9), the selection of the shaping functions must reflect the identity $S_{n}(s)-T_{n}(s)=I$. In general, $\bar{\sigma}\left[S_{n}(j \omega)\right]$ is asked to be small at low frequencies in order to endove the control system with good capabilities of tracking the set point, whose bandwidth is usually limited from above. This can be achieved by selecting a shaping function $W_{1}(s)$ such that $\bar{\sigma}\left[W_{1}^{-1}(j \omega)\right]$ is small at low frequencies and equal to 1 at higher frequencies. On the contrary, $\bar{\sigma}\left[T_{n}(j \omega)\right]$ is requested to be small at high frequencies in order to effectively counteract, for instance, the disturbances in the feedback path, as their spectra are usually located at high frequency. Consistently, the shaping function $W_{2}(s)$ can be selected so as to have $\bar{\sigma}\left[W_{2}^{-1}(j \omega)\right]$ equal to 1 at low frequencies and as small as possible at higher frequencies.

### 5.1.2 Uncertainty description

It is often more realistic to assume that the process model belongs to some specified set $\mathcal{G}$ rather than being perfectly known. Moreover, the so called nominal model $G_{n}(s)$ is usually taken as an element of the set $\mathcal{G}$ and therefore viewed as a first order approximation of the true model $G(s)$. Consistently, a description of the set $\mathcal{G}$ can be performed by parametrizing it by means of a transfer function $\Delta(s)$ belonging to a suitable set $\mathcal{D}_{\alpha}$ : the perturbations which $G_{n}(s)$ may undergo are then defined by the adopted parametrization and the structure of the set $\mathcal{D}_{\alpha}$.

Here reference will be made to unstructured perturbations only, that is to perturbations which are qualified only in terms of their amplitude as specified by the set

$$
\begin{equation*}
\mathcal{D}_{\alpha}:=\left\{\Delta(s) \mid \Delta(s) \in R H_{\infty}, \quad\|\Delta(s)\|_{\infty}<\alpha\right\} \tag{5.10}
\end{equation*}
$$



Figure 5.2: Different uncertainty models

Some particularly meaningful examples of parametrization of $\mathcal{G}$ are presented in the following equations, where the parameter $\Delta(s)$ is any element of the set $\mathcal{D}_{\alpha}$ defined by eq. (5.10):

$$
\begin{align*}
\mathcal{G} & :=\left\{G(s) \mid G(s)=G_{n}(s)+\Delta(s)\right\}  \tag{5.11}\\
\mathcal{G} & :=\left\{G(s) \mid G(s)=G_{n}(s)[I+\Delta(s)]\right\}  \tag{5.12}\\
\mathcal{G} & :=\left\{G(s) \mid G(s)=[I+\Delta(s)] G_{n}(s)\right\}  \tag{5.13}\\
\mathcal{G} & :=\left\{G(s) \mid G(s)=[I-\Delta(s)]^{-1} G_{n}(s)\right\}  \tag{5.14}\\
\mathcal{G} & :=\left\{G(s) \mid G(s)=\left[I-G_{n}(s) \Delta(s)\right]^{-1} G_{n}(s)\right\} \tag{5.15}
\end{align*}
$$

It is easy to verify that each one of the sets (5.11)-(5.15) is suited to describe meaningful types of uncertainties in a fairly natural way (see also fig. 5.2 where a block-scheme version of eqs. (5.11)-(5.15) is presented). The set (5.11) may model an uncertain location of right half plane zeros of $G(s)$ (as an example: $G_{n}(s)=(s-2) /(s+1)(s+2)$, $\Delta(s)=\varepsilon /(s+1)(s+2))$. The sets (5.12) and (5.13) may suitably account for neglected high frequency poles as well as right half plane zeros (as an example: $\Delta(s)=-\varepsilon s /(1+\varepsilon s)$ or $\Delta(s)=-2 /(1+s)$, so that $1+\Delta(s)=1 /(1+\varepsilon s)$ and $1+\Delta(s)=(s-1) /(1+s)$, respectively). These sets can obviously be exploited in describing the model uncertainties of both actuators and sensors as well. The set (5.14) can easily model neglected right half plane poles (as an example: $\Delta(s)=10 /(1+s)$, so that $\left.[1-\Delta(s)]^{-1}=(1+s) /(s-9)\right)$. Finally, the set (5.15) can easily account for the uncertain location of a right half plane pole (as an example: $\Delta(s)=\varepsilon$, $G_{n}(s)=1 /(s-1)$, so that $\left.G(s)=1 /(s-1-\varepsilon)\right)$.

### 5.1.3 Robust design

The design problem in an uncertain environment consists in selecting a controller $K(s)$ which ensures stability as well as satisfactory performances not only in nominal conditions (e.g., $G(s)=G_{n}(s)$ ), but also when the plant undergoes finite perturbations. As for the basic stability requirement, a controller $K(s)$ is said to guarantee robust stability if, given a set $\mathcal{D}_{\alpha}$, the control system is stable for each $G(s) \in \mathcal{G}$. In a similar way, a controller $K(s)$ is said to guarantee robust performances if, given a set $\mathcal{D}_{\alpha}$, the control system satisfies some specified performance requirements (like those defined through eqs. (5.6)-(5.9)) for each $G(s) \in \mathcal{G}$.

Within this framework a natural question arises, namely whether a control system which has been designed in nominal conditions can, for a given set $\mathcal{G}$ and some finite $\alpha$, guarantee robust performances and/or stability relative to the set $\mathcal{D}_{\alpha}$. Whenever possible, the answer to such a question is supplied by the so called procedures for the robustness analysis of a control system.

The same approach adopted for the nominal design problem can be exploited for the robust design problem, provided that reference is made to sensitivity functions defined in terms of $G(s)$ rather than of $G_{n}(s)$, namely

$$
\begin{align*}
S(s) & =[I-G(s) K(s)]^{-1}  \tag{5.16}\\
T(s) & =G(s) K(s)[I-G(s) K(s)]^{-1}  \tag{5.17}\\
V(s) & =K(s)[I-G(s) K(s)]^{-1} \tag{5.18}
\end{align*}
$$

Accordingly, the robust sensitivity performance, the robust complementary sensitivity performance and the robust control sensitivity performance are guaranteed if, given the sets $\mathcal{G}$ and $\mathcal{D}_{\alpha}$, the control system is stable for all $G(s) \in \mathcal{G}$ and

$$
\begin{array}{ll}
\left\|W_{1}(s) S(s)\right\|_{\infty}<1, & \forall G(s) \in \mathcal{G} \\
\left\|W_{2}(s) T(s)\right\|_{\infty}<1, & \forall G(s) \in \mathcal{G} \\
\left\|W_{3}(s) V(s)\right\|_{\infty}<1, & \forall G(s) \in \mathcal{G} \tag{5.21}
\end{array}
$$

respectively.
Not differently from the nominal design framework, it is possible to call for the simultaneous matching of two or even all the inequalities (5.19)-(5.21). Thus, for instance, a design problem could be stated requiring that

$$
\left\|\left[\begin{array}{l}
W_{1}(s) S(s)  \tag{5.22}\\
W_{2}(s) T(s)
\end{array}\right]\right\|_{\infty}<1, \quad \forall G(s) \in \mathcal{G}
$$

Finally, under some circumstances, the controller might be required to guarantee robust stability together with satisfactory performances (as specified by the inequalities presented in this section) in nominal conditions only (nominal design with robust stability).

The next section is devoted to showing how the design problems introduced in this section can all be reduced to a unique standard problem in the $R H_{\infty}$ context which is completely defined in terms of $G_{n}(s)$ and $\mathcal{G}$ only.

### 5.2 The standard problem

Consider the block-structure depicted in fig. 5.3 where $P(s)$ is the so called augmented


Figure 5.3: The standard 2-block configuration
system and $K(s)$ is the controller to be designed. It is now shown that the controller $\mathrm{K}(\mathrm{s})$ in fig. 5.1 which solves one of the design problems defined in Section 5.1 is the same controller that in fig. 5.3 guarantees stability and the boundedness of the $R H_{\infty}$ norm of the transfer function $T(z ; w, s)$ from $w$ to $z$. The interest in reformulating the original design problem in terms of the block structure of fig. 5.3 lies on the fact that the augmented plant $P(s)$ depends only on the nominal plant $G_{n}(s)$, on the particular set $\mathcal{G}$ of the given perturbations and on the performances requested to the control system.

The procedure at the basis of such a reformulation, i.e. the definition of the system $P(s)$ and signals $w$ and $z$, comes up to be very simple when dealing with a design problem in nominal conditions, henceforth referred to as nominal design. Such procedure is now presented at the light of simple, but illustrative, examples, which also serve as preliminaries in the cases not specifically considered herein.

Nominal design: sensitivity performance With reference to fig. 5.1 and fig. 5.3, define $z:=W_{1}(s) y$ and $w:=c^{o}$, or $z:=W_{1}(s) y$ and $w:=d_{r}$, or $z:=W_{1}(s) c$ and $w:=d_{c}$. Then, in the three cases,

$$
P(s)=\left[\begin{array}{cc}
W_{1}(s) & W_{1}(s) G_{n}(s) \\
I & G_{n}(s)
\end{array}\right]
$$

Nominal design: complementary sensitivity performance With reference to fig. 5.1 and fig. 5.3, define $z:=W_{2}(s) c$ and $w:=d_{r}$ or $w:=c^{0}$. Then, in both cases,

$$
P(s)=\left[\begin{array}{cc}
0 & W_{2}(s) G_{n}(s) \\
I & G_{n}(s)
\end{array}\right]
$$

Nominal design: input sensitivity performance With reference to fig. 5.1 and fig. 5.3 , define $z:=W_{3}(s) u_{p}$ and $w:=d_{r}$ or $w:=c^{o}$, or $w:=d_{c}$. Then, in all the three cases,

$$
P(s)=\left[\begin{array}{cc}
0 & W_{3}(s) \\
I & G_{n}(s)
\end{array}\right]
$$

Nominal design: joint sensitivity and complementary sensitivity performance
With reference to fig. 5.1 and fig. 5.3, define $z:=\left[y^{\prime} W_{1}(s)^{\prime} \quad c^{\prime} W_{2}(s)^{\prime}\right]^{\prime}$ and $w:=c^{o}$ or $w:=d_{r}$. Then, in both cases,

$$
P(s)=\left[\begin{array}{cc}
W_{1}(s) & W_{1}(s) G_{n}(s) \\
0 & W_{2}(s) G_{n}(s) \\
I & G_{n}(s)
\end{array}\right]
$$



Figure 5.4: Feedback connection of two systems

Restating in terms of the block-structure of fig. 5.3 the synthesis problems in uncertain conditions requires a preliminary result (usually called small gain theorem) which refers to the feedback connection of fig. 5.4.

Theorem 5.1 Let $G_{1}(s) \in R H_{\infty}$ be an assigned $p \times m$ transfer function and $G_{2}(s) \in$ $R H_{\infty}$ an arbitrary $m \times p$ transfer function with $\left\|G_{2}\right\|_{\infty}<\alpha \neq 0$. Then
i) The feedback connected system of fig. 5.4 is stable for any $G_{2}(s)$ if $\left\|G_{1}(s)\right\|_{\infty} \leq$ $\alpha^{-1}$
ii) If $\left\|G_{1}(s)\right\|_{\infty}>\alpha^{-1}$, there exists a transfer function $G_{2}(s)$ which destabilizes the feedback connected system of fig. 5.4.

Proof If $\left\|G_{1}(s)\right\|_{\infty} \leq \alpha^{-1}$ then, recalling Lemma 2.21, Remark 2.13 and Remark 2.16, it follows that $\bar{\sigma}\left[G_{1}(s) G_{2}(s)\right]<1, \forall R e(s) \geq 0$. Hence, Lemma 2.18 entails that all eigenvalues of $G_{1}(s) G_{2}(s)$ have modulus less than one in the closed right half plane, so that $\operatorname{det}\left[I-G_{1}(s) G_{2}(s)\right] \neq 0, \forall R e(s) \geq 0$. In view of Theorem 3.3 the conclusion is drawn that system in fig. 5.4 is (internally) stable.

On the contrary, suppose that $\left\|G_{1}(s)\right\|_{\infty}=\alpha^{-1}(1+\epsilon):=\rho^{-1}, \epsilon>0$ and consider the case $m \leq p$. Write a singular value decomposition of $G_{1}(j \omega)$ as $G_{1}(j \omega)=$ $U(j \omega) \Sigma(j \omega) V^{\sim}(j \omega)$ with $\Sigma(j \omega)=\left[S(j \omega)^{\prime} 0\right]^{\prime}$, where $S(j \omega)$ is a square and $m$ dimensional matrix. Letting $G_{2}(j \omega):=\rho V(j \omega) T U^{\sim}(j \omega)$, where $T=\left[\begin{array}{ll}I & 0\end{array}\right]$ is a $m \times p$ dimensional matrix, it follows that

$$
\begin{aligned}
\operatorname{det}\left[I-G_{1}(j \omega) G_{2}(j \omega)\right] & =\operatorname{det}\left[I-\rho U(j \omega)\left[\begin{array}{cc}
S(j \omega) & 0 \\
0 & 0
\end{array}\right] U^{\sim}(j \omega)\right] \\
& =\operatorname{det}\left[\left[\begin{array}{cc}
I-\rho S(j \omega) & 0 \\
0 & I
\end{array}\right]\right] \\
& =\operatorname{det}[I-\rho S(j \omega)]
\end{aligned}
$$

Since $\left\|G_{1}(s)\right\|_{\infty}=\rho^{-1}$, there exists a frequency $\bar{\omega}$ such that $\lim _{\omega \rightarrow \bar{\omega}} \bar{\sigma}\left[G_{1}(j \omega)\right]=\rho^{-1}$. Then, recalling that $S(j \omega)=\operatorname{diag}\left\{\sigma_{i}\left[G_{1}(j \omega)\right]\right\}$, it follows that one of the nonzero entries of $\rho S(j \omega)$ tends to one so that $\lim _{\omega \rightarrow \bar{\omega}} \operatorname{det}\left[I-G_{1}(j \omega) G_{2}(j \omega)\right]=0$. In view of Theorem 3.3 the system in fig. 5.4 is unstable. The case $m \geq p$ can be dealt with in a similar way.

Remark 5.1 (Stability of an uncertain matrix) Theorem 5.1 allows one to give a simple answer to a peculiar problem which concerns the stability of an uncertain linear system. Consider a system described by

$$
\dot{x}=\left(A+\Delta_{A}\right) x
$$



Figure 5.5: An uncertain system in feedback configuration
where $A$ is stable and $\Delta_{A}$ represents the uncertainty which is assumed to belong to a set of perturbations specified by

$$
\Delta_{A}=B \Omega C, \quad\|\Omega\| \leq \gamma^{-1}
$$

where $B$ and $C$ are two assigned matrices and $\gamma$ is an assigned positive scalar. Notice that it is possible, through a suitable choice of matrices $B$ and $C$, to effectively describe various situations, for instance the case in which only a parameter of the system dynamic matrix is really uncertain. Looking at fig. 5.5 it is apparent that the uncertain system can be viewed as a closed loop system obtained by performing the control law $u=\Omega y$ on the system with transfer function $G(s)=C(s I-A)^{-1} B$. Then, it can be concluded that, if $\|G(s)\|_{\infty}<\gamma$, the stability is guaranteed for any $\Omega$ such that $\|\Omega\| \leq \gamma^{-1}$.

Thanks to Theorem 2.14, the condition $\|G(s)\|_{\infty}<\gamma$ is equivalent to the existence of a symmetric positive semidefinite stabilizing solution $S_{s}$ of the Riccati equation (in the unknown $S$ )

$$
0=S A+A^{\prime} S+\gamma^{-2} S B B^{\prime} S+C^{\prime} C
$$

Remark 5.2 (Covariance bound) Consider the stochastic system described by

$$
\dot{x}=\left(A+\Delta_{A}\right) x+B_{1} w
$$

where $A$ is stable, $w$ is a zero-mean white noise with identity intensity. Similarly to Remark 5.1, the perturbation $\Delta_{A}$ is assumed here to be described by

$$
\Delta_{A}=B_{2} \Omega_{2} C \quad, \quad\left\|\Omega_{2}\right\| \leq \gamma^{-1}
$$

where $B_{2}$ and $C$ are specified matrices and $\gamma$ a positive scalar. It is also assumed that the system is stable for any perturbation $\Delta_{A}$ of the given form. Hence, it makes sense to tackle the problem of finding a meaningful (i.e. "small") upper bound of the asymptotic covariance matrix $X_{a}\left(\Omega_{2}\right)$ of the system state.

Let $\beta>0$ be fixed and assume for the moment that there exists a symmetric positive semidefinite and stabilizing solution $P_{s}(\beta)$ to the algebraic Riccati equation (in the unknown $P$ )

$$
0=P A^{\prime}+A P+\gamma^{-2} P C^{\prime} C P+B_{2} B_{2}^{\prime}+\beta^{-2} B_{1} B_{1}^{\prime}
$$

It is now proved that $\beta^{2} P_{s}(\beta)$ is a solution of the problem stated above. Recall that the asymptotic state covariance $X_{a}\left(\Omega_{2}\right)$ is the (unique) solution of the Lyapunov equation (in the unknown $X$ )

$$
0=\left(A+B_{2} \Omega_{2} C\right) X+X\left(A+B_{2} \Omega_{2} C\right)^{\prime}+B_{1} B_{1}^{\prime}
$$

This fact can be readily verified by exploiting arguments similar to those used in Remark 2.20 for the computation of $J_{2}$.


Figure 5.6: The covariance bound problem

Now, let $S\left(\Omega_{2}, \beta\right):=\beta^{2} P_{s}(\beta)-X_{a}\left(\Omega_{2}\right)$, and subtract the Lyapunov equation from the Riccati one. Then, $S\left(\Omega_{2}, \beta\right)$ satisfies the following equation (in the unknown $S$ )

$$
\begin{aligned}
0= & \left(A+B_{2} \Omega_{2} C\right) S+S\left(A+B_{2} \Omega_{2} C\right)^{\prime}+ \\
& +\left(\beta \gamma^{-1} P_{s} C^{\prime}-\gamma \beta B_{2} \Omega_{2}\right)\left(\beta \gamma^{-1} P_{s} C^{\prime}-\gamma \beta B_{2} \Omega_{2}\right)^{\prime}+ \\
& +\beta^{2} B_{2}\left(I-\gamma \Omega_{2} \Omega_{2}^{\prime}\right) B_{2}^{\prime}
\end{aligned}
$$

Being ( $A+B_{2} \Omega_{2} C$ ) stable, Lemma C. 1 leads to the conclusion that $S\left(\Omega_{2}, \beta\right) \geq 0$, i.e.

$$
\beta^{2} P_{s}(\beta) \geq X_{a}\left(\Omega_{2}\right), \quad \forall \Omega_{2},\left\|\Omega_{2}\right\|_{\infty} \leq \gamma^{-1}
$$

This inequality raises the interest in determining the value of $\beta$ to which the "minimum" value of the bound of $X_{a}\left(\Omega_{2}\right)$ corresponds. Let $\mathcal{B}$ be the set of $\beta$ 's for which the solution $P_{s}(\beta)$ of the Riccati equation actually exists. Observe that the existence of a solution $P_{s}(\beta)$ entails that the system shown in fig. 5.6 is stable for any $\Omega:=\left[\Omega_{1}^{\prime} \Omega_{2}^{\prime}\right]^{\prime},\|\Omega\| \leq \gamma^{-1}$. This is simply verified in view of Theorems 5.1, 2.13 and Remark 2.22 since $\|F(s, \beta)\|_{\infty}<\gamma$, with $F(s, \beta):=C(s I-A)^{-1} B(\beta)$, and $B(\beta):=\left[\beta^{-1} B_{1} B_{2}\right]$. Hence, $\forall \beta \in \mathcal{B}$ it results $\|F(s, \beta)\|_{\infty}<\gamma$. In order to show that $\mathcal{B} \neq \emptyset$, it is sufficient to observe that $\|F(s, \beta)\|_{\infty}$ is a monotonic nonincreasing function of $\beta$ and

$$
\lim _{\beta \rightarrow 0}\|F(s, \beta)\|_{\infty}=\infty, \quad \lim _{\beta \rightarrow \infty}\|F(s, \beta)\|_{\infty}=\left\|C(s I-A)^{-1} B_{2}\right\|_{\infty}<\gamma
$$

where the last inequality follows from the assumed stability of $A+B_{2} \Omega_{2} C,\left\|\Omega_{2}\right\| \leq \gamma^{-1}$.
In conclusion, there exists $\beta_{1}$ such that $\left\|F\left(s, \beta_{1}\right)\right\|_{\infty}=\gamma$ and $\mathcal{B}=\left(\beta_{1}, \infty\right)$. The problem of finding a meaningfully optimal value of $\beta$ (namely, a $\beta^{\circ}$ to which a "small" bound $\beta^{o^{2}} P_{s}\left(\beta^{\circ}\right)$ of $X_{a}\left(\Omega_{2}\right)$ corresponds) can easily be faced since it can be proven that for $\beta \in \mathcal{B}$

$$
\frac{d^{2}}{d \beta^{2}}\left[\beta^{2} P_{s}(\beta)\right] \geq 0
$$

This entails that, for instance, $\operatorname{trace}\left[\beta^{2} P_{s}(\beta)\right]$ is a convex function of $\beta$ so that the problem

$$
\inf _{\beta \in \mathcal{B}} \operatorname{trace}\left[\beta^{2} P_{s}(\beta)\right]
$$

actually admits a solution.


Figure 5.7: A control system with additive perturbations


Figure 5.8: The standard 3-block configuration

Uncertain design: robust stability Assume that one of the parametrizations of $\mathcal{G}$ in eq. (5.11)-(5.15) has been adopted to account for the uncertain knowledge of the plant. It is then obvious how important would be the choice of a controller $K(s)$ which guarantees closed loop stability for an assigned set of perturbations, i.e. for an assigned value of the scalar variable $\alpha$ (recall the definition of $\mathcal{D}_{\alpha}$, eq. (5.10)). Along these lines the control problem that spontaneously arises is a robust stability design problem. Such a problem is easily reformulated in that of determining (if any) a controller $K(s)$ that, with reference to the scheme of fig. 5.3 , guarantees stability and is such that the $R H_{\infty}$ norm of the transfer function $T(z, w ; s)$ is less than a suitable scalar $\beta$. The augmented plant $P(s)$ depends on the choice of the particular set $\mathcal{G}$, as now shown for the sets in eqs. (5.11) and (5.13). Dealing with the remaining cases is simple at the light of these examples. Considering the set $\mathcal{G}$ given by eq. (5.11) (additive perturbations) is equivalent to considering the control system of fig. 5.7. Then, define $z:=\vartheta$ and $w:=\varphi$ and observe that the control system of fig. 5.8 is completely equivalent to that of fig. 5.7 if

$$
P(s)=\left[\begin{array}{cc}
0 & I \\
I & G_{n}(s)
\end{array}\right]
$$

In view of Theorem 5.1, $K(s)$ guarantees stability for any $\Delta(s) \in \mathcal{D}_{\alpha}$ if and only if $\|T(z, w ; s)\|_{\infty}<\beta:=\alpha^{-1}$ in the system of fig. 5.3.

Analogously, consider the set $\mathcal{G}$ given by eq. (5.13) (multiplicative perturbations) and observe that this corresponds to considering the scheme in fig. 5.9 , which, in


Figure 5.9: A control system with multiplicative perturbations


Figure 5.10: A less primitive description of uncertainty
turn, is completely equivalent to that of fig. 5.8 if $z:=\vartheta, w:=\varphi$ and

$$
P(s)=\left[\begin{array}{cc}
0 & G_{n}(s) \\
I & G_{n}(s)
\end{array}\right]
$$

Hence, stability for any $\Delta(s) \in \mathcal{D}_{\alpha}$ is guaranteed if and only if, in the system of fig. $5.3,\|T(z, w ; s)\|_{\infty}<\beta:=\alpha^{-1}$.

Remark 5.3 At the light of what has been now shown, a less primitive class of uncertainty than the one defined in eq. (5.10) can readily be introduced by substituting the set $\mathcal{D}_{\alpha}$ with

$$
\begin{equation*}
\mathcal{D}_{\alpha}^{*}:=\left\{\Delta(s) \mid \Delta(s)=W_{5}(s) \Delta^{*}(s) W_{4}(s), \quad \Delta^{*}(s) \in R H_{\infty},\left\|\Delta^{*}(s)\right\|_{\infty}<\alpha\right\} \tag{5.23}
\end{equation*}
$$

where $W_{4}(s)$ and $W_{5}(s)$ are assigned elements of $R H_{\infty}$. Without any loss of generality, one can assume $\left\|W_{5}(s)\right\|_{\infty}=1$.

The two shaping functions $W_{4}(s)$ and $W_{5}(s)$ allows one to more finely qualify the perturbations which affect the control loop, in terms of both their harmonic components (nature of the functions) and their very structure (configuration of the matrices). With reference to the set $\mathcal{D}_{\alpha}^{*}$ given in eq. (5.23), it is possible to substitute $\Delta^{*}(s)$ for $\Delta(s)$ in fig. 5.8 , provided that $z:=W_{4}(s) \vartheta$ and $w:=W_{5}^{-1}(s) \varphi$ (recall fig. 5.7, 5.9, 5.10).

In conclusion, the controller $K(s)$ is stabilizing for any $\Delta(s) \in \mathcal{D}_{\alpha}^{*}$ if and only if $\|T(z, w ; s)\|_{\infty}<\beta:=\alpha^{-1}$ in the system of fig. 5.3 , where (recall fig. 5.7, 5.9) the augmented plant $P(s)$ is given, in the two considered cases, by

$$
P(s)=\left[\begin{array}{cc}
0 & W_{4}(s) \\
W_{5}(s) & G_{n}(s)
\end{array}\right]
$$

and

$$
P(s)=\left[\begin{array}{cc}
0 & W_{4}(s) G_{n}(s) \\
W_{5}(s) & G_{n}(s)
\end{array}\right]
$$



Figure 5.11: The control system with multiplicative perturbations
respectively.
Remark 5.4 Assume that one of the (single performance) nominal design problem previously introduced has been performed and let $K(s)$ be the resulting controller. In view of Theorem 5.1 such a controller also guarantees the robust stability relatively to suitable sets $\mathcal{G}$ (of the type of those defined in eqs. (5.11)-(5.15)) and $\mathcal{D}_{\alpha}^{*}$ (of the type of that defined in eq. (5.23)). Actually, the controller that satisfies the constraint $\|T(z, w ; s)\|_{\infty}<\beta$ in fig. 5.3 also guarantees the stability of the control system for any perturbation $\Delta^{*}(s),\left\|\Delta^{*}(s)\right\|_{\infty}<\beta^{-1}$ acting between the signals $z$ and $w$. Hence stability is ensured for any $\Delta(s) \in \mathcal{D}_{\alpha}^{*}, \alpha<\beta^{-1}$ with $W_{5}(s)=I$ and $W_{4}(s)$ depending on the shaping function which is possibly introduced in the formulation of the performance objective (recall eqs. (5.6)-(5.8)).

As an example, consider a controller that solves the nominal design problem in terms of the input sensitivity performance. It also guarantees the robust stability with respect to the perturbations in the set (5.11) with $W_{4}(s)=W_{3}(s)$. On the contrary, a controller that solves the nominal design problem for the sensitivity (complementary sensitivity) performance guarantees the robust stability with respect to the perturbations in the set (5.14) (resp. (5.13)) with $W_{4}(s)=W_{1}(s)\left(\operatorname{resp} . W_{4}(s)=W_{2}(s)\right)$.

Remark 5.5 If the shaping function $W_{4}(s)$ is suitably scaled, one can assume $\alpha=1$ in eq. (5.23) so that for given $W_{4}(s)$ and $W_{5}(s),\left\|W_{5}(s)\right\|_{\infty}=1$, the reference set is

$$
\begin{equation*}
\mathcal{D}^{*}:=\left\{\Delta(s) \mid \Delta(s)=W_{5}(s) \Delta^{*}(s) W_{4}(s), \quad \Delta^{*}(s) \in R H_{\infty},\left\|\Delta^{*}(s)\right\|_{\infty}<1\right\} \tag{5.24}
\end{equation*}
$$

Uncertain design: robust stability and nominal performances Consider the control system depicted in fig. 5.1 and assume that a description of the uncertainty, based on two sets $\mathcal{G}$ and $\mathcal{D}^{*}$, has been selected (recall eqs. (5.11)-(5.15) and (5.24)). The design problem is that of determining a controller $K(s)$ which guarantees the robust stability and the fulfillment, in nominal conditions, of preassigned performance requirements (of the type of those specified in eqs. (5.6)-(5.9)).

For example, if one is interested in the sensitivity performance (eq. (5.6)), being the set $\mathcal{G}$ specified by eq. (5.13), the problem is stated by requiring the determination of $K(s)$ such that (see fig. 5.11) $\left\|W_{1}(s) S_{n}(s)\right\|_{\infty}<1$ and $\left\|W_{4}(s) T_{n}(s) W_{5}(s)\right\|_{\infty}<1$. Recalling now what has been said about eq. (5.9), Lemma 2.21 and the fact that $\left\|W_{5}(s)\right\|_{\infty}=1$, it is apparent that the design specifications are met with if

$$
\left\|\left[\begin{array}{l}
W_{1}(s) S_{n}(s) \\
W_{4}(s) T_{n}(s)
\end{array}\right]\right\|_{\infty}<1
$$



Figure 5.12: The control system with additive perturbations
i.e. if $\|T(z, w ; s)\|_{\infty}<1$ in the block-scheme of fig. 5.3 , with

$$
z:=\left[\begin{array}{c}
z_{1} \\
z_{2}
\end{array}\right], P(s):=\left[\begin{array}{cc}
W_{1}(s) & W_{1}(s) G_{n}(s) \\
0 & W_{4}(s) G_{n}(s) \\
I & G_{n}(s)
\end{array}\right]
$$

On the contrary, if the set $\mathcal{G}$ is the one defined in eq. (5.11), the sensitivity performance requirement is met with if (see fig. 5.12)

$$
\left\|\left[\begin{array}{l}
W_{1}(s) S_{n}(s) \\
W_{4}(s) V_{n}(s)
\end{array}\right]\right\|_{\infty}<1
$$

i.e. if $\|T(z, w ; s)\|_{\infty}<1$ in the block-scheme of fig. 5.3, with

$$
z:=\left[\begin{array}{l}
z_{1} \\
z_{2}
\end{array}\right], \quad P(s):=\left[\begin{array}{cc}
W_{1}(s) & W_{1}(s) G_{n}(s) \\
0 & W_{4}(s) \\
I & G_{n}(s)
\end{array}\right]
$$

The statement of the design problems aimed at contemporarily achieving robust stability and robust sensitivity performances requires the result presented in the following lemma.

Lemma 5.1 Let $X(s)$ and $Y(s)$ be assigned elements of $R L_{\infty}$ and $\Psi(s)$ a generic element of $R L_{\infty}$ such that $\|\Psi(s)\|_{\infty}<1$. If

$$
\begin{equation*}
\sup _{\omega}\{\|X(j \omega)\|+\|Y(j \omega)\|\}<1 \tag{5.25}
\end{equation*}
$$

then

$$
\begin{gather*}
\sup _{\omega}\|X(j \omega)\|<1  \tag{5.26}\\
\sup _{\omega}\left\|Y(j \omega)[I-\Psi(j \omega) X(j \omega)]^{-1}\right\|<1 \tag{5.27}
\end{gather*}
$$

Proof Eq. (5.26) derives directly from eq. (5.25). If $X(s)=0$, eq. (5.27) follows trivially from eq. (5.25). Hence, let $X(s) \neq 0$. Remark 2.13 and Lemmas 2.21 and 2.18 (point 3), entail that, for each $\omega$,

$$
\begin{align*}
\left\|Y(j \omega)[I-\Psi(j \omega) X(j \omega)]^{-1}\right\| & \leq\|Y(j \omega)\| \cdot\left\|[I-\Psi(j \omega) X(j \omega)]^{-1}\right\| \\
& =\frac{\|Y(j \omega)\|}{\underline{\sigma}[I-\Psi(j \omega) X(j \omega)]} \tag{5.28}
\end{align*}
$$

But, for each $\omega$,

$$
\begin{align*}
\underline{\sigma}[I-\Psi(j \omega) X(j \omega)] & \geq \underline{\sigma}(I)-\bar{\sigma}[-\Psi(j \omega) X(j \omega)] \\
& \geq 1-\bar{\sigma}[-\Psi(j \omega) X(j \omega)] \\
& \geq 1-\bar{\sigma}[\Psi(j \omega)] \bar{\sigma}[X(j \omega)] \\
& \geq 1-\bar{\sigma}[X(j \omega)] \\
& \geq 1-\|X(j \omega)\| \tag{5.29}
\end{align*}
$$

The above relations have been written by exploiting Lemma 2.22 (first inequality), Lemma 2.21 (second inequality) and the assumption on the norm of $\Psi(j \omega)$ (last inequality). From eqs. (5.28) and (5.29) it follows that, for each $\omega$,

$$
\left\|Y(j \omega)[I-\Psi(j \omega) X(j \omega)]^{-1}\right\| \leq \frac{\|Y(j \omega)\|}{1-\|X(j \omega)\|}
$$

Thanks to eq. (5.25), the right hand side of this inequality is less than one, so that eq. (5.27) is proven.

Uncertain design: robust sensitivity performances First consider the case where the uncertainty is described by means of the sets $\mathcal{G}$ and $\mathcal{D}^{*}$ defined in eqs. (5.13) and (5.24), respectively (notice also fig. 5.11). Then,

$$
\begin{aligned}
S(s)= & {[I-G(s) K(s)]^{-1} } \\
= & {\left[I-\left(I+W_{5}(s) \Delta^{*}(s) W_{4}(s)\right) G_{n}(s) K(s)\right]^{-1} } \\
= & {\left[I-G_{n}(s) K(s)-W_{5}(s) \Delta^{*}(s) W_{4}(s) G_{n}(s) K(s)\right]^{-1} } \\
= & \left\{I I-W_{5}(s) \Delta^{*}(s) W_{4}(s) G_{n}(s) K(s)\left(I-G_{n}(s) K(s)\right)^{-1}\right] \\
& \left.\cdot\left[I-G_{n}(s) K(s)\right]\right\}^{-1} \\
= & S_{n}(s)\left[I-W_{5}(s) \Delta^{*}(s) W_{4}(s) T_{n}(s)\right]^{-1}
\end{aligned}
$$

Thus the robust sensitivity performance can be expressed as

$$
\begin{gather*}
\left\|W_{1}(s) S_{n}(s)\left[I-W_{5}(s) \Delta^{*}(s) W_{4}(s) T_{n}(s)\right]^{-1}\right\|_{\infty}<1 \\
\forall \Delta^{*}(s) \in R H_{\infty},\left\|\Delta^{*}(s)\right\|_{\infty}<1 \tag{5.30}
\end{gather*}
$$

whereas the robust stability requirement is stated as

$$
\begin{equation*}
\left\|W_{4}(s) T_{n}(s) W_{5}(s)\right\|_{\infty}<1 \tag{5.31}
\end{equation*}
$$

It is easy to verify that eqs. (5.30) and (5.31) are satisfied if

$$
\left\|\left[\begin{array}{l}
W_{1}(s) S_{n}(s)  \tag{5.32}\\
W_{4}(s) T_{n}(s)
\end{array}\right]\right\|_{\infty}<\frac{\sqrt{2}}{2}
$$

Actually, letting $X(s):=W_{4}(s) T_{n}(s) W_{5}(s)$ and $Y(s):=W_{1}(s) S_{n}(s)$, eq. (5.32) implies that (recall that $\left\|W_{5}(s)\right\|_{\infty}=1$ )

$$
\left\|\left[\begin{array}{c}
Y(s) \\
X(s)
\end{array}\right]\right\|_{\infty}^{2}<\frac{1}{2}
$$



Figure 5.13: Sensitivity and complementary sensitivity constraints
and this is equivalent to

$$
\sup _{\omega}\left\{\|X(j \omega)\|^{2}+\|Y(j \omega)\|^{2}\right\}<\frac{1}{2}
$$

This last relation implies (see fig. 5.13)

$$
\sup _{\omega}\{\|X(j \omega)\|+\|Y(j \omega)\|\}<1
$$

so that, thanks to Lemma 5.1, one can conclude about the correctness of what is claimed above. Notice that the simultaneous request of both robust stability and robust performance has lowered, not really surprisingly, the bound on the value of the norm.

Let now consider the case where the uncertainty is described by means of the sets $\mathcal{G}$ and $\mathcal{D}^{*}$ defined in eqs. (5.11) and (5.24), respectively (consider also fig. 5.12). Then, in strict analogy to what has been done in the previous case,

$$
S(s)=S_{n}(s)\left[I-W_{5}(s) \Delta^{*}(s) W_{4}(s) V_{n}(s)\right]^{-1}
$$

Thus, the robust sensitivity performance can be expressed in the following way

$$
\begin{gathered}
\left\|W_{1}(s) S_{n}(s)\left[I-W_{5}(s) \Delta^{*}(s) W_{4}(s) V_{n}(s)\right]^{-1}\right\|_{\infty}<1 \\
\forall \Delta^{*}(s) \in R H_{\infty},\left\|\Delta^{*}(s)\right\|_{\infty}<1
\end{gathered}
$$

whereas the robust stability requirement is formulated as

$$
\left\|W_{4}(s) V_{n}(s) W_{5}(s)\right\|_{\infty}<1
$$

The previously presented arguments can be exploited again to conclude that the goal of the design problem is achieved if the controller $K(s)$ is such that

$$
\left\|\left[\begin{array}{l}
W_{1}(s) S_{n}(s) \\
W_{4}(s) V_{n}(s)
\end{array}\right]\right\|_{\infty}<\frac{\sqrt{2}}{2}
$$

The foregoing discussion has shown that a number of meaningful control problems can be treated in a unified fashion. As a matter of fact, it has been shown that they
are all amenable to the problem of synthesizing a controller $K(s)$ which stabilizes the control system in fig. 5.3 and is such that the $R H_{\infty}$ norm of the transfer function $T(z, w ; s)$ from the input $w$ to the output $z$ is less than one. However, if such objective is not achievable, it makes sense to look for a controller that, besides stabilizing the control system, is such that $\|T(z, w ; s)\|_{\infty}$ is minimized.

The so resulting design problem is apparently well posed only if the augmented plant $P(s)$ is stabilizable: given this, its precise statement is presented in the remaining part of this section, whose development closely follows the one presented in Section 4.1 within the $\mathrm{RH}_{2}$ framework.

First of all, it is convenient to introduce the following time domain description of the process under control (augmented plant), which will be frequently quoted in the development of the present chapter.

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{5.33}\\
z & =C_{1} x+D_{11} w+D_{12} u  \tag{5.34}\\
y & =C_{2} x+D_{21} w+D_{22} u \tag{5.35}
\end{align*}
$$

The controller is constrained to be a finite dimensional, time invariant, linear system, described by:

$$
\begin{align*}
\dot{\xi} & =F \xi+G y  \tag{5.36}\\
u & =H \xi+E y \tag{5.37}
\end{align*}
$$

Hence,

$$
P(s)=\left[\begin{array}{c|cc}
A & B_{1} & B_{2} \\
\hline C_{1} & D_{11} & D_{12} \\
C_{2} & D_{21} & D_{22}
\end{array}\right] \quad, \quad K(s)=\left[\begin{array}{c|c}
F & G \\
\hline H & E
\end{array}\right]
$$

Of course, the feedback connection of system (5.33)-(5.35) with system (5.36),(5.37) must be well defined. For such a condition to be verified it is necessary that

$$
\begin{equation*}
\operatorname{det}\left[I-E D_{22}\right] \neq 0 \tag{5.38}
\end{equation*}
$$

so that the algebraic loop which is created by the insertion of the controller is automatically solvable. With reference to the block structure of fig. 5.3, the primary scope is that of determining $K(s)$ in such a way that the $R H_{\infty}$ norm of the transfer function $T(z, w ; s)$ is less than a specified positive value $\gamma$.

Notice that $T(z, w ; s) \in R H_{\infty}$ entails that such a function is stable: this objective is obviously satisfied if the internal stability of the closed loop system is ensured, i.e. if $K(s)$ in (5.36),(5.37) internally stabilizes system (5.33)-(5.35). This is equivalent to requiring the stability of the dynamic matrix of the closed-loop system, i.e.

$$
\begin{equation*}
 \tag{5.39}
\end{equation*}
$$

In view of the above considerations, it is useful to precisely formalize the concept of admissible controller.

Definition 5.1 (Admissible controller in $R H_{\infty}$ ) A controller $K(s)$ is said to be admissible in $R H_{\infty}$ for $P(s)$ if conditions (5.38)-(5.40,) are verified.


Figure 5.14: The full information problem

In Sections 5.3, 5.5 the attention will be focused on three main problems, each of them associated with a particular structure of the system $P(s)$ : they are referred to as the full information problem, the output estimation problem and the partial information problem. More in detail, the last problem will be tackled by exploiting the solutions of the former ones, which, in turn, are strictly related each to the other by structural relations: the complete picture put into sharp relief important duality and separation properties.

The main result relevant to these problems will concern the solution of two precise points, following the scheme formally presented in Problem 5.1 below: the existence of a controller such that $\|T(z, w ; s)\|_{\infty}<\gamma$ and the parametrization of such controllers. Problem 5.1 refers to the feedback configuration of fig. 5.3 and to the set $\mathcal{F}_{\infty \gamma}$ which represents the family of all admissible controllers in $R H_{\infty}$ for $P(s)$ such that $\|T(z, w ; s)\|_{\infty}<\gamma$.

Problem 5.1 (Standard problem in $R H_{\infty}$ ) Let a positive scalar $\gamma$ be fixed.
a) Find a necessary and sufficient condition for the existence of a controller $K(s)$ which is admissible in $R H_{\infty}$ for $P(s)$ and such that $\|T(z, w ; s)\|_{\infty}<\gamma$.
b) Find a family of controllers $\mathcal{F}_{\infty \gamma r} \subseteq \mathcal{F}_{\infty \gamma}$ whose elements generate the whole set of functions $T(z, w ; s)$ which are generated by the elements of $\mathcal{F}_{\infty \gamma}$.

Remark 5.6 An obvious necessary condition for the existence of a stabilizing controller (and therefore for the existence of an admissible controller in $R H_{\infty}$ for $P(s)$ ) is the stabilizability of the pair $\left(A, B_{2}\right)$ and the detectability of the pair $\left(A, C_{2}\right)$. The statement of Problem 5.1 makes sense only if both properties actually hold true.

In the forthcoming sections the parametrization of the controllers in the family $\mathcal{F}_{\infty \gamma}$ for all the considered cases will be presented. Such parametrization will be the subject of specified remarks which follow the main theorems concerning the solution of Problem 5.1. In other words, the issues relative to the family $\mathcal{F}_{\infty \gamma \gamma}$ and those relative to Problem 5.1 are treated separately. Actually, a unified treatment would further weight down the solution of Problem 5.1.

### 5.3 The full information problem

As in Section 4.2, the output signal $y$ is constituted by the state variable $x$ and the disturbance vector $w$ relevant to the controlled plant. Therefore the plant is described by the equations


Figure 5.15: The set $\mathcal{F}_{\infty \gamma r}$ of the controllers $K_{F I r}(s)$

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{5.41}\\
z & =C_{1} x+D_{12} u  \tag{5.42}\\
y & =\left[y_{1}^{\prime} y_{2}^{\prime}\right]^{\prime}  \tag{5.43}\\
y_{1} & =x  \tag{5.44}\\
y_{2} & =w \tag{5.45}
\end{align*}
$$

while the transfer function from $\left[\begin{array}{ll}w^{\prime} & u^{\prime}\end{array}\right]^{\prime}$ to $\left[\begin{array}{ll}z^{\prime} & y^{\prime}\end{array}\right]^{\prime}$ will be denoted with $P_{F I}(s)$. A complete answer to Problem 5.1 can be given under the following assumptions.

Assumption 5.1 The pair $\left[\left(A-B_{2} D_{12}^{\prime} C_{1}\right),\left(I-D_{12} D_{12}^{\prime}\right) C_{1}\right]$ is detectable and the pair $\left(A, B_{2}\right)$ is stabilizable.

Assumption 5.2 $D_{12}^{\prime} D_{12}=I$.
The main result of this section makes reference to the block-diagram of fig. 5.14, where $K_{F I}(s)$ is any controller $R H_{\infty}$ admissible for $P_{F I}(s)$ (recall Definition 5.1). In this diagram the transfer function from $w$ to $z$ is, as usual, denoted by $T(z, w ; s)$.

Theorem 5.2 (Full information) Consider Problem 5.1 relative to system (5.41)(5.45). Then, under Assumptions 5.1, 5.2, it has the solution
a) The existence of a symmetric, positive semidefinite and stabilizing solution $P_{\infty}$ of the Riccati equation (in the unknown $P$ )

$$
\begin{equation*}
0=P A_{c}+A_{c}^{\prime} P-P\left(B_{2} B_{2}^{\prime}-\gamma^{-2} B_{1} B_{1}^{\prime}\right) P+C_{1 c}^{\prime} C_{1 c} \tag{5.46}
\end{equation*}
$$

i.e., such that the matrix $A_{c c}$ given by

$$
\begin{equation*}
A_{c c}:=A_{c}-B_{2} B_{2}^{\prime} P_{\infty}+\gamma^{-2} B_{1} B_{1}^{\prime} P_{\infty} \tag{5.47}
\end{equation*}
$$

is stable. In eqs. (5.46), (5.47)

$$
\begin{equation*}
A_{c}:=A-B_{2} D_{12}^{\prime} C_{1}, \quad C_{1 c}:=\left(I-D_{12} D_{12}^{\prime}\right) C_{1} \tag{5.48}
\end{equation*}
$$

b) The set $\mathcal{F}_{\infty \gamma r}$ of the controllers $K_{F I r}(s)$ is defined by the diagram of fig. 5.15, where

$$
\begin{equation*}
F_{\infty}:=-B_{2}^{\prime} P_{\infty}-D_{12}^{\prime} C_{1} \tag{5.49}
\end{equation*}
$$

and $Q(s):=\Sigma\left(A_{q}, B_{q}, C_{q}, D_{q}\right)$ with $A_{q}$ stable and $\|Q(s)\|_{\infty}<\gamma$.


Figure 5.16: The modified full information problem

Proof First observe that the necessary condition for the problem at hand to make sense (recall Remark 5.6) is satisfied. Indeed, being measurable the state of the system, detectability of the pair $\left(A, C_{2}\right)$ trivially holds, while, on the other hand, stabilizability of the pair $\left(A, B_{2}\right)$ is guaranteed by Assumption 5.1.

Let now

$$
\begin{equation*}
u_{c}:=u+D_{12}^{\prime} C_{1} x \tag{5.50}
\end{equation*}
$$

which simply amounts to defining a control law as shown in fig. 5.16. The resulting system $P_{\text {FIc }}(s)$ is therefore described by

$$
\begin{align*}
\dot{x} & =A_{c} x+B_{1} w+B_{2} u_{c}  \tag{5.51}\\
z & =C_{1 c} x+D_{12} u_{c}  \tag{5.52}\\
y & =\left[y_{1}^{\prime} y_{2}^{\prime}\right]^{\prime}  \tag{5.53}\\
y_{1} & =x  \tag{5.54}\\
y_{2} & =w \tag{5.55}
\end{align*}
$$

Being measurable the state, solving Problem 5.1 relative to system $P_{\text {FIc }}(s)$ is equivalent to solving the same problem relative to system $P_{F I}(s)$. As for point (b) in particular, the solution relative to system $P_{F I}(s)$ follows from that relative to system $P_{F I c}(s)$ by recalling eq. (5.50).

The proof is organized into three main parts. In part (i) it will be shown that the observability of the pair $\left(A_{c}, C_{1 c}\right)$ can be assumed without any loss of generality. Then, in part (ii), the necessity of point a) will be proved. Finally, part (iii) is devoted to proving sufficiency of point $a$ ) and point $b$ ).

Part (i) With reference to fig. 5.16 , denote by $T(z, w ; s)$ the transfer function from $w$ to $z$ and suppose that a controller $K_{F I c}(s)$ admissible in $R H_{\infty}$ exists such that $\|T(z, w ; s)\|_{\infty}<\gamma$.

Now it will be shown that there is no loss of generality in assuming the pair $\left(A_{c}, C_{1 c}\right)$ observable rather than simply detectable. To this aim, let this pair already


Figure 5.17: The canonical decomposition of the modified problem
be in Kalman's canonical observability form, namely

$$
A_{c}=\left[\begin{array}{cc}
A_{c 1} & 0 \\
A_{c 2} & A_{c 3}
\end{array}\right], C_{1 c}=\left[\begin{array}{cc}
C_{1 c 1} & 0
\end{array}\right]
$$

with the pair $\left(A_{c 1}, C_{1 c 1}\right)$ observable, and, accordingly, decompose the matrices $B_{1}$ and $B_{2}$ as

$$
B_{1}=\left[\begin{array}{l}
B_{11} \\
B_{12}
\end{array}\right], B_{2}=\left[\begin{array}{l}
B_{21} \\
B_{22}
\end{array}\right]
$$

Finally, consider the state equations for $K_{F I c}(s)$

$$
\begin{align*}
\dot{\xi} & =F \xi+G_{1} x_{1}+G_{2} x_{2}+G_{3} w  \tag{5.56}\\
u_{c} & =H \xi+E_{1} x_{1}+E_{2} x_{2}+E_{3} w \tag{5.57}
\end{align*}
$$

where $x_{1}$ and $x_{2}$ are the components of the state vector $x=\left[x_{1}^{\prime} x_{2}^{\prime}\right]^{\prime}$ of system $P_{F I c}(s)$ in Kalman's canonical form. The overall system can then be viewed (see also fig. 5.17) as resulting from the feedback connection of the subsystem $P_{F I c 1}(s)$

$$
\begin{align*}
\dot{x}_{1} & =A_{c 1} x_{1}+B_{11} w+B_{21} u_{c}  \tag{5.58}\\
z & =C_{1 c 1} x_{1}+D_{12} u_{c} \tag{5.59}
\end{align*}
$$

with the controller $K_{\text {FIc1 }}(s)$ which, in turn, is constituted by $K_{\text {FIc }}(s)$ (see eqs. $(5.56),(5.57))$ and the subsystem $P_{F I c 2}(s)$

$$
\dot{x}_{2}=A_{c 2} x_{1}+A_{c 3} x_{2}+B_{12} w+B_{22} u_{c}
$$

Indeed, by substituting eq. (5.57) for $u_{c}$, the controller $K_{F I c 1}(s)$ for the system $P_{F I c 1}(s)$ (which is defined by eqs. (5.58),(5.59)) turns out to be described by

$$
\begin{aligned}
\dot{\xi}= & F \xi+G_{2} x_{2}+G_{1} x_{1}+G_{3} w \\
\dot{x}_{2}= & B_{22} H \xi+\left(A_{c 3}+B_{22} E_{2}\right) x_{2}+\left(A_{c 2}+B_{22} E_{1}\right) x_{1} \\
& +\left(B_{12}+B_{22} E_{3}\right) w \\
u_{c}= & H \xi+E_{2} x_{2}+E_{1} x_{1}+E_{3} w
\end{aligned}
$$

Obviously, the controller $K_{F I c 1}(s)$ is admissible in $R H_{\infty}$ for system $P_{F I c 1}(s)$ since $K_{F I c}(s)$ is such for $P_{F I c}(s)$.

Finally, observe that (the simple check is left to the reader) if $P_{1 \infty}$ is the symmetric, positive semidefinite and stabilizing solution of the Riccati equation (in the unknown $P_{1}$ )

$$
\begin{equation*}
0=P_{1} A_{c 1}+A_{c 1}^{\prime} P_{1}-P_{1}\left(B_{21} B_{21}^{\prime}-\gamma^{-2} B_{11} B_{11}^{\prime}\right) P_{1}+C_{1 c 1}^{\prime} C_{1 c 1} \tag{5.60}
\end{equation*}
$$

i.e., such that $A_{c c 1}:=A_{c 1}-\left(B_{21} B_{21}^{\prime}-\gamma^{-2} B_{11} B_{11}^{\prime}\right) P_{1 \infty}$ is stable, then

$$
P_{\infty}:=\left[\begin{array}{cc}
P_{1 \infty} & 0  \tag{5.61}\\
0 & 0
\end{array}\right]
$$

is a symmetric and positive semidefinite solution of eq. (5.46). Moreover, such a solution is stabilizing since the eigenvalues of the matrix $A_{c c}$, defined by eqs. (5.47),(5.48), are those of the matrix $A_{c 3}$ (which is stable by assumption, because the pair ( $A_{c}, C_{1 c}$ ) is detectable) together with those of the matrix $A_{c c 1}$.

Having proved that the existence of a controller $K_{\text {FIc }}(s)$ admissible in $R H_{\infty}$ for system $P_{F I c}(s)$ and such that $\|T(z, w ; s)\|_{\infty}<\gamma$ implies the existence of a controller $K_{F I c 1}(s)$ admissible in $R H_{\infty}$ for system $P_{F I c 1}(s)$ and such that $\|T(z, w ; s)\|_{\infty}<\gamma$, the necessity of the existence of the solution of the Riccati equation (5.46) is proved, thanks to eq. (5.61), once the necessity of the existence of the solution of eq. (5.60) has been ascertained. It is therefore possible to assume the pair ( $A_{c}, C_{1 c}$ ) to be observable from the very beginning.

Part (ii) With reference to system $P_{F I c}(s)$ (recall eqs. (5.51), (5.55)), consider the control law

$$
\begin{equation*}
u_{c}:=v-B_{2}^{\prime} P_{2} x \tag{5.62}
\end{equation*}
$$

where $P_{2}$ is the symmetric, positive semidefinite and stabilizing solution of the Riccati equation (in the unknown $P$ )

$$
\begin{equation*}
0=P A_{c}+A_{c}^{\prime} P-P B_{2} B_{2}^{\prime} P+C_{1 c}^{\prime} C_{1 c} \tag{5.63}
\end{equation*}
$$

Indeed, such a solution exists (recall Lemma C.4) because the pair ( $A_{c}, C_{1 c}$ ) is observable and the pair $\left(A_{c}, B_{2}\right)$ is stabilizable. This latter claim derives from the stabilizability assumption on the pair $\left(A, B_{2}\right)$ which is equivalent to that of the pair $\left(A_{c}, B_{2}\right)$, in view of the form of $A_{c}$. Moreover, $P_{2}>0$ due to the observability property.

From eqs. (5.51),(5.52) and (5.62) it follows

$$
\begin{equation*}
z_{L 0}=P_{c}(s) B_{1} w_{L}+U(s) v_{L} \tag{5.64}
\end{equation*}
$$

where

$$
\begin{aligned}
P_{c}(s) & :=\Sigma\left(A_{c}-B_{2} B_{2}^{\prime} P_{2}, I, C_{1 c}-D_{12} B_{2}^{\prime} P_{2}, 0\right) \\
U(s) & :=\Sigma\left(A_{c}-B_{2} B_{2}^{\prime} P_{2}, B_{2}, C_{1 c}-D_{12} B_{2}^{\prime} P_{2}, D_{12}\right)
\end{aligned}
$$

It is now shown that the quadruple $\left(A_{c}, B_{2}, C_{1 c}, D_{12}\right)$ verifies the assumptions which are required by Lemma C.5.

Indeed, from one side it is $D_{12}^{\prime} D_{12}=I$ because of Assumption 5.2, while, on the other side, from the same assumption and eq. (5.48) it is $C_{1 c}^{\prime} D_{12}=0$. Moreover, $P_{2}$ is the symmetric, positive definite and stabilizing solution of eq. (5.63). Thus, from Lemma C. 5 it follows that system

$$
F(s):=\left[\begin{array}{cc}
U(s) & U^{\perp}(s) \tag{5.65}
\end{array}\right]
$$

is square and inner and it results

$$
B_{1}^{\prime} P_{c}^{\sim}(s) F(s)=\left[\begin{array}{c|cc}
A_{c}-B_{2} B_{2}^{\prime} P_{2} & B_{2} & -P_{2}^{-1} C_{1 c}^{\prime} D_{12}^{\perp}  \tag{5.66}\\
\hline B_{1}^{\prime} P_{2} & 0 & 0
\end{array}\right]
$$

where

$$
U^{\perp}(s):=\Sigma\left(A_{c}-B_{2} B_{2}^{\prime} P_{2},-P_{2}^{-1} C_{1 c}^{\prime} D_{12}^{\perp}, C_{1 c}-D_{12} B_{2}^{\prime} P_{2}, D_{12}^{\perp}\right)
$$

and $D_{12}^{\perp}$ is such that the matrix [ $D_{12} D_{12}^{\perp}$ ] is orthogonal, i.e., such that $D_{12}^{\perp_{2}^{\prime}} D_{12}^{\perp}=I$ and $D_{12}^{\perp \prime} D_{12}=0$. Consider now the vector space

$$
\mathcal{Q}:=\left\{q=\left[\begin{array}{l}
q_{1} \\
q_{2}
\end{array}\right], q_{1} \in R H_{2}^{\perp}, q_{2} \in R L_{2}\right\}
$$

(which is a subspace of $R L_{2}$ ) and the operator $\Xi: \mathcal{Q} \rightarrow R H_{2}$ defined by (recall Definition 2.33)

$$
\Xi:\left[\begin{array}{l}
q_{1}  \tag{5.67}\\
q_{2}
\end{array}\right] \mapsto \Pi_{s} B_{1}^{\prime} P_{c}^{\sim}(s) F(s)\left[\begin{array}{l}
q_{1} \\
q_{2}
\end{array}\right]:=w
$$

Now it will be shown that the operator $\Xi^{*}: R H_{2} \rightarrow \mathcal{Q}$ defined by (again recall Definition 2.33)

$$
\Xi^{*}: w \mapsto\left[\begin{array}{c}
\Pi_{a} U^{\sim}(s)  \tag{5.68}\\
U^{\perp \sim}(s)
\end{array}\right] P_{c}(s) B_{1} w:=q
$$

is, consistently with the adopted symbol, the adjoint of the operator $\Xi$.
In fact, from $<w, \Xi q>=<\Xi^{*} w, q>, \forall q \in \mathcal{Q}, \forall w \in R H_{2}$, it follows that

$$
\begin{aligned}
<w, \Xi q> & =<w, \Pi_{s} B_{1}^{\prime} P_{c}^{\sim}(s) U(s) q_{1}>+<w, \Pi_{s} B_{1}^{\prime} P_{c}^{\sim}(s) U^{\perp}(s) q_{2}> \\
& =<w, B_{1}^{\prime} P_{c}^{\sim}(s) U(s) q_{1}>+<w, B_{1}^{\prime} P_{c}^{\sim}(s) U^{\perp}(s) q_{2}> \\
& =<U^{\sim}(s) P_{c}(s) B_{1} w, q_{1}>+<U^{\perp \sim}(s) P_{c}(s) B_{1} w, q_{2}> \\
& =<\Pi_{a} U^{\sim}(s) P_{c}(s) B_{1} w, q_{1}>+<U^{\perp \sim}(s) P_{c}(s) B_{1} w, q_{2}> \\
& =<\left[\begin{array}{c}
\Pi_{a} U^{\sim}(s) \\
U^{\perp \sim}(s)
\end{array}\right] P_{c}(s) B_{1} w,\left[\begin{array}{c}
q_{1} \\
q_{2}
\end{array}\right]> \\
& =<\Xi^{*} w, q>
\end{aligned}
$$

In the above equations the two identities $<\Pi_{a} \alpha, q_{1}>=<\alpha, q_{1}>$ and $<w, \beta>=<$ $w, \Pi_{s} \beta>$, which hold whenever $q_{1} \in R H_{2}^{\perp}$ and $w \in R H_{2}$, have been exploited now and then. Since $F(s)$ is square and inner, then $F^{\sim}(s) F(s)=F(s) F^{\sim}(s)=I$ and therefore, recalling Definition 2.30,

$$
\begin{equation*}
\|z\|_{2}=\left\|F^{\sim}(s) z_{L}\right\|_{2} \tag{5.69}
\end{equation*}
$$

Now suppose that a controller exists which is admissible in $R H_{\infty}$ for $P_{\text {FIc }}(s)$ and such that $\|T(z, w ; s)\|_{\infty}<\gamma$. Therefore, with reference to the block-scheme of fig. 5.16 , it is

$$
\sup _{\substack{w \in R H_{2} \\\|w\|_{2}=1}}\|z\|_{2}^{2}<\gamma^{2}
$$

from which it follows

$$
\sup _{\substack{w \in R H_{2} \\\|w\|_{2}=1}} \inf _{u_{c} \in R H_{2}}\|z\|_{2}^{2}<\gamma^{2}
$$

As a matter of fact, this inequality is fulfilled when $u_{c}$, rather than being suitably chosen in $\mathrm{RH}_{2}$, is the output of the controller $K_{F I c}(s)$ (observe that such an output belongs to $R H_{2}$, because the controller is admissible and $w \in R H_{2}$ ).

By recalling eq. (5.62) and $x \in R H_{2}$, it then follows

$$
\begin{equation*}
\sup _{\substack{w \in R H_{2} \\\|w\|_{2}=1}} \inf _{v \in R H_{2}}\|z\|_{2}^{2}<\gamma^{2} \tag{5.70}
\end{equation*}
$$

Eqs. (5.70), (5.69),(5.65) and (5.64) imply

$$
\begin{align*}
\gamma^{2} & >\sup _{\substack{w \in R H_{2} \\
\|w\|_{2}=1}} \inf _{v \in R H_{2}}\|z\|_{2}^{2} \\
& >\sup _{\substack{w \in R H_{2} \\
\|w\|_{2}=1}} \inf _{v \in R H_{2}}\left\|\left[\begin{array}{c}
U^{\sim}(s) P_{c}(s) B_{1} w+U^{\sim}(s) U(s) v \\
U^{\perp \sim}(s) P_{c}(s) B_{1} w+U^{\perp \sim}(s) U(s) v
\end{array}\right]\right\|_{2}^{2} \\
& >\sup _{\substack{w \in R H_{2} \\
\|w\|_{2}=1}} \inf _{v \in R H_{2}}\left\|\left[\begin{array}{c}
U^{\sim}(s) P_{c}(s) B_{1} w+v \\
U^{\perp \sim}(s) P_{c}(s) B_{1} w
\end{array}\right]\right\|_{2}^{2} \tag{5.71}
\end{align*}
$$

The identities $U^{\sim}(s) U(s)=I$ and $U^{\perp \sim}(s) U(s)=0$ which are consequences of $F(s)$ being inner, have been taken into account in writing down eq. (5.71). Moreover, by defining

$$
\xi:=v+\Pi_{s} U^{\sim}(s) P_{c}(s) B_{1} w
$$

and noticing that $\xi \in R H_{2}$ and

$$
\Pi_{s} U^{\sim}(s) P_{c}(s) B_{1} w+\Pi_{a} U^{\sim}(s) P_{c}(s) B_{1} w=U^{\sim}(s) P_{c}(s) B_{1} w
$$

from eqs. (5.71) and (5.68) it results

$$
\begin{align*}
\gamma^{2} & >\sup _{\substack{w \in R H_{2} \\
\|w\|_{2}=1}} \inf _{\xi \in R H_{2}}\|z\|_{2}^{2} \\
& >\sup _{\substack{w \in R H_{2} \\
\|w\|_{2}=1}} \inf _{\xi \in R H_{2}}\left\|\left[\begin{array}{c}
\Pi_{a} U^{\sim}(s) P_{c}(s) B_{1} w+\xi \\
U^{\perp \sim}(s) P_{c}(s) B_{1} w
\end{array}\right]\right\|_{2}^{2} \\
& >\sup _{\substack{w \in R H_{2} \\
\|w\|_{2}=1}} \inf _{\xi \in R H_{2}}\left\|\Xi^{*} w+\left[\begin{array}{c}
\xi \\
0
\end{array}\right]\right\|_{2}^{2} \tag{5.72}
\end{align*}
$$

But

$$
\begin{align*}
\left\|\Xi^{*} w+\left[\begin{array}{l}
\xi \\
0
\end{array}\right]\right\|_{2}^{2} & =\left\|\Xi^{*} w\right\|_{2}^{2}+\|\xi\|_{2}^{2}+2<\xi, \Pi_{a} U^{\sim}(s) P_{c}(s) B_{1} w> \\
& =\left\|\Xi^{*} w\right\|_{2}^{2}+\|\xi\|_{2}^{2} \tag{5.73}
\end{align*}
$$

since $\xi \in R H_{2}$, while $\Pi_{a} U^{\sim}(s) P_{c}(s) B_{1} w \in R H_{2}^{\perp}$. Therefore, from eqs. (5.72) and (5.73) it follows

$$
\begin{aligned}
\gamma^{2} & >\sup _{\substack{w \in R H_{2} \\
\|w\|_{2}=1}} \inf _{\xi \in R H_{2}}\|z\|_{2}^{2} \\
& >\sup _{\substack{w \in R H_{2} \\
\|w\|_{2}=1}} \inf _{\xi \in R H_{2}}\left\{\left\|\Xi^{*} w\right\|_{2}^{2}+\|\xi\|_{2}^{2}\right\}
\end{aligned}
$$

$$
\begin{aligned}
& >\sup _{\substack{w \in \mathrm{RH}_{2} \\
\|w\|_{2}=1}}\left\|\Xi^{*} w\right\|_{2}^{2} \\
& >\left\|\Xi^{*}\right\|^{2} \\
& >\|\Xi\|^{2}
\end{aligned}
$$

The last equality sign follows from Theorem 2.9. Then, by recalling Remark 2.13, it is

$$
\begin{equation*}
\sup _{\substack{q \in \mathbb{Q} \\\|q\|_{2}=1}}\|\Xi q\|_{2}<\gamma \tag{5.74}
\end{equation*}
$$

By exploiting Lemma G.2, applied to the operator $\Xi$ defined in eq. (5.67) and to the system specified in eq. (5.66), from eq. (5.74) the conclusion can be drawn that a symmetric, positive semidefinite and stabilizing solution $W$ exists for the Riccati equation in the unknown $P$ (recall that $D_{12}^{\frac{1}{2}} D_{12}^{\perp \prime}=I-D_{12} D_{12}^{\prime}$ )

$$
\begin{align*}
0= & P\left(A_{c}-B_{2} B_{2}^{\prime} P_{2}\right)+\left(A_{c}-B_{2} B_{2}^{\prime} P_{2}\right)^{\prime} P+ \\
& +\gamma^{-2} P P_{2}^{-1} C_{1 c}^{\prime} C_{1 c} P_{2}^{-1} P+P_{2} B_{1} B_{1}^{\prime} P_{2} \tag{5.75}
\end{align*}
$$

Lemma G. 2 also implies that $r_{s}\left(W L_{c}\right)<\gamma^{-2}$, where $L_{c}$ solves the Lyapunov equation (in the unknown $L$ )

$$
\begin{align*}
0= & L\left(A_{c}-B_{2} B_{2}^{\prime} P_{2}\right)^{\prime}+\left(A_{c}-B_{2} B_{2}^{\prime} P_{2}\right) L+ \\
& +P_{2}^{-1} C_{1 c}^{\prime} C_{1 c} P_{2}^{-1}+B_{2} B_{2}^{\prime} \tag{5.76}
\end{align*}
$$

By recalling that $P_{2}$ solves eq. (5.63), it is easy to verify that $P_{2}^{-1}$ satisfies the Lyapunov equation (5.76) so that it actually coincides with $L_{c}$, as such an equation admits a unique solution, thanks to Lemma C. 1 (all the eigenvalues of the matrix $A_{c}-B_{2} B_{2}^{\prime} P_{2}$ have negative real parts). Therefore, $L_{c}=P_{2}^{-1}$ and $r_{s}\left(W P_{2}^{-1}\right)<\gamma^{2}$, that is (recall Lemma B.11)

$$
\begin{equation*}
\gamma^{2} P_{2}-W>0 \tag{5.77}
\end{equation*}
$$

The Hamiltonian matrix $Z_{W}$ associated with the Riccati equation (5.75) is

$$
Z_{W}=\left[\begin{array}{cc}
A_{c}-B_{2} B_{2}^{\prime} P_{2} & \gamma^{-2} P_{2}^{-1} C_{1 c}^{\prime} C_{1 c} P_{2}^{-1} \\
-P_{2} B_{1} B_{1}^{\prime} P_{2} & -\left(A_{c}-B_{2} B_{2}^{\prime} P_{2}\right)^{\prime}
\end{array}\right]
$$

By letting

$$
T:=\left[\begin{array}{cc}
-\gamma^{2} I & P_{2}^{-1} \\
-\gamma^{2} P_{2} & 0
\end{array}\right]
$$

and recalling that $P_{2}$ solves eq. (5.63), it is easy to verify that

$$
Z_{\infty}:=T Z_{W} T^{-1}=\left[\begin{array}{cc}
A_{c} & \gamma^{-2} B_{1} B_{1}^{\prime}-B_{2} B_{2}^{\prime}  \tag{5.78}\\
-C_{1 c}^{\prime} C_{1 c} & -A_{c}^{\prime}
\end{array}\right]
$$

The matrix $Z_{\infty}$ is the Hamiltonian matrix associated with the Riccati equation (5.46). Now, remember that, being $W$ a stabilizing solution of eq. (5.75), it is

$$
Z_{W}=\left[\begin{array}{c}
I \\
W
\end{array}\right]=\left[\begin{array}{c}
I \\
W
\end{array}\right] V
$$

where $V$ is a stable matrix. By taking into account eq. (5.78) it then follows

$$
Z_{\infty} T\left[\begin{array}{c}
I \\
W
\end{array}\right]=T Z_{W}\left[\begin{array}{c}
I \\
W
\end{array}\right]=T\left[\begin{array}{c}
I \\
W
\end{array}\right] V
$$



Figure 5.18: The set $\mathcal{F}_{\infty y c r}$ of the controllers $K_{\text {FIcr }}(s)$
so that $\operatorname{Im}\left[T[I W]^{\prime}\right]$ is the $Z_{\infty}$-invariant subspace associated with the stable eigenvalues of such a matrix. Finally, being

$$
T\left[\begin{array}{c}
I \\
W
\end{array}\right]=\left[\begin{array}{c}
-\gamma^{2} I+P_{2}^{-1} W \\
-\gamma^{2} P_{2}
\end{array}\right]
$$

from Lemma C. 3 it can be concluded that

$$
\gamma^{2} P_{2}\left(\gamma^{2} I-P_{2}^{-1} W\right)^{-1}=\gamma^{2} P_{2}\left(\gamma^{2} P_{2}-W\right)^{-1} P_{2}
$$

is the (unique) symmetric and stabilizing solution $P_{\infty}$ of the Riccati equation (5.46). In view of eq. (5.77) such a solution is positive definite: hence $P_{\infty}$ possesses all the properties listed in the statement of the theorem and the necessity of point $a$ ) is proved.

Part (iii) Here reference is made to fig. 5.16 and the symmetric, positive semidefinite and stabilizing solution $P_{\infty}$ of eq. (5.46) is supposed to exist. Denote with $\mathcal{F}_{\infty \gamma c r}$ the set (described in fig. 5.18) of controllers $K_{\text {FIcr }}(s)$ for the system $P_{\text {FIc }}(s)$. By comparing fig. $5.14,5.15$ with fig. $5.16,5.18$ and recalling eqs. (5.49),(5.50) it is apparent that proving that the controllers $K_{\text {FIr }}(s)$ of the set $\mathcal{F}_{\infty \gamma r}$ are admissible in $R H_{\infty}$ for $P_{F I}(s)$ and $\|T(z, w ; s)\|_{\infty}<\gamma$ is completely equivalent to proving that the controllers $K_{\text {FIcr }}(s)$ of the set $\mathcal{F}_{\infty \gamma c r}$ are admissible in $R H_{\infty}$ for $P_{F I c}(s)$ and $\|T(z, w ; s)\|_{\infty}<\gamma$. Moreover, the set $\mathcal{F}_{\infty \gamma c r}$ solves point (b) of Problem 5.1 for system $P_{F I c}(s)$ if and only if the set $\mathcal{F}_{\infty \gamma r}$ solves point $b$ ) of the same problem for system $P_{F I}(s)$. As far as the sufficiency of point ( $a$ ) and point $b$ ) are concerned, it is therefore correct considering Problem 5.1 relative to the system $P_{\text {FIc }}(s)$ and the set $\mathcal{F}_{\infty \gamma c r}$. Consequently, it will first be shown that if a controller $K_{\text {FIcr }}(s)$ belongs to $\mathcal{F}_{\infty \gamma c r}$, then it is admissible in $R H_{\infty}$ for $P_{F I c}(s)$ and $\|T(z, w ; s)\|_{\infty}<\gamma$ (part (iii.1)). Second, it will be proved that if a controller admissible in $R H_{\infty}$ for $P_{F I c}(s)$ exists and $\|T(z, w ; s)\|_{\infty}<\gamma$, then there exists in the set $\mathcal{F}_{\infty \gamma c r}$ a controller which generates the same transfer function $T(z, w ; s)$ (part (iii.2)).

Part (iii.1) Consider the block-scheme of fig. 5.18. If the system $Q(s)$ with realization $Q(s):=\Sigma\left(A_{q}, B_{q}, C_{q}, D_{q}\right)$ is stable and such that $\|Q(s)\|_{\infty}<\gamma, r:=$ $w-\gamma^{-2} B_{1}^{\prime} P_{\infty} x, q:=u_{c}+B_{2}^{\prime} P_{\infty} x$, then (recall eqs. (5.51)-(5.55)) the resulting system (that is the system of fig. 5.16 with $K_{F I c}(s)$ given by the block-scheme in fig. $5.18)$ with transfer function $T(z, w ; s)$ is described by $T(z, w ; s)=\Sigma\left(A_{z}, B_{z}, C_{z}, D_{z}\right)$,
where

$$
\begin{align*}
A_{z} & =\left[\begin{array}{cc}
A_{c}-B_{2} B_{2}^{\prime} P_{\infty}-\gamma^{-2} B_{2} D_{q} B_{1}^{\prime} P_{\infty} & B_{2} C_{q} \\
-\gamma^{-2} B_{q} B_{1}^{\prime} P_{\infty} & A_{q}
\end{array}\right]  \tag{5.79}\\
B_{z} & =\left[\begin{array}{cc}
B_{1}+B_{2} D_{q} \\
B_{q}
\end{array}\right]  \tag{5.80}\\
C_{z} & =\left[\begin{array}{cc}
C_{1 c}-D_{12} B_{2}^{\prime} P_{\infty}-\gamma^{-2} D_{12} D_{q} B_{1}^{\prime} P_{\infty} & D_{12} C_{q}
\end{array}\right]  \tag{5.81}\\
D_{z} & =D_{12} D_{q} \tag{5.82}
\end{align*}
$$

The system (5.79)-(5.82) coincides with the one referred to in Lemma E. 3 and illustrated in fig. E.5. The conditions under which such a lemma can be applied are verified, so that it can be deduced that matrix $A_{z}$ is stable and $\|T(z, w ; s)\|_{\infty}<\gamma$.

Part(iii.2) Suppose that there exists a controller $K_{F I c}(s)$ which is admissible in $R H_{\infty}$ for $P_{F I c}(s)$ and such that $\|T(z, w ; s)\|_{\infty}<\gamma$. Let this controller be described by the equations

$$
\begin{align*}
\dot{\xi} & =L \xi+M_{1} w+M_{2} x  \tag{5.83}\\
u_{c} & =N \xi+O_{1} w+O_{2} x \tag{5.84}
\end{align*}
$$

Now define the variable $q:=u_{c}+B_{2}^{\prime} P_{\infty} x$ and let $w:=r+\gamma^{-2} B_{1}^{\prime} P_{\infty} x$. By putting together eqs. (5.83),(5.84) with eqs. (5.51)-(5.55) it is straightforward to verify that a possible realization of the transfer function $Q(s)$ from $r$ to $q$ is given by

$$
\begin{align*}
\dot{\vartheta}= & L \vartheta+\left(M_{2}+\gamma^{-2} M_{1} B_{1}^{\prime} P_{\infty}\right) \sigma+M_{1} r  \tag{5.85}\\
\dot{\sigma}= & B_{2} N \vartheta+\left(A_{c}+B_{2} O_{2}+\gamma^{-2} B_{1} B_{1}^{\prime} P_{\infty}+\right. \\
& \left.+\gamma^{-2} B_{2} O_{1} B_{1}^{\prime} P_{\infty}\right) \sigma+\left(B_{1}+B_{2} O_{1}\right) r  \tag{5.86}\\
q= & N \vartheta+\left(O_{2}+B_{2}^{\prime} P_{\infty}+\gamma^{-2} O_{1} B_{1}^{\prime} P_{\infty}\right) \sigma+O_{1} r \tag{5.87}
\end{align*}
$$

The controller $K_{\text {FIcr }}(s)$ described by these equations and by

$$
\begin{align*}
u_{c} & :=-B_{2}^{\prime} P_{\infty} x+q  \tag{5.88}\\
r & :=w-\gamma^{-2} B_{1}^{\prime} P_{\infty} x \tag{5.89}
\end{align*}
$$

is now shown to be an element of the set $\mathcal{F}_{\infty \gamma c r}$. The structure of such a controller apparently coincides with the one in fig. 5.18. Moreover, from eqs. (5.51)-(5.55), (5.85)-(5.89), it follows that, letting $\varepsilon:=\sigma-x$, the system resulting from the feedback connection of $P_{\text {FIc }}(s)$ with $K_{\text {FIcr }}(s)$ is described by the equations

$$
\begin{aligned}
\dot{x} & =A_{c} x+B_{1} w+B_{2} u_{c} \\
\dot{\vartheta} & =L \vartheta+\left(M_{2}+\gamma^{-2} M_{1} B_{1}^{\prime} P_{\infty}\right) \varepsilon+M_{2} x+M_{1} w \\
\dot{\varepsilon} & =\left(A_{c}-B_{2} B_{2}^{\prime} P_{\infty}+\gamma^{-2} B_{1} B_{1}^{\prime} P_{\infty}\right) \varepsilon \\
u_{c} & =N \vartheta+O_{2} x+O_{1} w+\left(O_{2}+B_{2}^{\prime} P_{\infty}+\gamma^{-2} O_{1} B_{1}^{\prime} P_{\infty}\right) \varepsilon \\
z & =C_{1 c} x+D_{12} u_{c}
\end{aligned}
$$

so that a comparison of these last equations with those relevant to the feedback connection of $P_{\text {FIc }}(s)$ with $K_{F I c}(s)$ (eqs. (5.51)-(5.55) and (5.83),(5.84)) leads to the following conclusions: First, the transfer functions from $w$ to $z$ are the same in both cases; second, the stability of the system which adopts $K_{\text {FIcr }}(s)$ is entailed by the stability of the system which utilizes $K_{F I c}(s)$, since matrix $A_{c}-B_{2} B_{2}^{\prime} P_{\infty}+$ $\gamma^{-2} B_{1} B_{1}^{\prime} P_{\infty}$ is stable.


Figure 5.19: The $R H_{\infty}$ performances of the controllers $L Q$ and $H$

It is only left to be proved that any controller possessing the structure shown in fig. 5.18, admissible in $R H_{\infty}$ for $P_{F I c}(s)$ and such that $\|T(z, w ; s)\|_{\infty}<\gamma$ belongs to the set $\mathcal{F}_{\infty \gamma c r}$ or, in other words, that $Q(s):=\left(A_{q}, B_{q}, C_{q}, D_{q}\right)$ is stable ( $A_{q}$ is stable) and $\|Q(s)\|_{\infty}<\gamma$. These conclusions are immediately drawn by exploiting Lemma E.3.

Example 5.1 Consider system (5.41)-(5.45) with

$$
\begin{gathered}
A=\left[\begin{array}{rr}
0 & 1 \\
-1 & -1
\end{array}\right], \quad B_{1}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad B_{2}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \\
C_{1}=\left[\begin{array}{l}
C_{11} \\
C_{12}
\end{array}\right]=\left[\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right], \quad D_{12}=\left[\begin{array}{l}
0 \\
1
\end{array}\right]
\end{gathered}
$$

A stabilizing controller is sought which makes small the effects of the disturbance $w$ over the first component $z_{1}$ of the output $z$. This goal can be attained by looking for the controller which is admissible in $R H_{\infty}$ for the given system and minimizes $\|T(z, w ; s)\|_{\infty}$. Notice that, in so doing, some kind of uncertainty in the knowledge of the system dynamics can be taken into account as well. It is in fact apparent (recall fig. 5.8), that perturbations of the matrix $A$ amounting to $\Delta_{A}=B_{1} \Omega C_{11}$ are effectively counteracted (under the stability point of view) by such a controller.

Eq. (5.46) admits a symmetric, positive semidefinite and stabilizing solution for $\gamma \geq$ $0.71:=\gamma_{M}$, while this does not happen when $\gamma \leq 0.705:=\gamma_{m}$, so that the minimal value attainable by $\|T(z, w ; s)\|_{\infty}$ belongs to the interval $\left(\gamma_{m}, \gamma_{M}\right]$.

For $\gamma=\gamma_{M}$ it results $F_{\infty}=-\left[\begin{array}{ll}0 & 0.89\end{array}\right]$. The performances of the resulting controller, labeled with $H$, are compared with those, labeled with $L Q$, relative to a controller designed within the $R H_{2}$ framework with $\Omega=0\left(F_{2}=-[00.41]\right)$. The quantities $\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty}$ and $\left\|T\left(z_{1}, w ; s\right)\right\|_{2}$ are plotted against the damping factor $\xi:=(1-\Omega) / 2$ in fig. 5.19 and 5.20 respectively. It is fairly apparent the more satisfactory behavior of the $H$ controller.

Remark 5.7 The proof of Theorem 5.2 presents a result which deserves some interest per se. More precisely, the existence of a symmetric, positive semidefinite and stabilizing solution of eq. (5.46), that is such that the matrix

$$
A_{c c}:=A_{c}-B_{2} B_{2}^{\prime} P_{\infty}+\gamma^{-2} B_{1} B_{1}^{\prime} P_{\infty}
$$



Figure 5.20: The $R H_{2}$ performances of the controllers $L Q$ and $H$
is stable, also entails the stability of the matrix

$$
A_{c \infty}:=A_{c}-B_{2} B_{2}^{\prime} P_{\infty}=A+B_{2} F_{\infty}
$$

which is the dynamic matrix of the control system resulting from the choice $Q(s)=0$ in the controller scheme of fig. 5.15 . The correctness of this result is implicitly guaranteed by the proof of Theorem 5.2 but can be verified with a different argument too. The matrix $P_{\infty}=P_{\infty}^{\prime} \geq 0$, being the stabilizing solution of eq. (5.46), is also the solution of the Lyapunov equation (in the unknown $P$ )

$$
0=P A_{c \infty}+A_{c \infty}^{\prime} P+W
$$

where

$$
W:=P_{\infty}\left(B_{2} B_{2}^{\prime}+\gamma^{-2} B_{1} B_{1}^{\prime}\right) P_{\infty}+C_{1 c}^{\prime} C_{1 c}
$$

If the pair $\left(A_{c \infty}, W\right)$ is detectable, then stability of matrix $A_{c \infty}$ follows from point (ii) of Lemma C.1. Detectability is readily proved by contradiction. Indeed, if $A_{c \infty} x=\lambda x, \operatorname{Re}(\lambda) \geq$ $0, W x=0, x \neq 0$, then $B_{1}^{\prime} P_{\infty} x=0$ and, consequently, $A_{c c} x=A_{c \infty} x=\lambda x$, so that $\lambda$ would be an eigenvalue, with nonnegative real part, of $A_{c c}$, which is stable by assumption.

The assumption that the pair $\left[\left(A-B_{2} D_{12}^{\prime} C_{1}\right),\left(I-D_{12} D_{12}^{\prime}\right) C_{1}\right]$ is detectable is not exploited in the proof of part (b) and of sufficiency of part $(a)$ of Theorem 5.2. Moreover, the existence of the stabilizing solution of eq. (5.46) implies, thanks to Remark 5.7 , the stabilizability of the pair $\left(A, B_{2}\right)$. It is therefore possible to state the following corollary which can be viewed as a useful side-product of the result presented in Theorem 5.2.

Corollary 5.1 Suppose that Assumption 5.2 holds. If, for a given positive scalar $\gamma$, there exists the symmetric, positive semidefinite and stabilizing solution of the Riccati equation (5.46), then there exists a controller $K_{F I}(s)$ which is admissible in $R H_{\infty}$ for $P_{F I}(s)$ and such that $\|T(z, w ; s)\|_{\infty}<\gamma$. Moreover, the set of controllers defined by point (b) in the statement of Theorem 5.2 constitutes the set $\mathcal{F}_{\infty \gamma r}$.

Remark 5.8 It is worth noticing that when $\gamma \rightarrow \infty$ the $R H_{\infty}$ controller of fig. 5.15 (with $Q(s)=0$ ) tends to the $R H_{2}$ controller of fig. 4.3 (with $Q(s)=0$ ).

Remark 5.9 (Parametrization of the set $\mathcal{F}_{\infty \gamma}$ ) A parametrization of the set $\mathcal{F}_{\infty \gamma}$ will be presented now. To this aim, consider the system $P_{F}(s)$ which results from system (5.41)-(5.45) after the control law

$$
\begin{equation*}
u=F_{\infty} x+v \tag{5.90}
\end{equation*}
$$

has been implemented, where $F_{\infty}$ is defined by eq. (5.49), namely the system

$$
\begin{aligned}
P_{F}(s) & =\left[\begin{array}{cc}
P_{F 11}(s) & P_{F 12}(s) \\
P_{F 21}(s) & P_{F 22}(s)
\end{array}\right] \\
& =\left[\begin{array}{c|cc}
A_{\mathrm{c} \infty} & B_{1} & B_{2} \\
\hline C_{1 \infty} & 0 & D_{12} \\
{\left[\begin{array}{l}
I \\
0
\end{array}\right]} & {\left[\begin{array}{l}
0 \\
I
\end{array}\right]} & {\left[\begin{array}{c}
0 \\
0
\end{array}\right]}
\end{array}\right]
\end{aligned}
$$

where

$$
A_{c \infty}:=A+B_{2} F_{\infty}, \quad C_{1 \infty}:=C_{1 c}-D_{12} B_{2}^{\prime} P_{\infty}
$$

having taken into account eq. (5.48). The set of controllers $K_{F}(s)$ which stabilize $P_{F}(s)$ obviously coincides with the set of controllers which stabilize $P_{F 22}(s)$. By mimicking the discussion in Remark 4.3 (recall that, thanks to Remark 5.7, matrix $A_{c \infty}$ is stable) and adopting the same kind of notation, it follows

$$
\left[I+\gamma^{-2} Q(s) B_{1}^{\prime} P_{\infty} \Phi(s) B_{2}\right]^{-1} Q(s)\left[I-\gamma^{-2} B_{1}^{\prime} P_{\infty} \Phi(s) B_{1}\right]=-\Theta(s) P_{F 21}(s)
$$

where $\Phi(s):=\left(s I-A_{c \infty}\right)^{-1}, \Theta(s) \in R H_{\infty}, Q(s) \in R H_{\infty}$. Letting

$$
\Delta(s):=I+\gamma^{-2} Q(s) B_{1}^{\prime} P_{\infty} \Phi(s) B_{2}
$$

the last equation can be rewritten as

$$
\Delta^{-1}(s) Q(s)\left[-\gamma^{-2} B_{1}^{\prime} P_{\infty} \quad I\right] P_{F 21}(s)=-\Theta(s) P_{F 21}(s)
$$

since

$$
P_{F 21}(s)=\left[\begin{array}{c}
\Phi(s) B_{1} \\
I
\end{array}\right]
$$

Assuming, for the moment being, that $\Delta^{-1}(s) \in R H_{\infty}$, it is easy to verify that a particular solution in $R H_{\infty}$ of such an equation is

$$
\bar{\Theta}(s)=\Delta^{-1}(s) Q(s)\left[\gamma^{-2} B_{1}^{\prime} P_{\infty} \quad-I\right]
$$

The general solution in $R H_{\infty}$ is therefore

$$
\Theta_{Q}(s)=\bar{\Theta}(s)+\hat{\Theta}(s)
$$

where $\hat{\Theta}(s)$ is any solution in $R H_{\infty}$ of the homogeneous equation

$$
\Theta(s) P_{F 21}(s)=0
$$

Letting $\Theta(s):=\left[\Delta^{-1}(s) \Theta_{1}(s) \quad \Theta_{2}(s)\right]$, this last equation implies that $\Theta_{2}(s)$ is given by $\Theta_{2}(s)=-\Delta^{-1}(s) \Theta_{1}(s) \Phi(s) B_{1}$ and thus

$$
\hat{\Theta}(s)=\Delta^{-1}(s) \Theta_{1}(s)\left[I \quad-\Phi(s) B_{1}\right]
$$

Notice that $\hat{\Theta}(s)$ is an element of $R H_{\infty}$ if and only if $\Theta_{1}(s)$ is such, since $\Delta^{-1}(s)$ has been assumed to belong to $R H_{\infty}$ and $\Phi(s) \in R H_{\infty}$ being $A_{c \infty}$ stable. Therefore, the set of functions $\Theta_{Q}(s)$ giving rise to controllers $K_{F}(s)$ which are admissible in $R H_{\infty}$ for $P_{F}(s)$ and such that $\|T(z, w ; s)\|_{\infty}<\gamma$ is defined by

$$
\Theta_{Q}(s)=\Delta^{-1}(s) \Lambda(s)
$$



Figure 5.21: The generic admissible controller for $P_{F I}(s)$
where

$$
\begin{aligned}
\Lambda(s):= & {\left[\Theta_{1}(s)+\gamma^{-2} Q(s) B_{1}^{\prime} P_{\infty} \quad-\left[Q(s)+\Theta_{1}(s) \Phi(s) B_{1}\right]\right] } \\
& \Theta_{1}(s) \in R H_{\infty}, Q(s) \in R H_{\infty},\|Q(s)\|_{\infty}<\gamma
\end{aligned}
$$

Thus, the transfer function (from $\left[x^{\prime} w^{\prime}\right]^{\prime}$ to $v$ ) of a generic controller which is admissible in $R H_{\infty}$ for $P_{F}(s)$ is (again, recall what has been presented in Remark 4.3)

$$
\begin{aligned}
K_{F}(s)= & -\left[I-\Theta_{Q}(s) P_{F 22}(s)\right]^{-1} \Theta_{Q}(s) \\
= & -\left[I-\Delta^{-1}(s) \Lambda(s) P_{F 22}(s)\right]^{-1} \Delta^{-1}(s) \Lambda(s) \\
= & -\left[\Delta(s)-\Lambda(s) P_{F 22}(s)\right]^{-1} \Lambda(s) \\
= & {\left[I-\Theta_{1}(s) \Phi(s) B_{2}\right]^{-1} \cdot } \\
& \cdot\left[-\gamma^{-2} Q(s) B_{1}^{\prime} P_{\infty}-\Theta_{1}(s) \quad Q(s)+\Theta_{1}(s) \Phi(s) B_{1}\right]
\end{aligned}
$$

By exploiting this equation, the generic controller which is admissible in $R H_{\infty}$ for $P_{F I}(s)$ can be represented by the block-scheme depicted in fig. 5.21. Such a scheme shows that

$$
\begin{aligned}
v_{L 0}= & \Theta_{1}(s) \Phi(s)\left[-\Phi^{-1}(s) x_{L}+B_{2} v_{L 0}+B_{1} w_{L}\right]+ \\
& +Q(s)\left(w_{L}-\gamma^{-2} B_{1}^{\prime} P_{\infty} x_{L}\right) \\
= & -\Theta_{1}(s) \Phi(s) x(0)+Q(s)\left(w_{L}-\gamma^{-2} B_{1}^{\prime} P_{\infty} x_{L}\right)
\end{aligned}
$$

having taken into consideration eqs. (5.41) and (5.90). Therefore, the effect of the parameter $\Theta_{1}(s)$ (which is responsible for the difference between the elements of the set $\mathcal{F}_{\infty \gamma}$ and those of the set $\mathcal{F}_{\infty \gamma r}$ ) on the control variable $u_{L}$ amounts to a term which depends on the initial conditions of system (5.41)-(5.45), only.

It remains to verify that $\Delta^{-1}(s) \in R H_{\infty}$. First, recall that $P_{\infty}$, being the stabilizing solution of eq. (5.46), is the stabilizing solution of the Riccati equation (in the unknown $P$ )

$$
0=P A_{c \infty}+A_{c \infty}^{\prime} P+P B_{2} B_{2}^{\prime} P+C_{1 c}^{\prime} C_{1 c}+\gamma^{-2} P B_{1} B_{1}^{\prime} P
$$

as well. In writing down this equation reference has been made to eqs. (5.48) and (5.49). This implies, thanks to Theorem 2.13, that system

$$
\Upsilon(s):=\left[\begin{array}{c|c}
A_{c \infty} & B_{2} \\
\hline C_{1 c} & 0 \\
\gamma^{-1} B^{\prime} P_{\infty} & 0
\end{array}\right]
$$

is such that $\|\Upsilon(s)\|_{\infty}<1$, which, in turn, entails (recall Lemma 2.23, Definition 2.23 and Remark 2.16) that

$$
\left\|\gamma^{-1} B_{1}^{\prime} P_{\infty} \Phi(s) B_{2}\right\|_{\infty}<1
$$



Figure 5.22: Proving stability of $\Delta^{-1}(s)$

Being $\|Q(s)\|_{\infty}<\gamma$, the system shown in fig. 5.22 is internally stable thanks to Theorem 5.1. The transfer function from $u$ to $e$ of such a system is precisely $\Delta^{-1}(s)$.

Remark 5.10 (Linear quadratic!differential game) The control problem whose solution has been presented in the preceding theorem may be viewed as a linear quadratic deterministic differential game ( $L Q D G$ ) with measurable state.

LQDG Problem Consider the $n$-th order system with initial state $x(0)=x_{0}$

$$
\begin{equation*}
\dot{x}=A x+\bar{B}_{1} \bar{w}+\bar{B}_{2} \bar{u} \tag{5.91}
\end{equation*}
$$

together with the cost functional

$$
J_{1}:=\int_{0}^{\infty}\left\{\left[\begin{array}{cc}
x^{\prime}(t) & \bar{u}^{\prime}(t)
\end{array}\right] L\left[\begin{array}{l}
x(t)  \tag{5.92}\\
\bar{u}(t)
\end{array}\right]-\gamma^{2} \bar{w}^{\prime}(t) R_{1} \bar{w}(t)\right\} d t
$$

where

$$
L:=\left[\begin{array}{cc}
Q & S \\
S^{\prime} & R_{2}
\end{array}\right] \geq 0, \quad R_{2}>0, \quad R_{1}>0
$$

Notice that the sign assumptions on $L$ and $R_{2}$ imply that

$$
\begin{equation*}
\hat{Q}:=Q-S R_{2}^{-1} S^{\prime} \geq 0 \tag{5.93}
\end{equation*}
$$

since $\hat{Q}=Z^{\prime} L Z$ with $Z^{\prime}=\left[I-S R_{2}^{-1}\right]$. Let $C_{11} \in R^{n \times n}$ be a factorization of $\hat{Q}$, so that

$$
\begin{equation*}
C_{11}^{\prime} C_{11}=\hat{Q} \tag{5.94}
\end{equation*}
$$

and define

$$
\begin{gather*}
C_{1}:=\left[\begin{array}{c}
C_{11} \\
R_{2}^{-\frac{1}{2}} S^{\prime}
\end{array}\right], \quad D_{12}:=\left[\begin{array}{c}
0 \\
I
\end{array}\right]  \tag{5.95}\\
u:=R_{2}^{\frac{1}{2}} \bar{u}, \quad z:=C_{1} x+D_{12} u, w:=R_{1}^{\frac{1}{2}} \bar{w}  \tag{5.96}\\
B_{2}:=\bar{B}_{2} R_{2}^{-\frac{1}{2}}, \quad B_{1}:=\bar{B}_{1} R_{1}^{-\frac{1}{2}} \tag{5.97}
\end{gather*}
$$

In view of eqs. (5.93), (5.97), eqs. (5.91)-(5.92) become

$$
\begin{gather*}
\dot{x}=A x+B_{1} w+B_{2} u  \tag{5.98}\\
J_{1}=\int_{0}^{\infty}\left[z^{\prime}(t) z(t)-\gamma^{2} w^{\prime}(t) w(t)\right] d t \tag{5.99}
\end{gather*}
$$

The solution of the differential game consists in finding the controller which stabilizes the resulting system and generates the control $u \in R H_{2}$ so as to minimize the functional $J_{1}$ corresponding to the worst input (disturbance) $w \in \mathrm{RH}_{2}$, that is corresponding to the input
$w$ which maximizes the same functional. Thus, the differential game is defined by eq. (5.98) and

$$
\begin{equation*}
\inf _{u \in R H_{2}} \sup _{w \in R H_{2}} J_{1} \tag{5.100}
\end{equation*}
$$

Notice that $\lim _{t \rightarrow \infty} x(t)=0$ since the resulting system is stable and $w \in R H_{2}$.
Assume that there exists the stabilizing solution $P_{\infty}=P_{\infty}^{\prime} \geq 0$ of the Riccati equation (5.46). Then, in view of eqs. (5.46),(5.48) and (5.98), it follows

$$
\begin{align*}
\frac{d}{d t}\left[x^{\prime}(t) P_{\infty} x(t)\right]= & -\left\|C_{1} x(t)\right\|^{2}-\gamma^{-2}\left\|B_{1}^{\prime} P_{\infty} x(t)\right\|^{2}+ \\
& +\left\|\left(D_{12}^{\prime} C_{1}+B_{2}^{\prime} P_{\infty}\right) x(t)\right\|^{2}+ \\
& +2<w(t), B_{1}^{\prime} P_{\infty} x(t)>+ \\
& +2<u(t), B_{2}^{\prime} P_{\infty} x(t)> \tag{5.101}
\end{align*}
$$

for generic inputs $u$ and $w$. In view of the preceding definitions, eq. (5.101) may also be written as

$$
\begin{aligned}
\frac{d}{d t}\left[x^{\prime}(t) P_{\infty} x(t)\right]= & -\|z(t)\|^{2}+\gamma^{2}\|w(t)\|^{2}+ \\
& -\gamma^{2}\left\|w(t)-\gamma^{-2} B_{1}^{\prime} P_{\infty} x(t)\right\|^{2}+ \\
& +\left\|u(t)-F_{\infty} x(t)\right\|^{2}
\end{aligned}
$$

By integrating both sides of this equation from 0 to $\infty$ and recalling that $x(\infty)=0$, it follows

$$
J_{1}=-\gamma^{2}\left\|w-\gamma^{-2} B_{1}^{\prime} P_{\infty} x\right\|_{2}^{2}+\left\|u-F_{\infty} x\right\|_{2}^{2}+x_{0}^{\prime} P_{\infty} x_{0}
$$

so that, consistently with eq. (5.100), the optimal control is given by

$$
u=F_{\infty} x
$$

while the worst disturbance is

$$
w=\gamma^{-2} B_{1}^{\prime} P_{\infty} x
$$

Remark 5.11 (Robust stabilizability) Consider a system characterized by a dynamic matrix which is not exactly known. More precisely, let

$$
\dot{x}=\left(A+\Delta_{A}\right) x+B_{2} u
$$

where the uncertainty introduced with the matrix $\Delta_{A}$ is defined by the equation

$$
\Delta_{A}:=B_{1} \Omega C_{1}
$$

Matrices $B_{1}$ and $C_{1}$ are thought as known, while an upper bound of the value of the norm is the only information available for matrix $\Omega$, that is

$$
\|\Omega\| \leq \gamma^{-1}
$$

with $\gamma$ a given positive scalar. Matrices $B_{1}, C_{1}, \Omega$ are constant. An algebraic control law $u=K x$ is sought such that the resulting control system is stable for each $\Omega,\|\Omega\| \leq \gamma^{-1}$. Apparently, the conditions under which this problem admits a solution constitute sufficient conditions for the robust stabilizability of the given system as well.

Suppose that there exists the symmetric, positive semidefinite and stabilizing solution $P_{\infty}$ of the Riccati equation (in the unknown $P$ )

$$
0=P A+A^{\prime} P+P\left(\gamma^{-2} B_{1} B_{1}^{\prime}-B_{2} B_{2}^{\prime}\right) P+C_{1}^{\prime} C_{1}
$$



Figure 5.23: The block-scheme for robust stabilizability

Then, $K=-B_{2}^{\prime} P_{\infty}$ solves the above stabilization problem. This claim is easily proved. First, consider the system $P_{x}(s)$ defined by

$$
\begin{aligned}
& \dot{x}=A x+B_{1} w+B_{2} u \\
& z=\left[\begin{array}{c}
z_{1} \\
z_{2}
\end{array}\right]=\left[\begin{array}{c}
C_{1} \\
0
\end{array}\right] x+\left[\begin{array}{l}
0 \\
I
\end{array}\right] u
\end{aligned}
$$

and depicted in fig. 5.23. Apparently, system $P_{x}(s)$ is obtained from the original one, with $\Omega=0$, by adding the performance output $z$ and the disturbance input $w$. Therefore, Corollary 5.1 can be applied to this system when $\Omega=0$, leading to the conclusion that the control law $u=-B_{2}^{\prime} P_{\infty} x$, corresponding to the choice $Q(s)=0$, is stabilizing and such that $\|T(z, w ; s)\|_{\infty}<\gamma$. From this last inequality it follows

$$
\begin{aligned}
\gamma^{2}>\|T(z, w ; s)\|_{\infty}^{2} & =\sup _{\substack{w \in R H_{2} \\
\|w\|_{2}=1}}\|z\|_{2}^{2} \\
& =\sup _{\substack{w \in R H_{2} \\
\|w\|_{2}=1}}\left\{\left\|z_{1}\right\|_{2}^{2}+\left\|z_{2}\right\|_{2}^{2}\right\} \\
& \geq \sup _{\substack{w \in R H_{2} \\
\|w\|_{2}=1}}\left\|z_{1}\right\|_{2}^{2} \\
& \geq\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty}^{2}
\end{aligned}
$$

In view of Theorem 5.1, the system in fig. 5.23 is stable, $\forall \Omega,\|\Omega\| \leq \gamma^{-1}$.
Example 5.2 Consider the system

$$
\begin{aligned}
\dot{x} & =\left(A+\Delta_{A}\right) x+B_{1} w_{1}+B_{2} u \\
z_{1} & =C_{1} x+u
\end{aligned}
$$

where

$$
\begin{gathered}
\Delta_{A}=\left[\begin{array}{l}
1 \\
1
\end{array}\right] \Omega\left[\begin{array}{ll}
0 & 1
\end{array}\right]=\left[\begin{array}{ll}
0 & \Omega \\
0 & \Omega
\end{array}\right], A=\left[\begin{array}{rr}
0 & 1 \\
-1 & -1
\end{array}\right] \\
B_{1}=\left[\begin{array}{l}
1 \\
0
\end{array}\right], \quad B_{2}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad C_{1}=\left[\begin{array}{ll}
0 & -2
\end{array}\right]
\end{gathered}
$$



Figure 5.24: The performances of the controllers $R$ and $H$


Figure 5.25: The disturbance feedforward problem

For $\Omega=0$ the Riccati equation (5.46) relevant to this system admits the positive semidefinite and stabilizing solution whenever $\gamma \geq 1.62$, while such a solution does not exist if $\gamma \leq 1.61$.

Taken $\gamma=1.62$, it results $F_{\infty}=-\left[\begin{array}{ll}367.68 & 227.58\end{array}\right]$. The corresponding controller will be denoted with the label $H$.

By proceeding along the same lines as in Remark 5.11, an output $z_{2}$ is defined as

$$
z_{2}:=\left[\begin{array}{ll}
0 & 1
\end{array}\right] x
$$

and an input $w_{2}$ is introduced, acting on the system through the matrix $\left[\begin{array}{ll}1 & 1\end{array}\right]^{\prime}$.
Corresponding to $\Omega=0$, the eq. (5.46) relevant to the new system characterized by the disturbance $w:=\left[\begin{array}{ll}w_{1} & w_{2}\end{array}\right]^{\prime}$ and the performance output $z:=\left[\begin{array}{ll}z_{1} & z_{2}\end{array}\right]^{\prime}$, admits a positive semidefinite and stabilizing solution for $\gamma \geq 2.46$, while such a solution does not exist if $\gamma \leq 2.45$. Taken $\gamma=2.46$ it results $F_{\infty r}=-[333.04$ 286.98]. The controller adopting such a $F_{\infty r}$ will be denoted with the label $R$. The graphs of $\left\|T\left(z_{1}, w_{1} ; s\right)\right\|_{\infty}$ corresponding to the adoption of the two controllers $R$ and $H$ are shown, as functions of $\Omega$, in fig. 5.24. The controller $R$ apparently behaves better for high values of $\Omega$.

Remark 5.12 (Disturbance feedforward) A problem similar to the one discussed in Remark 4.4 can be dealt with in the $R H_{\infty}$ context as well. Consider the system depicted in fig. 5.25 where $P_{D F}(s)$ is described by

$$
\begin{aligned}
& \dot{x}=A x+B_{1} w+B_{2} u \\
& z=C_{1} x+D_{12} u \\
& y=C_{2} x+w
\end{aligned}
$$



Figure 5.26: The set $\mathcal{F}_{\infty \gamma r}$ of the controllers $K_{D F r}(s)$
while $K_{D F}(s)$ is any controller which is admissible in $R H_{\infty}$ for $P_{D F}(s)$.
Assume that the pair $\left(A, B_{2}\right)$ is stabilizable, the pair $\left[\left(A-B_{2} D_{12}^{\prime} C_{1}\right),\left(I-D_{12} D_{12}^{\prime}\right) C_{1}\right]$ is detectable, the matrix $A-B_{1} C_{2}$ is stable and $D_{12}^{\prime} D_{12}=I$.

By proceeding as done in the $\mathrm{RH}_{2}$ setting (in particular, observe that Lemma E. 2 can be applied in the present $R H_{\infty}$ context as well) it is possible to directly give the solution of Problem 5.1 relative to system $P_{D F}(s)$
a) Existence of the symmetric, positive semidefinite and stabilizing solution $P_{\infty}$ of the Riccati equation (5.46);
b) The set $\mathcal{F}_{\infty \gamma r}$ of the controllers $K_{D F r}(s)$ is defined by the block-scheme of fig. 5.26 where

$$
N_{\infty}(s):=\left[\begin{array}{c|cc}
A-B_{1} C_{2}+B_{2} F_{\infty} & B_{1} & B_{2} \\
\hline F_{\infty} & 0 & I \\
-C_{2}-\gamma^{-2} B_{1}^{\prime} P_{\infty} & I & 0
\end{array}\right]
$$

$Q(s):=\Sigma\left(A_{q}, B_{q}, C_{q}, D_{q}\right)$, with $A_{q}$ stable, $\|Q(s)\|_{\infty}<\gamma$ and $F_{\infty}$ given by eq. (5.49).
Finally, notice that $\mathcal{F}_{\infty \gamma r}=\mathcal{F}_{\infty \gamma}$. In fact, what has been said concerning the problem in the $\mathrm{RH}_{2}$ setting can be applied with no modifications in the present framework, coming up to the same conclusions.

Example 5.3 Consider the system

$$
\begin{aligned}
\dot{x} & =\left(A+\Delta_{A}\right) x+B_{1} w+B_{2} u \\
z & =C_{1} x+D_{12} u \\
y & =C_{2} x+w
\end{aligned}
$$

where

$$
\begin{gathered}
A=\left[\begin{array}{rr}
0 & 1 \\
-1 & -1
\end{array}\right], \quad \Delta_{A}=L \Omega M, B_{1}=B_{2}=L=\left[\begin{array}{l}
0 \\
1
\end{array}\right] \\
C_{1}=C_{2}=\left[\begin{array}{ll}
1 & 0
\end{array}\right], M=\left[\begin{array}{ll}
0 & 1
\end{array}\right], \quad D_{12}=1
\end{gathered}
$$

The only uncertain system parameter is $\Omega$ : its nominal value is 0 and corresponds to a damping factor $\xi:=(1-\Omega) / 2=0.5$. A stabilizing controller is sought which makes the $R H_{\infty}$ norm of the transfer function from $w$ to $z$ small. If the perturbation $\Omega$ is ignored ( $\Omega=0$ ), it is possible to design a controller, denoted by the label $I C$ and with input $y$ and output $u$, which achieves a perfect indirect compensation of the disturbance, that is a controller such that $T(z, w ; s)=0$. More in detail, the transfer function of this controller is

$$
K_{I C}(s)=-\frac{1}{s^{2}+s+3}
$$



Figure 5.27: The performances of the controllers $I C$ and $R$

On the contrary, by exploiting what has been presented in Remark 5.11, a controller (denoted with the label $R$ ) can be designed which can account for $\Omega$ not being 0 in a more effective way. Notice that being $L=B_{1}$ it is useless introducing a further disturbance input, while being $M \neq C_{1}$ it is mandatory adding a new performance variable $z_{a}=M x$ to the formerly existing one $z$. Thus, the resulting problem has the same structure as those considered in Remark 5.12.

The Riccati equation relevant to the so restated problem does not admit the positive semidefinite and stabilizing solution when $\gamma \leq 0.7$, while such a solution exists if $\gamma \geq 0.71$. Associated with this last value of $\gamma$ it results $F_{\infty}=-\left[\begin{array}{ll}1 & 0.89\end{array}\right]$.

In fig. 5.27 the plots of $\|T(z, w ; s)\|_{\infty}$ corresponding to the above controllers are shown. The controller $R$ apparently behaves better for small values of the actual damping factor; moreover, its performance is somehow less sensitive (in terms of the $R H_{\infty}$ norm) to the variations of such a parameter. On the contrary, the controller $I C$ is to be preferred, as it should be expected, whenever the parameter $\Omega$ is rather precisely known.

Remark 5.13 The content of Remark 4.6 concerning the rank of $D_{12}$ and the condition $D_{12}^{\prime} D_{12}=I$ applies to the $R H_{\infty}$ setting with no changes. It is also apparent that the above given expressions get much simpler under the orthogonality assumption $D_{12}^{\prime} C_{1}=0$

Remark 5.14 Under Assumption 5.2, the detectability of the pair [ $\left(A-B_{2} D_{12}^{\prime} C_{1}\right),(I-$ $\left.\left.D_{12} D_{12}^{\prime}\right) C_{1}\right]$ (Assumption 5.1) is equivalent to asking the subsystem of $P_{F I}(s)$ corresponding to the transfer function $P_{F I 12}(s)$ between the input $u$ and the output $z$ (that is, system $\Sigma\left(A, B_{2}, C_{1}, D_{12}\right)$ not to have zeros in $R e(s) \geq 0$. This claim can easily be checked by means of the same arguments exploited in Remark 4.8.

### 5.4 The output estimation problem

In this section the problem of observing linear combinations of the state variables is dealt with in the $R H_{\infty}$ context. The system considered here is described by the state equations

$$
\begin{equation*}
\dot{x}=A x+B_{1} w+B_{2} u \tag{5.102}
\end{equation*}
$$



Figure 5.28: The output estimation problem

$$
\begin{align*}
& z=C_{1} x+u  \tag{5.103}\\
& y=C_{2} x+D_{21} w \tag{5.104}
\end{align*}
$$

As it has been done in Section 4.3, $P_{O E}(s)$ denotes the transfer function of this system, while $T(z, w ; s)$ is the transfer function from $w$ to $z$ in the block-scheme of fig. 5.28. The structure of system $P_{O E}(s)$ clearly indicates that solving Problem 5.1 for such a system can be really viewed as finding a "good" estimate (in the $R H_{\infty}$ sense) of the linear combination $-C_{1} x$ of the state variables. Indeed, should the controller $K_{O E}(s)$ be capable of zeroing $T(z, w ; s)$ so that $z_{L 0}=0$, then, apparently, $u_{L 0}=-C_{1} x_{L 0}$ would represent the best possible estimate of that linear combination.

The statement of the result below requires the following assumptions.
Assumption 5.3 The pair $\left[\left(A-B_{1} D_{21}^{\prime} C_{2}\right), B_{1}\left(I-D_{21}^{\prime} D_{21}\right)\right]$ is stabilizable and the pair $\left(A, C_{2}\right)$ is detectable.

Assumption 5.4 $D_{21} D_{21}^{\prime}=I$.
Assumption 5.5 $A-B_{2} C_{1}$ is stable.
It is now possible to prove the next theorem which makes reference to fig. 5.28.
Theorem 5.3 (Output estimation) Consider Problem 5.1 relative to system (5.102)(5.104). Then, under Assumptions 5.3-5.5, it has the solution
a) The existence of the symmetric, positive semidefinite and stabilizing solution $\Pi_{\infty}$ of the Riccati equation (in the unknown $\Pi$ )

$$
\begin{equation*}
0=\Pi A_{f}^{\prime}+A_{f} \Pi-\Pi\left(C_{2}^{\prime} C_{2}-\gamma^{-2} C_{1}^{\prime} C_{1}\right) \Pi+B_{1 f} B_{1 f}^{\prime} \tag{5.105}
\end{equation*}
$$

i.e., such that the matrix $A_{f c}$ given by

$$
\begin{equation*}
A_{f c}:=A_{f}-\Pi_{\infty}\left(C_{2}^{\prime} C_{2}-\gamma^{-2} C_{1}^{\prime} C_{1}\right) \tag{5.106}
\end{equation*}
$$

is stable. In eqs. (5.105), (5.106)

$$
\begin{equation*}
A_{f}:=A-B_{1} D_{21}^{\prime} C_{2}, B_{1 f}:=B_{1}\left(I-D_{21}^{\prime} D_{21}\right) \tag{5.107}
\end{equation*}
$$

b) The set $\mathcal{F}_{\infty \gamma r}$ of the controllers $K_{O E r}(s)$ is defined by the diagram of fig. 5.29, where

$$
M_{\infty}(s):=\left[\begin{array}{c|cc}
A_{f}-\Pi_{\infty} C_{2}^{\prime} C_{2}-B_{2} C_{1} & L_{\infty} & -B_{2}-\gamma^{-2} \Pi_{\infty} C_{1}^{\prime} \\
\hline C_{1} & 0 & I \\
C_{2} & I & 0
\end{array}\right]
$$



Figure 5.29: The set $\mathcal{F}_{\infty \gamma r}$ of the controllers $K_{O E r}(s)$

$$
\begin{align*}
& Q(s):=\Sigma\left(A_{q}, B_{q}, C_{q}, D_{q}\right) \text { with } A_{q} \text { stable and }\|Q(s)\|_{\infty}<\gamma, \text { having defined } \\
& \qquad L_{\infty}:=-\Pi_{\infty} C_{2}^{\prime}-B_{1} D_{21}^{\prime} \tag{5.108}
\end{align*}
$$

Proof First observe that Assumption 5.5 implies that the pair $\left(A, B_{2}\right)$ is stabilizable. This fact, together with the assumed detectability of the pair $\left(A, C_{2}\right)$, makes the necessary condition of Remark 5.6 satisfied.

Now consider the system $\hat{P}_{O E}(s):=P_{O E}^{\prime}(s)$. Then, from eqs. (5.102)-(5.104) it follows that such a system is given by

$$
\begin{aligned}
\dot{\xi} & =F \xi+G_{1} \zeta+G_{2} \eta \\
\omega & =H_{1} \xi+E \eta \\
\varphi & =H_{2} \xi+\zeta
\end{aligned}
$$

with

$$
\begin{equation*}
F:=A^{\prime}, \quad E:=D_{21}^{\prime} \tag{5.109}
\end{equation*}
$$

and, for $i=1,2$,

$$
\begin{equation*}
G_{i}:=C_{i}^{\prime}, \quad H_{i}:=B_{i}^{\prime} \tag{5.110}
\end{equation*}
$$

System $\hat{P}_{O E}(s)$ possesses the structure of system $P_{D F}(s)$ which has been introduced in Remark 5.12. Therefore, Problem 5.1 relative to system $\hat{P}_{O E}(s)$ is solved in the very same way, since the requirements there are satisfied by Assumptions 5.3-5.5. Thanks to Lemma E. 1 the results concerning system $P_{O E}$ follow by transposition of those concerning system $\hat{P}_{O E}(s)$, provided that eqs. (5.109) and (5.110) are taken into account.

Example 5.4 This example is aimed at pointing out an important difference existing between the solution of the estimation problem carried out in the $R H_{2}$ context and that obtained in the $R H_{\infty}$ one. Thus, consider system (5.102)-(5.104) with

$$
\begin{gathered}
A=\left[\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right], \quad B_{1}=I, \quad B_{2}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad C_{2}=\left[\begin{array}{ll}
1 & 0
\end{array}\right] \\
D_{21}=\left[\begin{array}{ll}
0 & 1
\end{array}\right], \quad C_{1}=\left[\begin{array}{ll}
1 & \alpha
\end{array}\right], \alpha>0
\end{gathered}
$$

Assumptions 5.3-5.5 (and therefore also Assumptions 4.3-4.5) are satisfied for each value of $\alpha$. The $\mathrm{RH}_{2}$ problem calls for the determination of the matrix $\Pi_{2}$ which is the stabilizing solution of the Riccati equation (4.64). Such a matrix is independent of $C_{1}$, so that also the matrix $L_{2}$ is such (recall Theorem 4.2). As a consequence of this fact and the structure of
system $M_{2}(s)$, the transfer function from the input $w$ to the state estimation error $\varepsilon:=x+x_{m}$ does not depend on $C_{1}$ (recall that $x_{m}$, the state of system $M_{2}(s)$, is the opposite of the estimate of the state of system $P_{O E}(s)$, since the relevant linear combination of the state variables is $\left.-C_{1} x\right)$. Indeed, for any $Q(s)$, it results

$$
T\left(\left(x+x_{m}\right), w ; s\right)=\left[s I-\left(A+L_{2} C_{2}\right)\right]^{-1}\left(L_{2} D_{21}+B_{1}\right)
$$

This property in no more verified in the $R H_{\infty}$ context. In fact, matrix $\Pi_{\infty}$, which is the stabilizing solution of the Riccati equation (5.105), and, consequently, matrix $L_{\infty}$ too, depends on $C_{1}$ (see Theorem 5.3). This fact, together with the structure of system $M_{\infty}(s)$, implies that the transfer function from the input $w$ to the state estimation error $\varepsilon:=x+x_{m}$, where $x_{m}$ is the state variable of system $M_{\infty}(s)$, is given by

$$
T\left(\left(x+x_{m}\right), w ; s\right)=\left[s I-\left(A+L_{\infty} C_{2}\right)\right]^{-1}\left(L_{\infty} D_{21}+B_{1}\right)
$$

when $Q(s)=0$. This function, though formally identical to that resulting in the $\mathrm{RH}_{2}$ context, depends on $C_{1}$, that is from the actual linear combination of the state of $P_{O E}(s)$ which one is willing to estimate.

With reference to the considered system, for instance, if $\alpha$ is set equal to 10 or equal to 100 , one obtains, corresponding to $\gamma=130, L_{\infty}=-\left[\begin{array}{ll}1.00 & 1.00\end{array}\right]^{\prime}$ and $L_{\infty}=-\left[\begin{array}{ll}6.81 & 8.17\end{array}\right]^{\prime}$, respectively.

Remark 5.15 (Parametrization of the set $\mathcal{F}_{\infty \gamma}$ ) Observe that $\mathcal{F}_{\infty \gamma}=\mathcal{F}_{\infty \gamma r}$. In fact, as in the proof of Theorem 5.3, the $R H_{\infty}$ admissible controllers for $P_{O E}(s)$ can be obtained by transposition of those which are $R H_{\infty}$ admissible for $\hat{P}_{O E}(s)$, this last system possessing the structure of system $P_{D F}(s)$ (see Remark 5.12). Therefore, being $\mathcal{F}_{\infty \gamma}=\mathcal{F}_{\infty \gamma r}$ for system $P_{D F}(s)$, the same conclusion must hold for system $P_{O E}(s)$.

Remark 5.16 (Optimal state filtering) The control problem associated with system $P_{O E}(s)$ can be viewed as an (optimal) state filtering problem in $R H_{\infty}$ for the stable system described by

$$
\begin{aligned}
\dot{x} & =A x+\bar{B}_{1} w_{1} \\
y & =C_{2} x+w_{2}
\end{aligned}
$$

The signals $w_{1}$ and $w_{2}$ belong to the set

$$
\mathcal{W}:=\left\{w_{\alpha} \mid w_{\alpha} \in R L_{2}[0, \infty),\left\|w_{\alpha}\right\|_{2} \leq \alpha\right\}
$$

A filter is sought for a given linear combination $S x$ of the state variables such that, denoting with $u$ the filter output, it results

$$
J=\min _{u \in R L_{2}[0, \infty)} \sup _{\substack{w_{1} \in \mathcal{W} \\ w_{2} \in \mathcal{W}}}\|u-S x\|_{2}
$$

Define $w:=\left[w_{1}^{\prime} w_{2}^{\prime}\right]^{\prime}, B_{1}:=\left[\bar{B}_{1} 0\right], D_{21}:=[0 I], C_{1}:=-S, z:=C_{1} x+u$. Then the problem can conveniently be stated by considering the system described by the equations

$$
\begin{aligned}
\dot{x} & =A x+B_{1} w \\
z & =C_{1} x+u \\
y & =C_{2} x+D_{21} w
\end{aligned}
$$

and seeking for the controller-filter such that the $R H_{\infty}$ norm of the transfer function from $w$ to $z$ is less than a given positive scalar $\gamma$. Indeed, observing that $\|w\|_{2} \leq \alpha \sqrt{2}:=\beta$, it results

$$
\sup _{\substack{w_{1} \in \mathcal{W} \\ w_{2} \in \mathcal{W}}}\|z\|_{2} \leq \sup _{\substack{w \in R H_{2} \\\|w\|_{2} \leq \beta}}\|z\|_{2}=\beta \sup _{\substack{w \in R H_{2} \\ w \neq 0}} \frac{\|z\|_{2}}{\|w\|_{2}}
$$



Figure 5.30: The $\mathrm{RH}_{2}$ norm of the four filters
so that $J \leq \beta\|T(z, w ; s)\|_{\infty}$. Thus, the smaller the value of $\gamma$, the tighter would be the upper bound of the maximum attainable value of $J$. What has previously been shown in Theorem 5.3 can now be applied to the so recast problem, since the underlying system possesses the very same structure of system $P_{O E}(s)$ and Assumptions 5.3-5.5 are verified (recall the definition of $D_{21}$, the stability of matrix $A$ and notice that $B_{1} D_{21}^{\prime}=0$ and $B_{2}=0$ ).

Example 5.5 Consider the system

$$
\begin{aligned}
\dot{x} & =A x+\bar{B}_{1} w_{1} \\
y & =C_{2} x+w_{2}
\end{aligned}
$$

where $w_{1}$ and $w_{2}$ are noises the features of which are to be specified later on, while

$$
A=\left[\begin{array}{cc}
0 & 1 \\
-1 & -1+\Omega
\end{array}\right], \quad \bar{B}_{1}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad C_{2}=\left[\begin{array}{ll}
1 & 0
\end{array}\right]
$$

The parameter $\Omega$ accounts for the uncertain knowledge of the system damping factor $\xi$. Indeed, the characteristic polynomial of $A$ is $\psi(s)=s^{2}+s(1-\Omega)+1=s^{2}+2 \xi \omega_{n} s+\omega_{n}^{2}$, with damping factor $\xi=(1-\Omega) / 2$ and natural frequency $\omega_{n}=1$.

The performances of four different filters for the whole state vector are now compared. Three of them, denoted with $K_{1}, K_{2}$ and $K_{3}$, respectively, are standard Kalman's filters which have been designed under the assumption that $w_{1}$ and $w_{2}$ are uncorrelated, zeromean, Gaussian white noises. In the considered cases, their intensity are $W_{1}=W_{2}=1$ for the $K_{1}$ filter, $W_{1}=0.5, W_{2}=1$ for the $K_{2}$ filter, $W_{1}=1, W_{2}=0.5$ for the $K_{3}$ filter.

On the contrary, the fourth filter, denoted with the label $H$, is designed accordingly to what has been presented in Remark 5.16. The adopted value for the scalar $\gamma$ is 1.001 (notice that for $\gamma \leq 1$ the positive semidefinite and stabilizing solution of the relevant Riccati equation does not exist).

The design of the four filters is carried out in nominal conditions ( $\Omega=0 \Leftrightarrow \xi=.5$ ), while their performances are evaluated by connecting them to the system perturbed in correspondence with various values of $\xi$.

Two performance criteria have been adopted. Both of them make reference to the transfer function $T(z, w ; s)$, where $w:=\left[\begin{array}{ll}w_{1} & w_{2}\end{array}\right]^{\prime}$ and $z$ is the state estimation error. The criteria are the $R H_{2}$ and $R H_{\infty}$ norms of $T(z, w ; s)$, respectively: their plots against $\xi$ are shown in fig. 5.30 and 5.31.


Figure 5.31: The $R H_{\infty}$ norm of the four filters

The two graphs show that the filter designed within the $R H_{\infty}$ framework is somehow less sensitive with respect to $\xi$ and entails lower values of the norm for small values of the damping factor, namely close to instability.

Remark 5.17 Analogously to what has been done in Remark 5.7 with reference to the Riccati equation (5.46), the existence of the symmetric, positive semidefinite and stabilizing solution of eq. (5.105), i. e. , such that

$$
A_{f c}:=A_{f}-\Pi_{\infty}\left(C_{2}^{\prime} C_{2}-\gamma^{-2} C_{1}^{\prime} C_{1}\right)
$$

is stable, also implies the stability of

$$
A_{f \infty}:=A_{f}-\Pi_{\infty} C_{2}^{\prime} C_{2}=A+L_{\infty} C_{2}
$$

which is the dynamic matrix of the $R H_{\infty}$ state observer discussed in Remark 5.16 when $Q(s)=0$. The proof of this claim is identical to the one presented in Remark 5.7.

It is easy to verify that the proof of point (b) and of the sufficiency part of point (a) of Theorem 5.3 does not exploit the assumption that the pair $\left[\left(A-B_{1} D_{21}^{\prime} C_{2}\right), B_{1}(I-\right.$ $\left.\left.D_{21}^{\prime} D_{21}\right)\right]$ is stabilizable. Moreover, the existence of the stabilizing solution of eq. (5.105) implies, thanks to Remark 5.17 , the detectability of the pair $\left(A, C_{2}\right)$. Therefore, it is possible to state the following corollary which has to be viewed as a useful side-product of Theorem 5.3.

Corollary 5.2 Suppose that Assumptions 5.4 and 5.5 hold. If, for a given positive scalar $\gamma$, there exists the symmetric, positive semidefinite and stabilizing solution of the Riccati equation (5.105), then there exists a controller $K_{O E}(s)$ which is admissible in $R H_{\infty}$ for $P_{O E}(s)$ and such that $\|T(z, w ; s)\|_{\infty}<\gamma$. Moreover, the set of controllers defined by point (b) in the statement of Theorem 5.3 constitutes the set $\mathcal{F}_{\infty \gamma r}$.

Remark 5.18 It is worth noticing that, also in the output estimation case, when $\gamma \rightarrow \infty$ the $R H_{\infty}$ controller of fig. 5.29 (with $Q(s)=0$ ) tends to the $R H_{2}$ controller of fig. 4.13 (with $Q(s)=0$ ).

Remark 5.19 (Robust parametric filtering) An interesting role is played by the Riccati equation (5.105) also when the problem at hand is filtering the state of a system affected by parametric uncertainty and contemporarily looking for a meaningful (namely, small) upper bound to the filtering error. More precisely, consider the $n$-th order stochastic system

$$
\begin{aligned}
\dot{x} & =\left(A+\Delta_{A}\right) x+B_{1} w \\
y & =C x+v
\end{aligned}
$$

with the matrix $A$ stable and the pair $\left(A, B_{1}\right)$ reachable. Moreover, the inputs $w$ and $v$ are zero mean, uncorrelated white noises with identity intensity. The perturbation $\Delta_{A}$ is defined by

$$
\Delta_{A}:=B_{2} \Omega C_{1}, \quad\|\Omega\| \leq \gamma^{-1}
$$

where the matrices $B_{2}, C_{1}$ and the positive scalar $\gamma$ are given. The problem of robust parametric filtering with cost $Q$ consists in finding a stable state filter of the form

$$
\dot{\xi}=A_{f} \xi+K_{f} y
$$

and a positive semidefinite matrix $Q$ which provides an upper bound for $X_{\varepsilon a}(\Omega)$, the asymptotic covariance matrix of the estimation error

$$
\varepsilon:=\xi-x
$$

correspondingly to any perturbation $\Delta_{A}$ of the above specified form. In other words, it is required that

$$
X_{\varepsilon a}(\Omega):=\lim _{t \rightarrow \infty} E\left[\varepsilon(t) \varepsilon^{\prime}(t)\right] \leq Q, \quad \forall\|\Omega\| \leq \gamma^{-1}
$$

Let $C_{2}(\beta):=\beta C, B(\beta):=\left[\begin{array}{ll}\beta^{-1} B_{1} & B_{2}\end{array}\right]$. Then it will be proved that, for each $\beta>0$ belonging to the set of the $\beta^{\prime} s$ for which the Riccati equation (in the unknown $\Pi$ )

$$
\begin{equation*}
0=A \Pi+\Pi A^{\prime}+\Pi\left(\gamma^{-2} C_{1}^{\prime} C_{1}-C_{2}^{\prime}(\beta) C_{2}(\beta)\right) \Pi+B(\beta) B^{\prime}(\beta) \tag{5.111}
\end{equation*}
$$

admits the stabilizing solution $\Pi_{\infty}(\beta)=\Pi_{\infty}^{\prime}(\beta) \geq 0$, the matrices

$$
\begin{align*}
& A_{f \beta}:=A+\Pi_{\infty}(\beta)\left(\gamma^{-2} C_{1}^{\prime} C_{1}-\beta^{2} C^{\prime} C\right)  \tag{5.112}\\
& K_{f \beta}:=\beta^{2} \Pi_{\infty}(\beta) C^{\prime}, \quad Q_{\beta}:=\beta^{2} \Pi_{\infty}(\beta) \tag{5.113}
\end{align*}
$$

define a solution of the above stated problem.
Remark 5.2 is exploited to prove this claim. In fact, let $\eta:=\left[\varepsilon^{\prime} \xi^{\prime}\right]^{\prime}$ and $\varphi:=\left[w^{\prime} v^{\prime}\right]^{\prime}$, then the connection of the system with the filter is described by the equation

$$
\dot{\eta}=\left(F+\Delta_{F}\right) \eta+G \varphi
$$

where

$$
\begin{gathered}
F:=\left[\begin{array}{cc}
A-K_{f} C & A_{f}+K_{f} C-A \\
-K_{f} C & A_{f}+K_{f} C
\end{array}\right], G:=\left[\begin{array}{cc}
-B_{1} & K_{f} \\
0 & K_{f}
\end{array}\right] \\
\Delta_{F}:=\bar{B} \Omega \bar{C}:=\left[\begin{array}{c}
B_{2} \\
0
\end{array}\right] \Omega\left[\begin{array}{ll}
C_{1} & -C_{1}
\end{array}\right]
\end{gathered}
$$

Of course, matrix $F$ is stable since matrices $A$ and $A_{f}$ are both stable. Corresponding to the given matrices $K_{f}$ and $A_{f}$, assume that there exists the stabilizing solution $\Theta_{\infty}(\beta)=$ $\Theta_{\infty}^{\prime}(\beta) \geq 0$ of the Riccati equation (in the unknown $\Theta$ )

$$
\begin{equation*}
0=\Theta F^{\prime}+F \Theta+\gamma^{-2} \Theta \bar{C}^{\prime} \bar{C} \Theta+\bar{B} \bar{B}^{\prime}+\beta^{-2} G G^{\prime} \tag{5.114}
\end{equation*}
$$

Then, in view of Remark 5.2 it follows

$$
X_{\eta a}(\Omega):=\lim _{t \rightarrow \infty} E\left[\eta(t) \eta^{\prime}(t)\right] \leq \beta^{2} \Theta_{\infty}(\beta), \quad \forall\|\Omega\| \leq \gamma^{-1}
$$

so that, by partitioning matrix $\Theta$ consistently with the structure of $F$, that is letting

$$
\Theta:=\left[\begin{array}{ll}
\Theta_{11} & \Theta_{12} \\
\Theta_{21} & \Theta_{22}
\end{array}\right]
$$

it results

$$
X_{\varepsilon a}(\Omega) \leq \beta^{2} \Theta_{\infty 11}(\beta), \quad \forall\|\Omega\| \leq \gamma^{-1}
$$

Further, if the matrices $A_{f}$ and $K_{f}$ are given by eqs. (5.112) and (5.113), then it is straightforward to verify that

$$
\bar{\Theta}(\beta):=\left[\begin{array}{cc}
\Pi_{\infty}(\beta) & 0 \\
0 & S_{\infty}(\beta)-\Pi_{\infty}(\beta)
\end{array}\right]
$$

satisfies eq. (5.114), provided that the stabilizing solution $S_{\infty}(\beta)=S_{\infty}^{\prime}(\beta) \geq 0$ exists of the Riccati equation (in the unknown $S$ )

$$
\begin{equation*}
0=A S+S A^{\prime}+\gamma^{-2} S C_{1}^{\prime} C_{1} S+B(\beta) B^{\prime}(\beta) \tag{5.115}
\end{equation*}
$$

Notice that $\bar{\Theta}(\beta)$ is a solution of the Lyapunov equation (in the unknown $\Theta$ )

$$
0=\Theta F^{\prime}+F \Theta+\gamma^{-2} \bar{\Theta}(\beta) \bar{C}^{\prime} \bar{C} \bar{\Theta}(\beta)+\bar{B} \bar{B}^{\prime}+\beta^{-2} G G^{\prime}
$$

Hence, from Lemma C.1, the stability of matrix $F$ entails that $\bar{\Theta}(\beta) \geq 0$. Moreover,

$$
F+\gamma^{-2} \bar{\Theta}(\beta) \bar{C}^{\prime} \bar{C}=\left[\begin{array}{cc}
A_{f \beta} & 0 \\
\star & A+\gamma^{-2} S_{\infty}(\beta) C_{1}^{\prime} C_{1}
\end{array}\right]
$$

where the symbol " $\star$ " denotes a matrix of no interest. Hence, $\bar{\Theta}(\beta)$ is also stabilizing, so that $\bar{\Theta}(\beta)=\Theta_{\infty}(\beta)$. Observe that, in view of Remark 5.2, it is

$$
X_{x a}(\Omega):=\lim _{t \rightarrow \infty} E\left[x(t) x^{\prime}(t)\right] \leq \beta^{2} S_{\infty}(\beta), \quad \forall\|\Omega\| \leq \gamma^{-1}
$$

As for the existence of the solutions $\Pi_{\infty}(\beta)$ and $S_{\infty}(\beta)$, first notice that the arguments exploited in Remark 5.2 can be applied here, thus allowing one to make reference to a nonempty set $\mathcal{B}$ of $\beta$ 's corresponding to which a solution $S_{\infty}(\beta)$ of eq. (5.115) exists.

Second, observe that for $\beta \in \mathcal{B}, S_{\infty}(\beta)$ is nonsingular because the pair $\left(A, B_{1}\right)$ is reachable. In fact, consider eq. (5.115) (which is solved by $S_{\infty}(\beta)$ ) and suppose that $S_{\infty}(\beta) x=0, x \neq 0$. If both sides of eq. (5.115) are premultiplied by $x^{\prime}$ and postmultiplied by $x$, then, in view of the definition of $B(\beta)$, one can conclude that $B_{1}^{\prime} x=0$ and $S_{\infty}(\beta) A^{\prime} x=0$. If both sides of eq. (5.115) are now premultiplied by $x^{\prime} A$ and postmultiplied by $A^{\prime} x$, one can conclude that $B_{1}^{\prime} A^{\prime} x=0$ and $S_{\infty}(\beta)\left(A^{\prime}\right)^{2} x=0$. By iterating these operations it follows

$$
x^{\prime}\left[\begin{array}{lllll}
B_{1} & A B_{1} & A^{2} B_{1} & \cdots & A^{n-1} B_{1}
\end{array}\right]=0
$$

which contradicts the reachability of the pair $\left(A, B_{1}\right)$.
Since $S_{\infty}(\beta)$ is nonsingular, its inverse satisfies the Riccati equation (in the unknown $S^{-1}$ ) which results from premultiplying and postmultiplying by $S^{-1}$ eq. (5.115), namely

$$
\begin{equation*}
0=S^{-1} A+A^{\prime} S^{-1}+\gamma^{-2} C_{1}^{\prime} C_{1}+S^{-1} B(\beta) B^{\prime}(\beta) S^{-1} \tag{5.116}
\end{equation*}
$$

Third, recall that for $\beta \in \mathcal{B}, A+\gamma^{-2} S_{\infty}(\beta) C_{1}^{\prime} C_{1}$ is stable, so that if both sides of eq. (5.115) with $S=S_{\infty}(\beta)$ are premultiplied by $S_{\infty}^{-1}(\beta)$ one gets

$$
S_{\infty}^{-1}(\beta)\left[A+\gamma^{-2} S_{\infty}(\beta) C_{1}^{\prime} C_{1}\right] S_{\infty}(\beta)=-\left[A+B(\beta) B^{\prime}(\beta) S_{\infty}^{-1}(\beta)\right]^{\prime}
$$

from which it follows that matrix

$$
A_{s}:=-\left[A+B(\beta) B^{\prime}(\beta) S_{\infty}^{-1}(\beta)\right]
$$

is stable. Then the Riccati equation (in the unknown $V$ )

$$
\begin{align*}
0= & V\left[A+B(\beta) B^{\prime}(\beta) S_{\infty}^{-1}(\beta)\right]+\left[A+B(\beta) B^{\prime}(\beta) S_{\infty}^{-1}(\beta)\right]^{\prime} V+ \\
& +V B(\beta) B^{\prime}(\beta) V-C_{2}^{\prime}(\beta) C_{2}(\beta) \tag{5.117}
\end{align*}
$$

which can equivalently be written as

$$
0=V A_{s}+A_{s}^{\prime} V-V B(\beta) B^{\prime}(\beta) V+C_{2}^{\prime}(\beta) C_{2}(\beta)
$$

admits the stabilizing solution $V_{\infty}(\beta)=V_{\infty}^{\prime}(\beta) \geq 0$ since matrix $A_{s}$ is stable (recall Lemma C.4). Therefore, $\Lambda_{\infty}(\beta):=S_{\infty}^{-1}(\beta)+V_{\infty}(\beta)$ is positive definite because $S_{\infty}^{-1}$ is positive definite. By summing up eq. (5.117) to eq. (5.116) one can conclude that $\Lambda_{\infty}(\beta)$ solves the Riccati equation (in the unknown $\Lambda$ )

$$
\begin{equation*}
0=\Lambda A+A^{\prime} \Lambda+\Lambda B(\beta) B^{\prime}(\beta) \Lambda+\gamma^{-2} C_{1}^{\prime} C_{1}-C_{2}^{\prime}(\beta) C_{2}(\beta) \tag{5.118}
\end{equation*}
$$

Thus, its inverse (which actually exists because $\Lambda_{\infty}(\beta)>0$ ) solves the equation which derives from eq. (5.118) after both sides have been premultiplied and postmultiplied by $\Lambda^{-1}$. Such an equation coincides with eq. (5.111) which therefore admits a symmetric, positive definite solution given by

$$
\Pi_{\infty}(\beta):=\Lambda_{\infty}^{-1}(\beta)=\left[S_{\infty}^{-1}(\beta)+V_{\infty}(\beta)\right]^{-1}
$$

It is only left to be proved that such a solution is the stabilizing one. By exploiting the fact that $\Pi_{\infty}(\beta)$ solves eq. (5.111), it results

$$
\begin{aligned}
A_{s}-B(\beta) B^{\prime}(\beta) V_{\infty}(\beta)= & -A-B(\beta) B^{\prime}(\beta) \Pi_{\infty}^{-1}(\beta) \\
= & -\left[A \Pi_{\infty}(\beta)+B(\beta) B^{\prime}(\beta)\right] \Pi_{\infty}^{-1}(\beta) \\
= & \Pi_{\infty}(\beta)\left[A^{\prime}+\left(\gamma^{-2} C_{1}^{\prime} C_{1}-C_{2}^{\prime}(\beta) C_{2}(\beta)\right) .\right. \\
& \left.\cdot \Pi_{\infty}(\beta)\right] \Pi_{\infty}^{-1}(\beta)
\end{aligned}
$$

so that $A^{\prime}+\left(\gamma^{-2} C_{1}^{\prime} C_{1}-C_{2}^{\prime}(\beta) C_{2}(\beta)\right) \Pi_{\infty}(\beta)$ is stable because $A_{s}-B(\beta) B^{\prime}(\beta) V_{\infty}(\beta)$ is stable (recall that $V_{\infty}(\beta)$ is the stabilizing solution).

In order to tighten as much as possible the so obtained bound for $X_{\varepsilon a}(\Omega)$, the best value for the parameter $\beta$ should be found: this task may be performed by means of a suitable (iterative) searching method.

Example 5.6 Consider the system described in Example 5.5. A filter, labeled with the symbol $R$, is to be designed according to what has been presented in Remark 5.19. More precisely, letting $w_{1}$ and $w_{2}$ be zero mean white noises with identity intensity, a filter is sought which supplies an as small as possible upper bound to the trace of the asymptotic estimation error covariance matrix. This goal corresponds to minimizing trace $\left[\beta^{2} \Pi_{\infty}(\beta)\right]$ with respect to the parameter $\beta$, where $\Pi_{\infty}(\beta)$ is the stabilizing solution of eq. (5.111). A suitable computer routine gives $\beta^{\circ}=1.13$. The performance of the corresponding filter are compared with those of the Kalman's filter $K_{1}$ and observer $H$ which have been designed in Example 5.5. The same performance criteria (namely, $R H_{2}$ and $R H_{\infty}$ norms of the transfer function from $w:=\left[w_{1} w_{2}\right]$ ) are adopted. The more satisfactory behavior of the robust (with respect to parameter variations) filter for small values of the system damping factor is put into evidence in fig. 5.32 and 5.33.

Remark 5.20 (Full control) The problem considered in Remark 4.12 can be dealt with also within the $R H_{\infty}$ context. The problem refers to the system $P_{F C}(s)$ depicted in fig. 5.34 and described by the equations

$$
\begin{aligned}
\dot{x} & =A x+B_{1} w+\left[\begin{array}{ll}
I & 0
\end{array}\right] u \\
z & =C_{1} x+\left[\begin{array}{ll}
0 & I
\end{array}\right] u \\
y & =C_{2} x+D_{21} w \\
u & =\left[\begin{array}{ll}
u_{1}^{\prime} & u_{2}^{\prime}
\end{array}\right]^{\prime}
\end{aligned}
$$



Figure 5.32: The $\mathrm{RH}_{2}$ norm of the three filters

In fig. 5.34 $K_{F C}(s)$ is a generic controller admissible in $R H_{\infty}$ for $P_{F C}(s)$. Assume that the pair $\left(A, C_{2}\right)$ is detectable, the pair $\left[\left(A-B_{1} D_{21}^{\prime} C_{2}\right), B_{1}\left(I-D_{21}^{\prime} D_{21}\right)\right]$ is stabilizable and $D_{21} D_{21}^{\prime}=I$. By exploiting also in the present framework what has been utilized in the $R H_{2}$ context (in particular, note that Lemma E. 1 can be applied both in the $R H_{2}$ and $R H_{\infty}$ settings), it is possible to directly claim that Problem 5.1 relative to system $P_{F C}(s)$ has the solution
a) Existence of the symmetric, positive semidefinite and stabilizing solution $\Pi_{\infty}$ of the Riccati equation (5.105);
b) The set $\mathcal{F}_{\infty \gamma r}$ of the controllers $K_{F C r}(s)$ is defined by the block-scheme of fig. 5.35 where $Q(s):=\Sigma\left(A_{q}, B_{q}, C_{q}, D_{q}\right)$, with $A$ stable and $\|Q(s)\|_{\infty}<\gamma$.
Finally, the set $\mathcal{F}_{\infty \gamma}$ of the controllers which are $R H_{\infty}$ admissible for $P_{F C}(s)$ can be easily obtained by exploiting (via transposition) what has been shown in Remark 5.9 with reference to system $P_{F I}(s)$. The block-scheme which defines such a set is shown in fig. 5.36 (note that it is the "transpose" of fig. 5.21) where $\Theta_{1}(s) \in R H_{\infty}, \Psi(s):=\left(s I-A-L_{\infty} C_{2}\right)^{-1}$.

Remark 5.21 The comments on the matrix $D_{21}$ presented in Remark 4.13 for the $\mathrm{RH}_{2}$ setting can be done with no changes in the present context too with reference to Assumption 5.4. Further, many of the above given formulas greatly simplify if the additional orthogonality assumption $D_{21} B_{1}^{\prime}=0$ is done.

Remark 5.22 Under Assumption 5.4, the stabilizability of the pair [ $\left(A-B_{1} D_{21}^{\prime} C_{2}\right), B_{1}(I-$ $\left.\left.D_{21}^{\prime} D_{21}\right)\right]$ (Assumption 5.3) is equivalent to requiring that the subsystem of $P_{O E}(s)$ corresponding to the transfer function $P_{O E 21}$ between the input $w$ and the output $y$, that is system $\Sigma\left(A, B_{1}, C_{2}, D_{21}\right)$, does not have zeros in the closed right half plane. The proof of this claim exploits the same arguments adopted in proving an analogous result in Remark 4.15 .

Remark 5.23 The discussion in Remark 4.16 concerning a possible direct dependence of the output variable $y$ on the control variable $u$ applies to the $R H_{\infty}$ setting as well. In fact, suppose that eq. (5.104) is substituted by

$$
\begin{equation*}
y=C_{2} x+D_{21} w+D_{22} u \tag{5.1.19}
\end{equation*}
$$



Figure 5.33: The $R H_{\infty}$ norm of the three filters


Figure 5.34: The full control problem

Letting

$$
\begin{equation*}
\bar{y}:=y-D_{22} u=C_{2} x+D_{21} w \tag{5.120}
\end{equation*}
$$

the solution to Problem 5.1 relative to the system $\bar{P}_{O E}(s)$, namely to the system described by eqs. (5.102), (5.103), (5.120), is supplied by Theorem 5.3. If $\bar{K}_{O E}(s)$ is a $R H_{\infty}$ admissible controller for $\bar{P}_{O E}(s)$, then the controller $K_{O E}(s)$ defined in the block scheme of fig. 5.37 is apparently $R H_{\infty}$ admissible for $P_{O E}(s)$, provided that it is well defined, that is, provided that matrix $I+D_{q} D_{22}$ is nonsingular.

### 5.5 The partial information problem

The partial information problem in $R H_{\infty}$ is discussed now. As in Section 4.4, only the output variable $y$ can be measured: consistently, the system under control is described by

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{5.121}\\
z & =C_{1} x+D_{12} u  \tag{5.122}\\
y & =C_{2} x+D_{21} w \tag{5.123}
\end{align*}
$$

while its transfer function is denoted by $P(s)$. Notice that system (5.121)-(5.123) is less general than system (5.33)-(5.35) however, the latter can be redrawn to the former by exploiting Remark 5.24. The solution of Problem 5.1 relative to system (5.121)-(5.123) will be presented under the following assumptions.


Figure 5.35: The set $\mathcal{F}_{\infty \gamma r}$ of the controllers $K_{F C r}(s)$


Figure 5.36: The generic admissible controller for $P_{F C}(s)$

Assumption 5.6 The pair $\left[\left(A-B_{2} D_{12}^{\prime} C_{1}\right),\left(I-D_{12} D_{12}^{\prime}\right) C_{1}\right]$ is detectable and the pair $\left[\left(A-B_{1} D_{21}^{\prime} C_{2}\right), B_{1}\left(I-D_{21}^{\prime} D_{21}\right)\right]$ is stabilizable.

Assumption 5.7 The pair $\left(A, B_{2}\right)$ is stabilizable and the pair $\left(A, C_{2}\right)$ is detectable.
Assumption $5.8 D_{12}^{\prime} D_{12}=I$.
Assumption $5.9 D_{21} D_{21}^{\prime}=I$.
The result in Theorem 5.4 below makes reference to the block-scheme of fig. 5.38 relative to which the transfer function from $w$ to $z$ is denoted by $T(z, w ; s)$, while $K(s)$ is a generic $R H_{\infty}$ admissible controller for $P(s)$.
Theorem 5.4 (Partial information problem) Consider Problem 5.1 relative to system (5.121)-(5.123). Then, under Assumptions 5.6-5.9, it has the solution
a1) Existence of the symmetric, positive semidefinite and stabilizing solution $P_{\infty}$ of the Riccati equation (in the unknown $P$ )

$$
\begin{equation*}
0=P A_{c}+A_{c}^{\prime} P-P\left(B_{2} B_{2}^{\prime}-\gamma^{-2} B_{1} B_{1}^{\prime}\right) P+C_{1 c}^{\prime} C_{1 c} \tag{5.124}
\end{equation*}
$$

that is, such that

$$
\begin{equation*}
A_{c c}:=A_{c}-B_{2} B_{2}^{\prime} P_{\infty}+\gamma^{-2} B_{1} B_{1}^{\prime} P_{\infty} \tag{5.125}
\end{equation*}
$$



Figure 5.37: The controller structure when $D_{22} \neq 0$


Figure 5.38: The partial information problem
is stable. In the above equations

$$
\begin{equation*}
A_{c}:=A-B_{2} D_{12}^{\prime} C_{1}, C_{1 c}:=\left(I-D_{12} D_{12}^{\prime}\right) C_{1} \tag{5.126}
\end{equation*}
$$

a2) Existence of the symmetric, positive semidefinite and stabilizing solution $\Pi_{\infty}$ of the Riccati equation (in the unknown $\Pi$ )

$$
\begin{equation*}
0=\Pi A_{f}^{\prime}+A_{f} \Pi-\Pi\left(C_{2}^{\prime} C_{2}-\gamma^{-2} C_{1}^{\prime} C_{1}\right) \Pi+B_{1 f} B_{1 f}^{\prime} \tag{5.127}
\end{equation*}
$$

that is, such that

$$
\begin{equation*}
A_{f f}:=A_{f}-\Pi_{\infty} C_{2}^{\prime} C_{2}+\gamma^{-2} \Pi_{\infty} C_{1}^{\prime} C_{1} \tag{5.128}
\end{equation*}
$$

is stable. In the above equations

$$
\begin{equation*}
A_{f}:=A-B_{1} D_{21}^{\prime} C_{2}, \quad B_{1 f}:=B_{1}\left(I-D_{21}^{\prime} D_{21}\right) \tag{5.129}
\end{equation*}
$$

a3)

$$
\begin{equation*}
r_{s}\left(P_{\infty} \Pi_{\infty}\right)<\gamma^{2} \tag{5.130}
\end{equation*}
$$

b) The set $\mathcal{F}_{\infty \gamma r}$ of the controllers $K_{r}(s)$ is defined by the diagram of fig. 5.39, where

$$
S_{\infty}(s):=\left[\begin{array}{c|cc}
A_{c c}+Z_{\infty} L_{\infty}\left(C_{2}+\gamma^{-2} D_{21} B_{1}^{\prime} P_{\infty}\right) & -Z_{\infty} L_{\infty} & +Z_{\infty} B_{2 \infty} \\
\hline F_{\infty} & 0 & I \\
-C_{2}-\gamma^{-2} D_{21} B_{1}^{\prime} P_{\infty} & I & 0
\end{array}\right]
$$



Figure 5.39: The set $\mathcal{F}_{\infty \gamma r}$ of the controllers $K_{\text {FIr }}(s)$


Figure 5.40: The equivalent output estimation problem
where

$$
\begin{align*}
F_{\infty} & :=-B_{2}^{\prime} P_{\infty}-D_{12}^{\prime} C_{1}  \tag{5.131}\\
L_{\infty} & :=-\Pi_{\infty} C_{2}^{\prime}-B_{1} D_{21}^{\prime}  \tag{5.132}\\
Z_{\infty} & :=\left(I-\gamma^{-2} \Pi_{\infty} P_{\infty}\right)^{-1}  \tag{5.133}\\
B_{2 \infty} & :=B_{2}+\gamma^{-2} \Pi_{\infty} C_{1}^{\prime} D_{12}
\end{align*}
$$

and $Q(s):=\Sigma\left(A_{q}, B_{q}, C_{q}, D_{q}\right)$ with $A_{q}$ stable and $\|Q(s)\|_{\infty}<\gamma$.
Proof First observe that Assumption 5.7 coincides with the necessary condition of Remark 5.6.

Sufficiency of parts a1) - a3) and part b) Assume that points (a1)-(a3) hold true and define the variables $v$ and $q$ as

$$
\begin{align*}
w & :=r+\gamma^{-2} B_{1}^{\prime} P_{\infty} x  \tag{5.134}\\
q & :=u-F_{\infty} x \tag{5.135}
\end{align*}
$$

By utilizing these equations into eqs. (5.121) and (5.123) one obtains the system (see fig. 5.40)

$$
P_{t}(s):=\left[\begin{array}{c|cc}
A_{t} & B_{1} & B_{2}  \tag{5.136}\\
\hline-F_{\infty} & 0 & I \\
C_{2 t} & D_{21} & 0
\end{array}\right]
$$

where

$$
\begin{align*}
A_{t} & :=A+\gamma^{-2} B_{1} B_{1}^{\prime} P_{\infty}  \tag{5.137}\\
C_{2 t} & :=C_{2}+\gamma^{-2} D_{21} B_{1}^{\prime} P_{\infty} \tag{5.138}
\end{align*}
$$

It is easy to ascertain that Lemma E. 4 can be applied to the couple of systems shown in fig. 5.38 and 5.40 , so that this part of the proof can be carried over by making reference to the system in fig. 5.40 only.

Observe that system $P_{t}(s)$ has the same structure as system $P_{O E}(s)$ dealt with in Section 5.4, where the output estimation problem was discussed. When dealing with system $P_{t}(s)$ Assumptions 5.3-5.5 become: $(\alpha 1)$ the pair $\left(A_{t}-B_{1} D_{21}^{\prime} C_{2 t}, B_{1 f}\right)$ is stabilizable; $(\alpha 2)$ the pair $\left(A_{t}, C_{2 t}\right)$ is detectable; $(\alpha 3) D_{21} D_{21}^{\prime}=I ;(\alpha 4)$ the matrix $A_{t}+B_{2} F_{\infty}$ is stable. The stability of this last matrix trivially follows from eq. $(5.125),(5.126),(5.131),(5.137)$. Condition $(\alpha 3)$ is Assumption 5.9.

Condition ( $\alpha 1$ ) is immediately derived by acknowledging that $A_{t}-B_{1} D_{21}^{\prime} C_{2 t}=$ $A_{f}+\gamma^{-2} B_{1 f} B_{1}^{\prime} P_{\infty}$, so that such a condition is equivalent to the stabilizability of the pair $\left(A_{f}+\gamma^{-2} B_{1 f} B_{1}^{\prime} P_{\infty}, B_{1 f}\right)$, which, in turn, is implied by Assumption 5.6, thanks to the invariance of the stabilizability property with respect to state feedback. The discussion on condition ( $\alpha 2$ ) is postponed.

Therefore, the sufficiency part of the theorem and point (b) are proved once the existence is ascertained of a symmetric and positive semidefinite solution $\Pi_{t \infty}$ of the Riccati equation (in the unknown $\Pi_{t}$ ) which is relevant to the output estimation problem for system $P_{t}(s)$, namely, the equation

$$
\begin{equation*}
0=A_{t t} \Pi_{t}+\Pi_{t} A_{t t}^{\prime}+\Pi_{t}\left(\gamma^{-2} F_{\infty}^{\prime} F_{\infty}-C_{2 t}^{\prime} C_{2 t}\right) \Pi_{t}+B_{1 f} B_{1 f}^{\prime} \tag{5.139}
\end{equation*}
$$

where

$$
\begin{equation*}
A_{t t}:=A_{t}-B_{1} D_{21}^{\prime} C_{2 t} \tag{5.140}
\end{equation*}
$$

Moreover, the solution $\Pi_{t \infty}$ must be stabilizing, i. e. such that the matrix $A_{t t}+$ $\Pi_{t \infty}\left(\gamma^{-2} F_{\infty}^{\prime} F_{\infty}-C_{2 t}^{\prime} C_{2 t}\right)$ is stable. To this aim, consider the Hamiltonian matrix associated with the Riccati equation (5.139)

$$
J_{t \infty}:=\left[\begin{array}{cc}
A_{t t}^{\prime} & \gamma^{-2} F_{\infty}^{\prime} F_{\infty}-C_{2 t}^{\prime} C_{2 t} \\
-B_{1 f} B_{1 f}^{\prime} & -A_{t t}
\end{array}\right]
$$

and the one associated with the Riccati equation (5.127)

$$
J_{\infty}:=\left[\begin{array}{cc}
A_{f}^{\prime} & \gamma^{-2} C_{1}^{\prime} C_{1}-C_{2}^{\prime} C_{2} \\
-B_{1 f} B_{1 f}^{\prime} & -A_{f}
\end{array}\right]
$$

By exploiting the fact that $P_{\infty}$ solves eq. (5.124), it is not difficult to verify that the two Hamiltonian matrices $J_{t \infty}$ and $J_{\infty}$ are similar. Indeed, it results

$$
J_{t \infty}=T J_{\infty} T^{-1}
$$

with

$$
T:=\left[\begin{array}{cc}
I & -\gamma^{-2} P_{\infty} \\
0 & I
\end{array}\right]
$$

Being $\Pi_{\infty}$ a solution of eq. (5.127), one has

$$
J_{\infty} \operatorname{Im}\left[\left[\begin{array}{c}
I \\
\Pi_{\infty}
\end{array}\right]\right] \subseteq \operatorname{Im}\left[\left[\begin{array}{c}
I \\
\Pi_{\infty}
\end{array}\right]\right]
$$

so that

$$
J_{t \infty} T \operatorname{Im}\left[\left[\begin{array}{c}
I \\
\Pi_{\infty}
\end{array}\right]\right]=T J_{\infty} \operatorname{Im}\left[\left[\begin{array}{c}
I \\
\Pi_{\infty}
\end{array}\right]\right] \subseteq T \operatorname{Im}\left[\left[\begin{array}{c}
I \\
\Pi_{\infty}
\end{array}\right]\right]
$$

Therefore, the subspace

$$
\mathcal{S}:=T \operatorname{Im}\left[\left[\begin{array}{c}
I \\
\Pi_{\infty}
\end{array}\right]\right]=\operatorname{Im}\left[\left[\begin{array}{c}
I-\gamma^{-2} P_{\infty} \Pi_{\infty} \\
\Pi_{\infty}
\end{array}\right]\right]
$$

is $J_{t \infty}$-invariant. It is also complementary to $\operatorname{Im}\left[\left[\begin{array}{ll}0 & I\end{array}\right]^{\prime}\right]$ because condition (a3) guarantees that the matrix $I-\gamma^{-2} P_{\infty} \Pi_{\infty}$ is nonsingular (see the proof of Lemma B.11). Finally, the subspace $\mathcal{S}$ is generated by the (generalized) eigenvectors of $J_{t \infty}$ relative to the eigenvalues with negative real parts. Indeed, being $\Pi_{\infty}$ a stabilizing solution,

$$
J_{\infty} \operatorname{Im}\left[\left[\begin{array}{c}
I \\
\Pi_{\infty}
\end{array}\right]\right]=\operatorname{Im}\left[\left[\begin{array}{c}
I \\
\Pi_{\infty}
\end{array}\right]\right] \Lambda
$$

with $\Lambda$ stable, so that $J_{t \infty} \mathcal{S}=\mathcal{S} \Lambda$. Therefore, in view of Lemma C.2, it can be concluded that there exists a stabilizing solution $\Pi_{t \infty}$ of eq. (5.139) which is given by

$$
\begin{equation*}
\Pi_{t \infty}=\Pi_{\infty}\left(I-\gamma^{-2} P_{\infty} \Pi_{\infty}\right)^{-1} \tag{5.141}
\end{equation*}
$$

Moreover, such a solution turns out to be positive semidefinite thanks to Lemma B.11.
It is left to prove condition ( $\alpha 2$ ), namely, the detectability of the pair ( $A_{t}, C_{2 t}$ ). Remark 5.17 applied to eq. (5.139) leads to the conclusion that matrix $A_{t t}-\Pi_{t} C_{2 t}^{\prime} C_{2 t}$ is stable, which, in turn, implies that the pair $\left(A_{t t}, C_{2 t}\right)$ is detectable. By recalling the definition of the matrix $A_{t t}$ (see eq. (5.140)), the detectability of the pair ( $A_{t}, C_{2 t}$ ), i. e. condition ( $\alpha 2$ ), follows too, since the stabilizability property of the pair ( $A_{t}^{\prime}, C_{2 t}^{\prime}$ ) is invariant with respect to the state feedback.

Finally, the form of the set $\mathcal{F}_{\infty \gamma r}$ given in the statement follows by applying Theorem 5.3 to system $P_{t}(s)$.

Necessity of part a) With reference to fig. 5.38, assume that there exists a controller $K(s)$ admissible in $R H_{\infty}$ for system $P(s)$ such that $\|T(z, w ; s)\|_{\infty}<\gamma$. Then, the controller

$$
K_{F I}(s):=K(s)\left[\begin{array}{ll}
C_{2} & D_{21}
\end{array}\right]
$$

is admissible in $R H_{\infty}$ for the system $\bar{P}_{F I}(s)$ which is the system defined by eqs. (5.121) and (5.122) with the output vector constituted by $\left[x^{\prime} w^{\prime}\right]^{\prime}$. Moreover, $\|T(z, w ; s)\|_{\infty}<$ $\gamma$ also for system $\bar{P}_{F I}(s)$. This system coincides with the system $P_{F I}(s)$ considered in Section 5.3. The Assumptions 5.1 and 5.2 which are required for the result relevant to the full information problem (Theorem 5.2), are verified since they constitute a subset of Assumptions 5.6-5.9. Therefore, Theorem 5.2 guarantees the existence of the solution of eq. (5.124) endowed with all the properties specified in the statement of the present theorem. The request of the existence of the solution of eq. (5.127) can be proved in a similar manner by making reference to the system in fig. 5.41 , where $\hat{P}(s):=P^{\prime}(s)$ and $\hat{K}(s):=K^{\prime}(s)$. Thus, system $\hat{P}(s)$ is described by

$$
\begin{align*}
\dot{\xi} & =A^{\prime} \xi+C_{1}^{\prime} \hat{z}+C_{2}^{\prime} \hat{y}  \tag{5.142}\\
\hat{w} & =B_{1}^{\prime} \xi+D_{21}^{\prime} \hat{y}  \tag{5.143}\\
\hat{u} & =B_{2}^{\prime} \xi+D_{12}^{\prime} \hat{z} \tag{5.144}
\end{align*}
$$

From Lemma E.1, $\hat{K}(s)$ is admissible in $R H_{\infty}$ for $\hat{P}(s)$ and $\|T(\hat{w}, \hat{z} ; s)\|_{\infty}<\gamma$ because $K(s)$ is admissible in $R H_{\infty}$ for $P(s)$ and $\|T(z, w ; s)\|_{\infty}<\gamma$. The above exploited argument can be applied in the very same way leading to the conclusion that the controller

$$
\hat{K}_{F I}(s):=\hat{K}(s)\left[\begin{array}{ll}
B_{2}^{\prime} & D_{12}^{\prime}
\end{array}\right]
$$



Figure 5.41: The transposed partial information problem
is admissible in $R H_{\infty}$ for system $\hat{P}_{F I}(s)$ (described by eqs. (5.142) and (5.143) and having the vector $\left[\xi^{\prime} \hat{z}^{\prime}\right]^{\prime}$ as measurable output) and is such that $\|T(\hat{w}, \hat{z} ; s)\|_{\infty}<\gamma$. System $\hat{P}_{F I}(s)$ has the same structure as system $P_{F I}(s)$ which has been defined for the full information problem (see Section 5.3). Since the assumptions required by Theorem 5.2 to be applied to system $\hat{P}_{F I}(s)$ are a subset of Assumptions 5.6-5.9, it is possible to claim that there exists a solution of eq. (5.127) which possesses the properties requested by the statement of the present theorem.

Further, it is easy to check that also the conditions for Lemma E. 4 to be applied are verified, so that, with reference to the block-scheme of fig. 5.40 , it is possible to conclude that the controller $K(s)$ is admissible in $R H_{\infty}$ for system $P_{t}(s)$ and $\|T(q, r ; s)\|_{\infty}<\gamma$ (recall the definition of $P_{t}(s)$ given in eqs. (5.136)-(5.138)). As already said, system $P_{t}(s)$ has the same structure as system $P_{O E}(s)$ which was introduced in Section 5.4 when dealing with the output estimation problem. Assumptions 5.3 - 5.5 which are relevant to Theorem 5.3, are now shown to be verified. These assumptions, if expressed in terms of system $P_{t}(s)$, are precisely those which have been denoted as $(\alpha 1),(\alpha 2),(\alpha 3)$ and $(\alpha 4)$ in the sufficiency part of the proof.

First observe that condition ( $\alpha 2$ ) (detectability of the pair $\left(A_{t}, C_{2 t}\right)$ ) is obviously satisfied. In fact, it is apparently necessary in view of the already established stability of the system in fig. 5.40 (again, recall the structure of system $P_{t}(s)$ given in eq. (5.136)). Second, in the sufficiency part of the proof it has already been shown that condition ( $\alpha 1$ ) is equivalent to Assumption 5.6. Condition ( $\alpha 3$ ) is straightforward. Finally, condition ( $\alpha 4$ ) precisely amounts to requiring the stability of matrix $A_{c c}$ (see eq. (5.125)).

Thus, the conditions for Theorem 5.3 to be applied to system $P_{t}(s)$ are satisfied. Therefore, this theorem guarantees the existence of a positive semidefinite and stabilizing solution $\Pi_{t \infty}$ of the Riccati equation (5.139). By applying in reverse order what has been shown in passing from $\Pi_{\infty}$ (a solution of eq. (5.127)) to $\Pi_{t \infty}$ (a solution of eq. (5.139)), it follows

$$
\begin{equation*}
\Pi_{\infty}:=\Pi_{t \infty}\left(I+\gamma^{-2} P_{\infty} \Pi_{t \infty}\right)^{-1} \tag{5.145}
\end{equation*}
$$

Note that the existence of the inverse in eq. (5.145) is guaranteed by the fact that matrix $P_{\infty} \Pi_{t \infty}$ has nonnegative eigenvalues, thanks to Lemma B.10. From eq. (5.145) it follows

$$
\begin{equation*}
\Pi_{\infty}=\left(I-\gamma^{-2} \Pi_{\infty} P_{\infty}\right) \Pi_{t \infty} \tag{5.146}
\end{equation*}
$$

so that

$$
\Pi_{t \infty}=\left(I-\gamma^{-2} \Pi_{\infty} P_{\infty}\right)^{-1} \Pi_{\infty}
$$

since the existence of the inverse in this last equation can easily be proved. In fact, let, by contradiction,

$$
\begin{equation*}
x^{\prime}\left(I-\gamma^{-2} \Pi_{\infty} P_{\infty}\right)=0, \quad x \neq 0 \tag{5.147}
\end{equation*}
$$

so that, from eq. (5.146), also $x^{\prime} \Pi_{\infty}=0$. But this would imply, in view of eq. (5.147), $x=0$. By recalling that $\Pi_{t \infty} \geq 0$, one can conclude that condition ( $a 3$ ) in the statement is verified, thanks to eq. (5.141) and Lemma B.11.

The following result, which constitutes an interesting side-product of Theorem 5.4 , can be stated in view of what has been presented in connection with the full information problem (Corollary 5.1) and the output estimation problem (Corollary 5.2).

Corollary 5.3 Assume that Assumptions 5.8 and 5.9 hold. If, for a given positive $\gamma$, there exist the symmetric, positive semidefinite and stabilizing solutions of the Riccati equations (5.124) and (5.127) and such solutions satisfy eq. (5.130), then there exists a $R H_{\infty}$ admissible controller $K(s)$ for $P(s)$, corresponding to which $\|T(z, w ; s)\|_{\infty}<$ $\gamma$. Moreover, the set of controllers described in point (b) of Theorem 5.4 constitutes the set $\mathcal{F}_{\infty \gamma r}$.

Example 5.7 Consider system (5.121) - (5.123) with

$$
\begin{gathered}
A=\left[\begin{array}{ll}
0 & 1 \\
2 & 1
\end{array}\right], \quad B_{1}=\left[\begin{array}{l}
0 \\
0
\end{array}\right], \quad B_{2}=\left[\begin{array}{l}
0 \\
1
\end{array}\right] \\
C_{1}=\left[\begin{array}{ll}
0 & 0
\end{array}\right], \quad C_{2}=\left[\begin{array}{ll}
-1 & 1
\end{array}\right], \quad D_{12}=D_{21}=1
\end{gathered}
$$

Because of the particular form of matrices $B_{1}$ and $C_{1}$, the Riccati equations (5.124) and (5.127) do not depend on $\gamma$. The stabilizing solutions of these equations are

$$
P_{\infty}=\left[\begin{array}{ll}
4 & 4 \\
4 & 4
\end{array}\right], \quad \Pi_{\infty}=\left[\begin{array}{cc}
4 & 8 \\
8 & 16
\end{array}\right]
$$

respectively. Moreover, $r_{s}\left(P_{\infty} \Pi_{\infty}\right)=144$ so that, according to Corollary 5.3, an admissible controller in $R H_{\infty}$ exists for the system if $\gamma>12$ (recall eq. (5.130)).

Now consider the nominal plant $G_{n}(s):=(s-1) /\left(s^{2}-s-2\right)$ subject to perturbations of the form (5.10),(5.11). Then, what is the maximum value of $\alpha$ corresponding to which a stabilizing controller exists? This question can be answered by applying Theorem 5.1 to an enlarged plant which coincides with the controlled system considered in the present example. Thus, a stabilizing controller exists only provided that $\alpha \leq 1 / 12$.

Remark 5.24 (Loop shifting) Theorem 5.4 can be applied to systems exhibiting a more general structure than that of system (5.121) - (5.123). Indeed, consider system $P^{(1)}(s)$ given by

$$
P^{(1)}(s):=\left[\begin{array}{c|cc}
A^{(1)} & B_{1}^{(1)} & B_{2}^{(1)} \\
\hline C_{1}^{(1)} & D_{11}^{(1)} & D_{12}^{(1)} \\
C_{2}^{(1)} & D_{21}^{(1)} & D_{22}^{(1)}
\end{array}\right]
$$

and depicted in fig. 5.42. Assume that the variables $w^{(1)}, u^{(1)}, z^{(1)}, y^{(1)}$ of such a system have dimension $q, m, r, p$, respectively, and suppose that $\operatorname{rank}\left[D_{12}^{(1)}\right]=m, \operatorname{rank}\left[D_{21}^{(1)}\right]=p$.

Then, if one performs the 6 operations described below (the effects of which can be easily understood by making reference to fig. 5.42), it is generically possible to successively transform system $P^{(1)}(s)$ into systems

$$
P^{(i)}(s):=\left[\begin{array}{c|cc}
A^{(i)} & B_{1}^{(i)} & B_{2}^{(i)} \\
\hline C_{1}^{(i)} & D_{11}^{(i)} & D_{12}^{(i)} \\
C_{2}^{(i)} & D_{21}^{(i)} & D_{22}^{(i)}
\end{array}\right]
$$



Figure 5.42: The loop shifting
which are characterized by the inputs $w^{(i)}$ and $u^{(i)}$ and the outputs $z^{(i)}$ and $y^{(i)}$, in such a way that

$$
P^{(7)}(s):=\left[\begin{array}{c|cc}
A^{(7)} & B_{1}^{(7)} & B_{2}^{(7)} \\
\hline C_{1}^{(7)} & 0 & D_{12}^{(7)} \\
C_{2}^{(7)} & D_{21}^{(7)} & 0
\end{array}\right]
$$

with $D_{12}^{(7) \prime} D_{12}^{(7)}=I, D_{21}^{(7)} D_{21}^{(7) \prime}=I$.
Operation 1 : Let

$$
\begin{aligned}
y^{(2)} & :=y^{(1)}-D_{22}^{(1)} u^{(1)} \\
w^{(2)} & :=w^{(1)}, \quad u^{(2)}:=u^{(1)}, \quad z^{(2)}:=z^{(1)}
\end{aligned}
$$

so that

$$
\begin{aligned}
& A^{(2)}=A^{(1)}, \quad B_{i}^{(2)}=B_{i}^{(1)}, \quad i=1,2 \\
& C_{i}^{(2)}=C_{i}^{(1)}, \quad i=1,2 \\
& D_{11}^{(2)}=D_{11}^{(1)}, \quad D_{12}^{(2)}=D_{12}^{(1)}, \quad D_{21}^{(2)}=D_{21}^{(1)}
\end{aligned}
$$

and

$$
D_{22}^{(2)}=0
$$

Operation 2 : Find four matrices $V, V^{\perp}, W, W^{\perp}$ with orthonormal columns and such that

$$
\begin{aligned}
\operatorname{Im}[V] & =\operatorname{Im}\left[D_{12}^{(2)}\right], \quad \operatorname{Im}\left[V^{\perp}\right]=\operatorname{Im}\left[D_{12}^{(2)}\right]^{\perp} \\
\operatorname{Im}[W] & =\operatorname{Im}\left[D_{21}^{(2) \prime}\right], \quad \operatorname{Im}\left[W^{\perp}\right]=\operatorname{Im}\left[D_{21}^{(2) \prime}\right]^{\perp}
\end{aligned}
$$

Such matrices do exist thanks to the assumptions $\operatorname{rank}\left[D_{12}^{(1)}\right]=m, \operatorname{rank}\left[D_{21}^{(1)}\right]=p$ and the fact that $D_{12}^{(2)}=D_{12}^{(1)}$ and $D_{21}^{(2)}=D_{21}^{(1)}$. Obviously, if $D_{12}^{(2)}$ is a square matrix (hence nonsingular), $V^{\perp}=\emptyset$ and $V=I$. A similar comment is in order if $D_{21}^{(2)}$ is a square matrix (hence nonsingular).

Then it follows

$$
\begin{aligned}
& V_{z} D_{12}^{(2)}:=\left[\begin{array}{c}
V^{\perp^{\prime}} \\
V^{\prime}
\end{array}\right] D_{12}^{(2)}=\left[\begin{array}{c}
0 \\
S_{u}
\end{array}\right] \\
& D_{21}^{(2)} V_{w}:=D_{21}^{(2)}\left[\begin{array}{ll}
W^{\perp} & W
\end{array}\right]=\left[\begin{array}{ll}
0 & S_{y}
\end{array}\right]
\end{aligned}
$$

Observe that matrices $V_{w}$ and $V_{z}$ are orthogonal and matrices $S_{u}$ and $S_{y}$ are nonsingular. Then let

$$
\begin{aligned}
w^{(3)} & :=V_{w}^{\prime} w^{(2)}, u^{(3)}:=S_{u} u^{(2)} \\
z^{(3)} & :=V_{z} z^{(2)}, y^{(3)}:=S_{y}^{-1} y^{(2)}
\end{aligned}
$$

so that

$$
\begin{aligned}
& A^{(3)}=A^{(2)}, \quad B_{1}^{(3)}=B_{1}^{(2)} V_{w}, \quad B_{2}^{(3)}=B_{2}^{(2)} S_{u}^{-1} \\
& C_{1}^{(3)}=V_{z} C_{1}^{(2)}, \quad C_{2}^{(3)}=S_{y}^{-1} C_{2}^{(2)}, \quad D_{11}^{(3)}=V_{z} D_{11}^{(2)} V_{w} \\
& D_{22}^{(3)}=0
\end{aligned}
$$

and

$$
D_{12}^{(3)}=V_{z} D_{12}^{(2)} S_{u}^{-1}=\left[\begin{array}{l}
0 \\
I
\end{array}\right], \quad D_{21}^{(3)}=S_{y}^{-1} D_{21}^{(2)} V_{w}=\left[\begin{array}{ll}
0 & I
\end{array}\right]
$$

Note that, thanks to the orthogonality of matrices $V_{z}$ and $V_{w}$, the transfer functions from $w^{(3)}$ to $z^{(3)}$ and from $w^{(2)}$ to $z^{(2)}$ have the same $R H_{\infty}$ norm.
Operation 3: Consider the following partition of the $r \times q$ matrix $D_{11}^{(3)}$

$$
D_{11}^{(3)}:=\left[\begin{array}{ll}
D_{1111} & D_{1112} \\
D_{1121} & D_{1122}
\end{array}\right]
$$

where the submatrix $D_{1122}$ is $m \times p$ (recall that the assumptions on the rank of matrices $D_{12}^{(1)}$ and $D_{21}^{(1)}$ imply $r \geq m$ and $q \geq p$ ). Let $\Lambda_{\infty}$ be the matrix (whose existence is guaranteed by Theorem F. 1 which supplies its expression, too) which minimizes $\|\Delta(\Lambda)\|$, where

$$
\Delta(\Lambda):=\left[\begin{array}{cc}
D_{1111} & D_{1112} \\
D_{1121} & D_{1122}+\Lambda
\end{array}\right]
$$

If

$$
\alpha:=\max \left(\left\|\left[\begin{array}{l}
D_{1111} \\
D_{1121}
\end{array}\right]\right\|,\left\|\left[\begin{array}{l}
D_{1111}^{\prime} \\
D_{1112}^{\prime}
\end{array}\right]\right\|\right)
$$

then, thanks to the above quoted theorem, $\left\|\Delta\left(\Lambda_{\infty}\right)\right\|=\alpha$. By defining

$$
\begin{aligned}
w^{(4)} & :=w^{(3)}, \quad u^{(4)}:=u^{(3)}-\Lambda_{\infty} y^{(3)} \\
z^{(4)} & :=z^{(3)}, \quad y^{(4)}:=y^{(3)}
\end{aligned}
$$

one obtains

$$
\begin{aligned}
& A^{(4)}=A^{(3)}+B_{2}^{(3)} \Lambda_{\infty} C_{2}^{(3)} \\
& B_{1}^{(4)}=B_{1}^{(3)}+B_{2}^{(3)} \Lambda_{\infty} D_{21}^{(3)}, B_{2}^{(4)}=B_{2}^{(3)} \\
& C_{1}^{(4)}=C_{1}^{(3)}+D_{12}^{(3)} \Lambda_{\infty} C_{2}^{(3)}, \quad C_{2}^{(4)}=C_{2}^{(3)} \\
& D_{11}^{(4)}=\Delta\left(\Lambda_{\infty}\right) \\
& D_{12}^{(4)}=D_{12}^{(3)}=\left[\begin{array}{c}
0 \\
I
\end{array}\right], D_{21}^{(4)}=D_{21}^{(3)}=\left[\begin{array}{ll}
0 & I
\end{array}\right] \\
& D_{22}^{(4)}=0
\end{aligned}
$$

Now consider the feedback connection of system $P^{(3)}(s)$ with the controller defined by $u_{L}^{(3)}=$ $\Lambda(s) y_{L}^{(3)}$, with $\lim _{\omega \rightarrow \infty} \Lambda(j \omega)=\Lambda_{\infty}$. Then, the scalar $\alpha$ constitutes a lower bound to the $R H_{\infty}$ norm of the transfer function from $w^{(3)}$ to $z^{(3)}$. Indeed, as $\omega \rightarrow \infty$, the limiting value of such a function is precisely

$$
D_{11}^{(4)}=D_{11}^{(3)}+D_{12}^{(3)} \Lambda_{\infty} D_{21}^{(3)}=\Delta\left(\Lambda_{\infty}\right)
$$

Operation 4 : Define the variables $z^{(5)}$ and $w^{(5)}$ by means of the equation

$$
\left[\begin{array}{c}
z^{(5)} \\
w^{(4)}
\end{array}\right]:=\Gamma_{1} \Theta \Gamma_{2}\left[\begin{array}{c}
w^{(5)} \\
z^{(4)}
\end{array}\right]
$$

where

$$
\Gamma_{1}:=\left[\begin{array}{cc}
I & 0 \\
0 & \gamma^{-1} I
\end{array}\right], \quad \Gamma_{2}:=\left[\begin{array}{cc}
\gamma I & 0 \\
0 & I
\end{array}\right]
$$

while

$$
\Theta:=\left[\begin{array}{ll}
\Theta_{11} & \Theta_{12} \\
\Theta_{21} & \Theta_{22}
\end{array}\right]
$$

is any matrix such that: (i) $\Theta^{\prime} \Theta=\Theta^{\prime}=I$; (ii) $\Theta_{12}$ and $\Theta_{21}$ are square matrices of dimension $r$ and $q$, respectively; (iiii) $0<\gamma^{-1}<\alpha^{-1}$. In particular, from these features it follows that $\left\|\Theta_{22}\right\| \leq\|\Theta\|=1$ (recall Lemma 2.16), so that the matrices $\Psi_{1}:=(\gamma I-$ $\left.D_{11}^{(4)} \Theta_{22}\right)^{-1}$ and $\Psi_{2}:=\left(\gamma I-\Theta_{22} D_{11}^{(4)}\right)^{-1}$ are well defined because $\left\|D_{11}^{(4)}\right\|=\alpha$ (recall also Lemma 2.18, point (2) and Lemma 2.21). Further let

$$
u^{(5)}:=u^{(4)}, y^{(5)}:=y^{(4)}
$$

In so doing, one obtains

$$
\begin{aligned}
& A^{(5)}=A^{(4)}+B_{1}^{(4)} \Psi_{2} \Theta_{22} C_{1}^{(4)} \\
& B_{1}^{(5)}=\gamma B_{1}^{(4)} \Psi_{2} \Theta_{21}, \quad B_{2}^{(5)}=B_{2}^{(4)}+B_{1}^{(4)} \Psi_{2} \Theta_{22} D_{12}^{(4)} \\
& C_{1}^{(5)}=\gamma \Theta_{12} \Psi_{1} C_{1}^{(4)}, \quad C_{2}^{(5)}=C_{2}^{(4)}+D_{21}^{(4)} \Psi_{2} \Theta_{22} C_{1}^{(4)} \\
& D_{11}^{(5)}=\gamma\left(\Theta_{11}+\Theta_{12} D_{11}^{(4)} \Psi_{2} \Theta_{21}\right), \quad D_{22}^{(5)}=D_{21}^{(4)} \Psi_{2} \Theta_{22} D_{12}^{(4)} \\
& D_{12}^{(5)}=\gamma \Theta_{12} \Psi_{1} D_{12}^{(4)}, \quad D_{21}^{(5)}=\gamma D_{21}^{(4)} \Psi_{2} \Theta_{21}
\end{aligned}
$$

Notice that if the feedback connection of system $P^{(4)}(s)$ with the controller described by $u_{L}^{(4)}=\Lambda(s) y_{L}^{(4)}$ results in a stable system and the $R H_{\infty}$ norm of the transfer function $T\left(z^{(4)}, w^{(4)} ; s\right)$ from the input $w^{(4)}$ to the output $z^{(4)}$ is less than $\gamma$, then also the feedback connection of system $P^{(5)}(s)$ with the controller described by by $u_{L}^{(5)}=\Lambda(s) y_{L}^{(5)}$ results in a stable system and the $R H_{\infty}$ norm of the transfer function $T\left(z^{(5)}, w^{(5)} ; s\right)$ from the input $w^{(5)}$ to the output $z^{(5)}$ is less than $\gamma$. In fact, notice that

$$
\Gamma_{1} \Theta \Gamma_{2}=\left[\begin{array}{cc}
\gamma \Theta_{11} & \Theta_{12} \\
\Theta_{21} & \gamma^{-1} \Theta_{22}
\end{array}\right]
$$

and $\left\|\gamma^{-1} \Theta_{22}\right\| \leq \gamma^{-1}$, so that it is possible to apply Theorem 5.1 and conclude about stability. As for the transfer functions norms, observe that, from the definition of the new variables $z^{(5)}$ and $w^{(5)}$, it follows

$$
\left[\begin{array}{c}
z^{(5)} \\
\gamma w^{(4)}
\end{array}\right]=\Theta\left[\begin{array}{c}
\gamma w^{(5)} \\
z^{(4)}
\end{array}\right]
$$

so that, thanks to the properties of $\Theta$,

$$
\left\|z^{(5)}\right\|_{2}^{2}+\gamma^{2}\left\|w^{(4)}\right\|_{2}^{2}=\left\|z^{(4)}\right\|_{2}^{2}+\gamma^{2}\left\|w^{(5)}\right\|_{2}^{2}
$$

that is

$$
\left\|z^{(5)}\right\|_{2}^{2}-\gamma^{2}\left\|w^{(5)}\right\|_{2}^{2}=\left\|z^{(4)}\right\|_{2}^{2}-\gamma^{2}\left\|w^{(4)}\right\|_{2}^{2}
$$

The left (resp. right) hand side of this equation is negative if and only if $\left\|T\left(z^{(5)}, w^{(5)} ; s\right)\right\|_{\infty}<$ $\gamma$ (resp. $\left\|T\left(z^{(4)}, w^{(4)} ; s\right)\right\|_{\infty}<\gamma$ ) (see Theorem 2.12), so that the $R H_{\infty}$ norm of the first transfer function is less than $\gamma$ if and only if the norm of the second one is such. Now let

$$
\begin{aligned}
& \Theta_{11}:=-\gamma^{-1} D_{11}^{(4)}, \quad \Theta_{12}:=\left(I-\gamma^{-2} D_{11}^{(4)} D_{11}^{(4) \prime}\right)^{1 / 2} \\
& \Theta_{21}:=\left(I-\gamma^{-2} D_{11}^{(4) \prime} D_{11}^{(4)}\right)^{1 / 2}, \quad \Theta_{22}:=\gamma^{-1} D_{11}^{(4) \prime}
\end{aligned}
$$

By applying Lemma F. 1 it is easy to verify that the matrix $\Theta$ satisfies the equation $\Theta^{\prime} \Theta=$ $\Theta \Theta^{\prime}=I$.

Finally, one obtain

$$
D_{11}^{(5)}=0
$$

Operation 5 : While performing Operation 4 it might be happened that $D_{22}^{(5)} \neq 0$. Should this be the case, one has to apply the procedure outlined in Operation 1 above by letting

$$
\begin{aligned}
y^{(6)} & :=y^{(5)}-D_{22}^{(5)} u^{(5)} \\
w^{(6)} & :=w^{(5)}, u^{(6)}:=u^{(5)}, \quad z^{(6)}:=z^{(5)}
\end{aligned}
$$

so that

$$
\begin{aligned}
& A^{(6)}=A^{(5)}, \quad B_{i}^{(6)}=B_{i}^{(5)}, \quad i=1,2 \\
& C_{i}^{(6)}=C_{i}^{(5)}, i=1,2 \\
& D_{11}^{(6)}=0, \quad D_{12}^{(6)}=D_{12}^{(5)}, \quad D_{21}^{(6)}=D_{21}^{(5)}
\end{aligned}
$$

and

$$
D_{22}^{(6)}=0
$$

Operation 6 : While performing Operation 4 it might be also happened that the matrices $D_{12}^{(5) \prime} D_{12}^{(5)}$ and/or $D_{21}^{(5)} D_{21}^{(5) \prime}$, though positive definite, are different from the identity matrices. Should this be the case, one has to apply the procedure outlined in Operation 2 above by letting

$$
\begin{aligned}
w^{(7)} & :=\hat{V}_{w}^{\prime} w^{(6)}, u^{(7)}:=\hat{S}_{u} u^{(6)} \\
z^{(7)} & :=\hat{V}_{z} z^{(6)}, \quad y^{(7)}:=\hat{S}_{y}^{-1} y^{(6)}
\end{aligned}
$$

where the matrices $\hat{V}_{w}, \hat{S}_{u}, \hat{V}_{z}$ and $\hat{S}_{y}$ are computed in such a way as to comply with the same requirements expressed in Operation 2, so that

$$
\begin{aligned}
& A^{(7)}=A^{(6)}, \quad B_{1}^{(7)}=B_{1}^{(6)} \hat{V}_{w}, \quad B_{2}^{(7)}=B_{2}^{(6)} \hat{S}_{u}^{-1} \\
& C_{1}^{(7)}=\hat{V}_{z} C_{1}^{(6)}, \quad C_{2}^{(7)}=\hat{S}_{y}^{-1} C_{2}^{(6)}, \quad D_{11}^{(7)}=0 \\
& D_{22}^{(7)}=0
\end{aligned}
$$

and

$$
D_{12}^{(7)}=\hat{V}_{z} D_{12}^{(6)} \hat{S}_{u}^{-1}=\left[\begin{array}{l}
0 \\
I
\end{array}\right], \quad D_{21}^{(7)}=\hat{S}_{y}^{-1} D_{21}^{(6)} \hat{V}_{w}=\left[\begin{array}{ll}
0 & I
\end{array}\right]
$$

Finally, given a controller $K^{(7)}(s)$ which is admissible in $R H_{\infty}$ for system $P^{(7)}(s)$ and such that $\left\|T\left(z^{(7)}, w^{(7)} ; s\right)\right\|_{\infty}<\gamma$, it is fairly apparent how to design a controller $K^{(1)}(s)$ which is admissible in $R H_{\infty}$ for system $P^{(1)}(s)$ and such that $\left\|T\left(z^{(1)}, w^{(1)} ; s\right)\right\|_{\infty}<\gamma$ (see also fig. 5.42).

Example 5.8 Consider system $P^{(1)}(s)$ with

$$
\begin{array}{ll}
A^{(1)}=\left[\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right], \quad B_{1}^{(1)}=\left[\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right], \quad B_{2}^{(1)}=\left[\begin{array}{l}
0 \\
1
\end{array}\right] \\
C_{1}^{(1)}=\left[\begin{array}{ll}
1 & 0 \\
1 & 1
\end{array}\right], \quad C_{2}^{(1)}=\left[\begin{array}{ll}
1 & 0
\end{array}\right], \quad D_{22}^{(1)}=1 \\
D_{11}^{(1)}=\left[\begin{array}{ll}
1 & 1 \\
1 & 0
\end{array}\right], \quad D_{12}^{(1)}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad D_{21}^{(1)}=\left[\begin{array}{ll}
1 & 1
\end{array}\right]
\end{array}
$$

By performing the operations described in Remark 5.24 one obtains

$$
\begin{aligned}
V_{z} & =I, \quad V_{w}=\frac{1}{\sqrt{2}}\left[\begin{array}{rr}
1 & 1 \\
-1 & 1
\end{array}\right], S_{u}=1, S_{y}=\sqrt{2} \\
\Lambda_{\infty} & =-\frac{1}{\sqrt{2}}, \gamma>\sqrt{2}, \quad \Theta_{11}=-\frac{1}{\gamma \sqrt{2}}\left[\begin{array}{ll}
0 & 2 \\
1 & 0
\end{array}\right] \\
\Theta_{12} & =\frac{1}{\gamma \sqrt{2}}\left[\begin{array}{cc}
\sqrt{2 \gamma^{2}-4} & 0 \\
0 & \sqrt{2 \gamma^{2}-1}
\end{array}\right] \\
\Theta_{21} & =\frac{1}{\gamma \sqrt{2}}\left[\begin{array}{cc}
\sqrt{2 \gamma^{2}-1} & 0 \\
0 & \sqrt{2 \gamma^{2}-4}
\end{array}\right] \\
\Theta_{22} & =\frac{1}{\gamma \sqrt{2}}\left[\begin{array}{ll}
0 & 1 \\
2 & 0
\end{array}\right] \\
D_{22}^{(6)} & =0, \hat{V}_{z}=\hat{V}_{w}=I \\
\hat{S}_{u} & =\frac{\gamma \sqrt{2}}{\sqrt{2 \gamma^{2}-1}}, \quad \hat{S}_{y}=\frac{\gamma}{\sqrt{\gamma^{2}-2}}
\end{aligned}
$$

Remark 5.25 It is worth noticing that the well known separation property of the control system eigenvalues (which has been discussed in Remark 4.17 within the framework of the $R H_{2}$ partial information problem) does not hold anymore in the $R H_{\infty}$ setting. However, it is possible to read Theorem 5.4 in the light of a weak separation principle. In fact, the controller resulting from the choice $Q(s)=0$ in fig. 5.39 , may be written as

$$
\begin{aligned}
\dot{\xi} & =A \xi+B_{2} u+B_{1} \hat{w}_{w}+Z_{\infty} L_{\infty}\left(C_{2} \xi+D_{21} \hat{w}_{w}-y\right) \\
u & =F_{\infty} \xi
\end{aligned}
$$

where

$$
\hat{w}_{w}:=\gamma^{-2} B_{1}^{\prime} P_{\infty} \xi
$$

In view of eqs. (5.131) - (5.133), (5.138), (5.141), it is easy to verify that $Z_{\infty} L_{\infty}=$ $-\Pi_{t \infty} C_{2 t}^{\prime}-B_{1} D_{21}^{\prime}$, so that it coincides with the gain of the filter for the linear combination $F_{\infty} x$ (recall the solution of the output estimation problem for system $P_{t}(s)$ defined by eqs. (5.136)-(5.138)). Further, $\hat{w}_{w}$ can be seen as the worst input for the control problem
with full information (recall Remark 5.10 concerning differential games and see fig. 5.15, too). Therefore, it can be concluded that the controller for the partial information problem is the filter for the control law of the full information problem when the worst disturbance $w_{p}:=\gamma^{-2} B_{1}^{\prime} P_{\infty} x$ acts on the system. This claim is correct in the $R H_{2}$ setting too (recall Remark 4.19), if $w_{w}$ is set equal to zero (hence $\hat{w}_{w}=0$ as well).

Remark 5.26 (Parametrization of the set $\mathcal{F}_{\infty \gamma}$ ) Observe that $\mathcal{F}_{\infty \gamma}=\mathcal{F}_{\infty \gamma r}$. In fact, the set of controllers which are admissible in $R H_{\infty}$ for $P(s)$ has been shown to coincide with the set of controllers which are admissible in $R H_{\infty}$ for $P_{t}(s)$ (see the proof of Theorem 5.4). Therefore, the claim is true thanks to Remark 5.15 which can be exploited because the structure of $P_{t}(s)$ is identical to that of $P_{O E}(s)$.

Remark 5.27 Theorem 5.3 of Section 5.4 (which is relative to the output estimation problem in the $R H_{\infty}$ setting) can be derived as a particular case of Theorem 5.4. In fact, it suffices to set $D_{12}=I$ in Assumptions 5.6-5.9 in order to conclude that Assumptions $5.3-5.5$ are satisfied. In particular, it results $C_{1 c}=0$ so that the (unique) symmetric, positive semidefinite and stabilizing solution of the Riccati equation (5.124) is $P_{\infty}=0$, which implies $Z_{\infty}=I$. In such a context, the conclusions of Theorem 5.4 are immediately redrawn to those of Theorem 5.3. In a similar way, the solution of the disturbance feedforward problem dealt with in Remark 5.12 can be viewed as the solution of a particular case of the partial information problem, namely the case in which $D_{21}=I$. Indeed, letting $D_{21}=I$ in Assumptions 5.6-5.9, it is easy to verify that the assumptions made in Remark 5.12 are satisfied. In particular, it results $B_{1 f}=0$ and hence $\Pi_{\infty}=0, Z_{\infty}=I$, so that the conclusions of Theorem 5.4 coincide with those illustrated in Remark 5.12.

Remark 5.28 (Partial information and parametric perturbations) Consider Problem 5.1 relative to the system $P_{\Omega}(s)$ described by

$$
\begin{align*}
\dot{x} & =\left(A+\Delta_{A}\right) x+\bar{B}_{1} w_{1}+B_{2} u  \tag{5.148}\\
z_{1} & =\bar{C}_{1} x+\bar{D}_{12} u  \tag{5.149}\\
y & =C_{2} x+\bar{D}_{21} w_{1} \tag{5.150}
\end{align*}
$$

The structure of the perturbation $\Delta_{A}$ is taken to be

$$
\begin{equation*}
\Delta_{A}:=L \Omega M, \quad\|\Omega\|<\beta^{-1} \tag{5.151}
\end{equation*}
$$

The robust (with respect to the parametric perturbations) control problem in $R H_{\infty}$ consists in designing a controller which, for all $\Omega,\|\Omega\|<\beta^{-1}$, is such that: (i) it is admissible in $R H_{\infty}$ relative to $P_{\Omega}(s)$; (ii) the $R H_{\infty}$ norm of the transfer function from the input $w_{1}$ to the output $z_{1}$ (from now on denoted by $T_{\Omega}\left(z_{1}, w_{1} ; s\right)$ ) is less than a certain positive scalar $\alpha$. This problem can be tackled by considering the system $P(s)$

$$
\begin{aligned}
\dot{\xi} & =A \xi+B_{1} w+B_{2} u \\
z & =C_{1} \xi+D_{12} u \\
y & =C_{2} \xi+D_{21} w
\end{aligned}
$$

where

$$
\begin{aligned}
w & :=\left[\begin{array}{l}
w_{1} \\
w_{2}
\end{array}\right], z:=\left[\begin{array}{l}
z_{1} \\
z_{2}
\end{array}\right] \\
B_{1} & :=\left[\begin{array}{ll}
\bar{B}_{1} & L
\end{array}\right], \quad D_{21}:=\left[\begin{array}{ll}
\bar{D}_{21} & 0
\end{array}\right] \\
D_{12} & :=\left[\begin{array}{c}
\bar{D}_{12} \\
0
\end{array}\right], \quad C_{1}:=\left[\begin{array}{l}
\bar{C}_{1} \\
M
\end{array}\right]
\end{aligned}
$$

Observe that $P(s)=P_{\Omega}(s)$ if $w_{2}=\Omega z_{2}$. If Assumptions 5.6-5.9 are verified for system $P_{0}(s)$, that is for system $P_{\Omega}(s)$ in nominal conditions ( $\Omega=0$ ), then they are verified also for


Figure 5.43: The performances of the controllers $R$ and $H$
system $P(s)$. Therefore Theorem 5.4 can be applied to such a system with $\gamma:=\min [\alpha, \beta]$. The solution of the relevant problem is then a solution also for the original problem. In fact, the controller which solves Problem 5.1 for system $P(s)$ is such that $\|T(z, w ; s)\|_{\infty}<\gamma$ which implies $J:=\|z\|_{2}^{2}-\gamma^{2}\|w\|_{2}^{2}<0, \forall w \in R H_{2}$. This controller, when applied to system $P_{\Omega}(s)$, is such that $\left\|T_{\Omega}\left(z_{1}, w_{1} ; s\right)\right\|_{\infty}<\alpha$, which implies $J_{1}:=\left\|z_{1}\right\|_{2}^{2}-\alpha^{2}\left\|w_{1}\right\|_{2}^{2}<0, \forall w_{1} \in R H_{2}$. Indeed, by recalling that $w_{2}=\Omega z_{2}$, it follows

$$
\begin{aligned}
J_{1} & =\left\|z_{1}\right\|_{2}^{2}-\alpha^{2}\left\|w_{1}\right\|_{2}^{2} \\
& \leq\left\|z_{1}\right\|_{2}^{2}-\gamma^{2}\left\|w_{1}\right\|_{2}^{2} \\
& \leq\|z\|_{2}^{2}-\left\|z_{2}\right\|_{2}^{2}-\gamma^{2}\|w\|_{2}^{2}+\gamma^{2}\left\|w_{2}\right\|_{2}^{2} \\
& \leq J-\left\|z_{2}\right\|_{2}^{2}+\gamma^{2}\left\|w_{2}\right\|_{2}^{2} \\
& \leq J-\left\|z_{2}\right\|_{2}^{2}+\gamma^{2}\left\|\Omega z_{2}\right\|_{2}^{2} \\
& \leq J-\left\|z_{2}\right\|_{2}^{2}\left(1-\gamma^{2}\|\Omega\|^{2}\right) \\
& \leq J-\left\|z_{2}\right\|_{2}^{2}\left(1-\gamma^{2} \beta^{-2}\right) \\
& \leq J<0
\end{aligned}
$$

It is left to be proved that the controller which solves Problem 5.1 for $P(s)$ stabilizes $P_{\Omega}(s), \forall \Omega,\|\Omega\|<\beta^{-1}$ as well. This result is a direct consequence of Theorem 5.1.
Example 5.9 Consider the system (5.148)-(5.151) with

$$
\begin{aligned}
A & =\left[\begin{array}{rr}
0 & 1 \\
-1 & -1
\end{array}\right], \quad \bar{B}_{1}=\left[\begin{array}{c}
\varphi^{-1} \\
0
\end{array}\right], \quad B_{2}=\left[\begin{array}{l}
0 \\
\varphi
\end{array}\right], \quad L=\left[\begin{array}{l}
0 \\
1
\end{array}\right] \\
M & =\left[\begin{array}{ll}
0 & 1
\end{array}\right], \quad C_{1}=\left[\begin{array}{l}
\varphi^{-1} \\
\varphi^{-1}
\end{array}\right], \quad C_{2}=\left[\begin{array}{ll}
\varphi & 0
\end{array}\right] \\
\bar{D}_{12} & =\bar{D}_{21}=1, \quad \beta=1, \quad \varphi:=\sqrt{2}
\end{aligned}
$$

A controller is sought which, for each $\Omega,\|\Omega\|<\beta^{-1}$, is admissible in $R H_{\infty}$ for system $P_{\Omega}(s)$ and is such that $\left\|T\left(z_{1}, w_{1} ; s\right)\right\|_{\infty}<\alpha$. Letting $\alpha=\beta$, a controller, labeled with $R$, is designed according to what has been shown in Remark 5.28. On the contrary, the label $H$ denotes the controller resulting from Theorem 5.4 when applied to system $P_{0}(s)$ and for the choice $Q(s)=0$. The performances of the two controllers are compared in fig. 5.43 where $\left\|T\left(z_{1}, w_{1} ; s\right)\right\|_{\infty}$ is plotted against the system damping factor $\xi:=(1-\Omega) / 2$. The controller $R$ is weakly sensitive and apparently behaves in a much better way for low values of $\xi$.


Figure 5.44: The performances of the controllers $R$ and $H(M=M S)$

Example 5.10 Consider the system defined by the matrices

$$
\begin{aligned}
& A=\left[\begin{array}{rr}
0 & 1 \\
-1 & -1
\end{array}\right]+L \Omega M, \quad B_{1}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad B_{2}=\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right] \\
& L=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad C_{1}=\left[\begin{array}{ll}
1 & 2
\end{array}\right], \quad C_{2}=D_{12}=\left[\begin{array}{ll}
1 & 0
\end{array}\right], \quad D_{21}=1
\end{aligned}
$$

while $M=M S:=[20]$ or $M=M U:=[02]$. The parameter $\Omega$ has nominal value equal to 0 and describes the uncertainty which is supposed to affect the system under consideration. A stabilizing controller is sought which makes small the effect of the disturbance $w$ on the output $z=C_{1} x+D_{12} u$. By utilizing the first component only of the control variable $u$ it is possible to design a controller, labeled with $I C$, which makes zero the transfer function from $w$ to $z$ in nominal conditions. The transfer function of such a controller (which achieves the perfect indirect compensation of the disturbance) is

$$
K_{I C}(s)=-\frac{1+2 s}{s^{2}+3 s+3}
$$

If, on the contrary, the uncertain knowledge of the parameter $\Omega$ has to be somehow taken into account, a controller, labeled with $R$, can be designed according to the discussion in Remark 5.28 . Consistently, a new variable $z_{a}:=M x+u_{2}$ is added to the system. When $M=M S$ one finds, independently of the value of $\gamma, P_{\infty}=\Pi_{\infty}=0$ because the matrices $A_{c}$ and $A_{f}$ are stable and $C_{1 c}=0, B_{1 f}=0$. Correspondingly,

$$
F_{\infty}=\left[\begin{array}{rr}
-1 & -2 \\
-2 & 0
\end{array}\right], \quad L_{\infty}=\left[\begin{array}{r}
0 \\
-1
\end{array}\right]
$$

This is no more the case when $M=M U$, so that, for $\gamma=0.275$ (notice that for $\gamma=0.270$ a solution of eq. (5.25) with the required properties does not exist) one obtains

$$
F_{\infty}=\left[\begin{array}{rr}
7.67 & -4.43 \\
-30.87 & 6.67
\end{array}\right], \quad L_{\infty}=\left[\begin{array}{r}
0 \\
-1
\end{array}\right]
$$

The performances of the two controllers are compared in fig. 5.44 and 5.45 for the considered cases. The choice of the set of values of the parameter $\Omega$ reflects the need of stability for the control system resulting from the insertion of the controller $I C$. In both cases, the better behavior of the controller $R$ is apparent, especially for small values of $\Omega$.


Figure 5.45: The performances of the controllers $R$ and $H \quad(M=M U)$

Remark 5.29 It is worth noticing that, also in the partial information case, when $\gamma \rightarrow \infty$ the $R H_{\infty}$ controller of fig. 5.39 (with $Q(s)=0$ ) tends to the $R H_{2}$ controller of fig. 4.19 (with $Q(s)=0$ ).

Remark 5.30 The contents of Remark 5.13 (as for matrix $D_{12}$ ) and of Remark 5.21 (as for matrix $D_{21}$ ) apply with no changes to the problem dealt with in the present section.

Remark 5.31 In view of Remarks 5.14 and 5.22 it can be said that, under Assumptions 5.8 and 5.9, Assumption 5.6 amounts to requiring that the two subsystems of $P(s)$ having transfer functions $P_{12}(s)$ (that is, system $\Sigma\left(A, B_{2}, C_{1}, D_{12}\right)$ with input $u$ and output $\left.z\right)$ and $P_{21}(s)$ (that is, system $\Sigma\left(A, B_{1}, C_{2}, D_{21}\right)$ with input $w$ and output $y$ ), respectively, do not have zeros in the closed right half plane.

### 5.6 The operatorial approach

This section is aimed at tackling the partial information problem in the $R H_{\infty}$ setting from a point of view which is rather different from the one adopted in Section 5.5. In fact, reference will be made to the theory of linear operators applied to dynamical systems. The minimum norm problem will be presented by constraining the attention to the scalar case only, as the general multivariable situation is substantially more complex to be handled.

Recall that the problem at hand consists in designing a controller $K(s)$ which, with reference to fig. 5.46 , is admissible in $R H_{\infty}$ for $P(s)$ and such that the $R H_{\infty}$ norm of the transfer function $T(z, w ; s)$ from $w$ to $z$ is less than a given positive scalar $\gamma$. Indeed, it will be shown that, having restrained the attention to the case in which the variables $w, u, z, y$ are scalar, it is fairly easy to design that controller $K^{o}(s)$ which minimizes $\|T(z, w ; s)\|_{\infty}$. Preliminarily, partition the matrix $P(s)$ into the four (scalar) transfer functions $P_{i j}(s)$, so that

$$
\left[\begin{array}{l}
z_{L} \\
y_{L}
\end{array}\right]=\left[\begin{array}{ll}
P_{11}(s) & P_{12}(s) \\
P_{21}(s) & P_{22}(s)
\end{array}\right]\left[\begin{array}{l}
w_{L} \\
u_{L}
\end{array}\right]
$$



Figure 5.46: The partial information problem
and assume that system $P(s)$ can be made internally stable by implementing a feedback loop from $y$ to $u$. It is obvious that a $R H_{\infty}$ admissible controller for $P(s)$ exists if and only if such an assumption is verified. This is equivalent to the internal stabilizability (from $y$ ) of the subsystem $P_{22}(s)$. Further, if $K(s)$ is a $R H_{\infty}$ admissible controller for $P_{22}(s)$, then it is $R H_{\infty}$ admissible also for $P(s)$ and the set of controllers endowed with this property coincides with the set of controllers which stabilize $P_{22}(s)$.

Second, observe that, if $K(s)$ is any controller corresponding to which the control system in fig. 5.46 is well defined, that is if

$$
\lim _{s \rightarrow \infty} \operatorname{det}\left[I-D_{22} K(s)\right] \neq 0
$$

where $D_{22}:=\lim _{s \rightarrow \infty} P_{22}(s)$, then

$$
\begin{equation*}
T(z, w ; s)=P_{11}(s)+P_{12}(s)\left[I-K(s) P_{22}(s)\right]^{-1} K(s) P_{21}(s) \tag{5.152}
\end{equation*}
$$

Hence, the problem of designing a $R H_{\infty}$ admissible controller for $P(s)$ such that $\|T(z, w ; s)\|_{\infty}<\gamma$ is equivalent to selecting a controller which verifies this last inequality among those which internally stabilize $P_{22}(s)$.

Theorem 3.7 supplies the parametrization of the set of all controllers which stabilize $P_{22}(s)$ on the basis of a double coprime factorization in $R H_{\infty}$ of $P_{22}(s)$ (also recall Theorem 3.4). More precisely, if

$$
\begin{gather*}
P_{22}(s)=N(s) M^{-1}(s)=\hat{M}^{-1}(s) \hat{N}(s)  \tag{5.153}\\
{\left[\begin{array}{cc}
\hat{X}(s) & -\hat{Y}(s) \\
-\hat{N}(s) & \hat{M}(s)
\end{array}\right]\left[\begin{array}{cc}
M(s) & Y(s) \\
N(s) & X(s)
\end{array}\right]=I} \tag{5.154}
\end{gather*}
$$

where all the functions $N(s), M(s), \hat{M}(s), \hat{N}(s), \hat{X}(s), \hat{Y}(s), X(s), Y(s)$ belong to $R H_{\infty}$, then any controller $K(s)$ which stabilizes $P_{22}(s)$ can be given the form

$$
\begin{equation*}
K(s)=[\hat{X}(s)-Q(s) \hat{N}(s)]^{-1}[\hat{Y}(s)-Q(s) \hat{M}(s)] \tag{5.155}
\end{equation*}
$$

where $Q(s):=\Sigma\left(A_{q}, B_{q}, C_{q}, D_{q}\right)$ is any element of $R H_{\infty}$ such that $\operatorname{det}\left[I-D_{22} D_{q}\right] \neq 0$.
Thanks to the explicit expression of the stabilizing controllers given by eq. (5.155), it is easy to prove that the transfer function $T(z, w ; s)$ is a linear function of the parameter $Q(s)$, as stated in the following theorem where reference is made to the functions

$$
\begin{align*}
& T_{1}(s):=P_{11}(s)+P_{12}(s) M(s) \hat{Y}(s) P_{21}(s)  \tag{5.156}\\
& T_{2}(s):=P_{12}(s) M(s)  \tag{5.157}\\
& T_{3}(s):=\hat{M}(s) P_{21}(s) \tag{5.158}
\end{align*}
$$

Theorem 5.5 The functions $T_{i}(s), i=1,2,3$ defined in eqs. (5.156)-(5.158) belong to $R H_{\infty}$. Moreover, if $K(s)$ is given by eq. (5.155), then

$$
T(z, w ; s)=T_{1}(s)-T_{2}(s) Q(s) T_{3}(s)
$$

Proof Consider a realization of $P(s)$

$$
P(s):=\left[\begin{array}{c|cc}
A & B_{1} & B_{2} \\
\hline C_{1} & D_{11} & D_{12} \\
C_{2} & D_{21} & D_{22}
\end{array}\right]
$$

with the pair $\left(A, B_{2}\right)$ stabilizable and the pair $\left(A, C_{2}\right)$ detectable. Further, let $F$ and $H$ be two matrices such that $A+B_{2} F$ and $A+H C_{2}$ are stable. Consistently with the discussion in the proof of Theorem 3.4, choose

$$
\begin{aligned}
M(s) & :=\Sigma\left(A+B_{2} F, B_{2}, F, I\right) \\
\hat{M}(s) & :=\Sigma\left(A+H C_{2}, H, C_{2}, I\right) \\
\hat{Y}(s) & :=\Sigma\left(A+H C_{2}, H,-F, 0\right)
\end{aligned}
$$

Denote with $\varphi_{i}$ (resp. $\eta_{i}$ ), $i=1, \cdots, 5$ the state (resp. output) variables of the five systems $P_{21}(s), \hat{Y}(s), M(s), P_{12}(s), P_{11}(s)$, taken in their order of appearance. If $\vartheta_{1}$ and $\nu_{1}$ denote the output and input variables of system $T_{1}(s)$, respectively, then such a system is described by (recall eq. (5.156))

$$
\begin{aligned}
& P_{21}(s)=\left\{\begin{array}{l}
\dot{\varphi}_{1}=A \varphi_{1}+B_{1} \nu_{1} \\
\eta_{1}=C_{2} \varphi_{1}+D_{21} \nu_{1}
\end{array}\right. \\
& \hat{Y}(s)=\left\{\begin{array}{l}
\dot{\varphi}_{2}=\left(A+H C_{2}\right) \varphi_{2}+H \eta_{1} \\
\eta_{2}=-F \varphi_{2}
\end{array}\right. \\
& M(s)=\left\{\begin{array}{l}
\dot{\varphi}_{3}=\left(A+B_{2} F\right) \varphi_{3}+B_{2} \eta_{2} \\
\eta_{3}=F \varphi_{3}+\eta_{2}
\end{array}\right. \\
& P_{12}(s)=\left\{\begin{array}{l}
\dot{\varphi}_{4}=A \varphi_{4}+B_{2} \eta_{3} \\
\eta_{4}=C_{1} \varphi_{4}+D_{12} \eta_{3}
\end{array}\right. \\
& P_{11}(s)=\left\{\begin{array}{l}
\dot{\varphi}_{5}=A \varphi_{5}+B_{1} \nu_{1} \\
\eta_{5}=C_{1} \varphi_{5}+D_{11} \nu_{1}
\end{array}\right. \\
& \vartheta_{1}=\eta_{4}+\eta_{5}
\end{aligned}
$$

Letting $\varepsilon_{1}:=\varphi_{1}-\varphi_{5}$ and $\varepsilon_{2}:=\varphi_{4}-\varphi_{3}$, it is easy to verify that

$$
\begin{aligned}
& \dot{\varepsilon}_{1}=A \varepsilon_{1} \\
& \dot{\varepsilon}_{2}=A \varepsilon_{2}
\end{aligned}
$$

and that a realization of the transfer function $T_{1}(s)$ can be built up by exploiting the state variables $\varphi_{2}, \varphi_{3}, \varphi_{5}$ only. Further, by defining $\varepsilon_{3}:=\varphi_{5}+\varphi_{2}$ and $\varepsilon_{4}:=\varphi_{5}+\varphi_{3}$, one obtains

$$
\begin{align*}
& \dot{\varepsilon}_{3}=\left(A+H C_{2}\right) \varepsilon_{3}+\left(B_{1}+H D_{21}\right) \nu_{1}  \tag{5.159}\\
& \dot{\varepsilon}_{4}=-B_{2} F \varepsilon_{3}+\left(A+B_{2} F\right) \varepsilon_{4}+B_{1} \nu_{1}  \tag{5.160}\\
& \vartheta_{1}=-D_{12} F \varepsilon_{3}+\left(C_{1}+D_{12} F\right) \varepsilon_{4}+D_{11} \nu_{1} \tag{5.161}
\end{align*}
$$

These equations define a simpler realization of $T_{1}(s)$ which allows one to claim that $T_{1}(s) \in R H_{\infty}$, since the eigenvalues of the dynamic matrix of the system described by eqs. (5.159)-(5.161) are those of matrix $A+B_{2} F$ together with those of matrix $A+H C_{2}$, both matrices being stable by construction.

Now, let $\varphi_{i}$ (resp. $\eta_{i}$ ), $i=6,7$ be the state (resp. output) variables of the two systems $M(s)$ and $P_{12}(s)$, taken in their order of appearance. If $\vartheta_{2}$ and $\nu_{2}$ denote the output and input variables of system $T_{2}(s)$, respectively, then such a system is described by (recall eq. (5.157))

$$
\begin{gathered}
M(s)=\left\{\begin{array}{l}
\dot{\varphi}_{6}=\left(A+B_{2} F\right) \varphi_{6}+B_{2} \nu_{2} \\
\eta_{6}=F \varphi_{6}+\nu_{2}
\end{array}\right. \\
P_{12}(s)=\left\{\begin{array}{l}
\dot{\varphi}_{7}=A \varphi_{7}+B_{2} \eta_{6} \\
\eta_{7}=C_{1} \varphi_{7}+D_{12} \eta_{6}
\end{array}\right. \\
\vartheta_{2}=\eta_{7}
\end{gathered}
$$

Letting $\varepsilon_{5}:=\varphi_{6}-\varphi_{7}$, it is easy to verify that

$$
\dot{\varepsilon}_{5}=A \varepsilon_{5}
$$

and that a realization of the transfer function $T_{2}(s)$ can be built up by exploiting the state variables $\varphi_{6}$ only. The dynamic matrix of such a realization is $A+B_{2} F$, so that $T_{2}(s) \in R H_{\infty}$ too.

Finally, let $\varphi_{i}$ (resp. $\eta_{i}$ ) , $i=8,9$ be the state (resp. output) variables of the two systems $P_{21}(s)$ and $\hat{M}(s)$, taken in their order of appearance. If $\vartheta_{3}$ and $\nu_{3}$ denote the output and input variables of system $T_{3}(s)$, respectively, then such a system is described by (recall eq. (5.158))

$$
\begin{gathered}
P_{21}(s)=\left\{\begin{array}{l}
\dot{\varphi}_{8}=A \varphi_{8}+B_{1} \nu_{3} \\
\eta_{8}=C_{2} \varphi_{8}+D_{21} \nu_{3}
\end{array}\right. \\
\hat{M}(s)=\left\{\begin{array}{l}
\dot{\varphi}_{9}=\left(A+H C_{2}\right) \varphi_{9}+H \eta_{8} \\
\eta_{9}=C_{2} \varphi_{9}+\eta_{8}
\end{array}\right. \\
\vartheta_{3}=\eta_{9}
\end{gathered}
$$

Letting $\varepsilon_{6}:=\varphi_{8}+\varphi_{9}$, it is easy to verify that

$$
\begin{aligned}
& \dot{\varepsilon}_{6}=\left(A+H C_{2}\right) \varepsilon_{6}+\left(B_{1}+H D_{21}\right) \nu_{3} \\
& \vartheta_{3}=C_{2} \varepsilon_{6}+D_{21} \nu_{3}
\end{aligned}
$$

so that a realization of the transfer function $T_{3}(s)$ is characterized the dynamic matrix $A+H C_{2}$. Hence, $T_{3}(s) \in R H_{\infty}$.

As for the expression of $T(z, w ; s)$, first observe that, in view of eqs. (5.153)(5.155), it is

$$
\begin{aligned}
{\left[I-K(s) P_{22}(s)\right]^{-1}=} & {\left[I-[\hat{X}(s)-Q(s) \hat{N}(s)]^{-1}\right.} \\
& \left.\cdot[\hat{Y}(s)-Q(s) \hat{M}(s)] N(s) M^{-1}(s)\right]^{-1} \\
= & \left\{[\hat{X}(s)-Q(s) \hat{N}(s)]^{-1}[\hat{X}(s)-Q(s) \hat{N}(s)-\right.
\end{aligned}
$$

$$
\begin{aligned}
& \left.\left.-\hat{Y}(s) N(s) M^{-1}(s)+Q(s) \hat{M}(s) N(s) M^{-1}(s)\right]\right\}^{-1} \\
= & \left\{[\hat{X}(s)-Q(s) \hat{N}(s)]^{-1}[\hat{X}(s) M(s)-\hat{Y}(s) N(s)-\right. \\
& \left.-Q(s)[\hat{N}(s) M(s)-\hat{M}(s) N(s)] M^{-1}\right\}^{-1} \\
= & M(s)[\hat{X}(s)-Q(s) \hat{N}(s)]
\end{aligned}
$$

and

$$
\begin{aligned}
{\left[I-K(s) P_{22}(s)\right]^{-1}=} & {\left[I-[\hat{X}(s)-Q(s) \hat{N}(s)]^{-1}[\hat{Y}(s)-\right.} \\
& \left.-Q(s) \hat{M}(s)] N(s) M^{-1}(s)\right]^{-1} \\
= & \left\{[\hat{X}(s)-Q(s) \hat{N}(s)]^{-1}[\hat{X}(s)-\right. \\
& -Q(s) \hat{N}(s)-\hat{Y}(s) N(s) M^{-1}(s)+ \\
& \left.\left.+Q(s) \hat{M}(s) N(s) M^{-1}(s)\right]\right\}^{-1} \\
= & \left\{[\hat{X}(s)-Q(s) \hat{N}(s)]^{-1}[\hat{X}(s) M(s)-\right. \\
& -\hat{Y}(s) N(s)-Q(s)[\hat{N}(s) M(s)- \\
& \left.-\hat{M}(s) N(s)]] M^{-1}\right\}^{-1} \\
= & M(s)[\hat{X}(s)-Q(s) \hat{N}(s)]
\end{aligned}
$$

so that

$$
\left[I-K(s) P_{22}(s)\right]^{-1} K(s)=M(s)[\hat{Y}(s)-Q(s) \hat{M}(s)]
$$

Then, by recalling also eqs. (5.152),(5.156)-(5.158), it follows

$$
\begin{aligned}
T(z, w ; s)= & P_{11}(s)+P_{12}(s)\left[I-K(s) P_{22}(s)\right]^{-1} K(s) P_{21}(s) \\
= & P_{11}(s)+P_{12}(s) M(s)[\hat{Y}(s)-Q(s) \hat{M}(s)] P_{21}(s) \\
= & P_{11}(s)+P_{12}(s) M(s) \hat{Y}(s) P_{21}(s)- \\
& -P_{12}(s) M(s) Q(s) \hat{M}(s) P_{21}(s) \\
= & T_{1}(s)-T_{2}(s) Q(s) T_{3}(s)
\end{aligned}
$$

In view of this theorem the problem of designing a $R H_{\infty}$ admissible controller $K(s)$ for system $P(s)$ such that $\|T(z, w ; s)\|_{\infty}<\gamma$ is apparently equivalent to that of finding a function $Q(s) \in R H_{\infty}$ such that $\left\|T_{1}(s)-T_{2}(s) Q(s) T_{3}(s)\right\|_{\infty}<\gamma$. Indeed, it will be shown that, under a suitable assumption, a function $Q^{\circ}(s) \in R H_{\infty}$ can be found which minimizes such a norm: the so called model matching problem is solved in this way and the original control problem has got an optimal answer (in the sense of making as small as possible $\left.\|T(z, w ; s)\|_{\infty}\right)$ at the same time.

The following simple examples show that solving the model matching problem is not trivial in general, as there are cases where the solution exists and other ones where it does not.

Example 5.11 Let

$$
T_{3}(s)=1, \quad T_{2}(s)=\frac{s-1}{s+1}
$$

while $T_{1}(s)$ is any element of $R H_{\infty}$. Then it follows, for any $Q(s) \in R H_{\infty}$ (recall Definition 2.24),

$$
\left\|T_{1}(s)-T_{2}(s) Q(s)\right\|_{\infty} \geq\left|T_{1}(1)-T_{2}(1) Q(1)\right|=\left|T_{1}(1)\right|
$$

so that

$$
\inf _{Q(s) \in R H_{\infty}}\left|T_{1}(s)-T_{2}(s) Q(s) \|:=\alpha \geq\left|T_{1}(1)\right|\right.
$$

Chosen

$$
Q^{o}(s):=\frac{T_{1}(s)-T_{1}(1)}{T_{2}(s)}
$$

it is easy to verify that $Q^{\circ}(s) \in R H_{\infty}$ (observe that the zero of $T_{2}(s)$ is canceled out because the numerator of $Q^{\circ}(s)$ vanishes for $\left.s=1\right)$. On the other side,

$$
T_{1}(s)-T_{2}(s) Q^{o}(s)=T_{1}(1)
$$

and therefore $\alpha \leq\left\|T_{1}(s)-T_{2}(s) Q^{\circ}(s)\right\|_{\infty}=\left|T_{1}(1)\right| \leq \alpha$, so that

$$
\left\|T_{1}(s)-T_{2}(s) Q^{o}(s)\right\|_{\infty}=\alpha
$$

Example 5.12 Let $T_{3}(s)=1$ and

$$
T_{1}(s)=\frac{1}{s+1}, \quad T_{2}(s)=\frac{1}{(s+1)^{2}}
$$

Define

$$
Q_{\varepsilon}(s):=\frac{s+1}{\varepsilon s+1}, \quad 0<\varepsilon<1
$$

so that

$$
T_{1}(s)-T_{2}(s) Q_{\varepsilon}(s)=\frac{\varepsilon s}{(s+1)(\varepsilon s+1)}
$$

The diagram of $\left|T_{1}(j \omega)-T_{2}(j \omega) Q_{\varepsilon}(j \omega)\right|$ lies always below $\varepsilon$. Therefore, it is $\| T_{1}(s)-$ $T_{2}(s) Q_{\varepsilon}(s) \|_{\infty} \leq \varepsilon$, so that

$$
\inf _{\varepsilon}\left\|T_{1}(s)-T_{2}(s) Q_{\varepsilon}(s)\right\|_{\infty}=0
$$

This value of the norm is attained in correspondence of a $Q_{\varepsilon}^{o}(s)$ which is such that $T_{1}(s)-$ $T_{2}(s) Q_{\varepsilon}^{o}(s)=0$, namely $Q_{\varepsilon}^{o}(s)=s+1 \notin R H_{\infty}$.

The result to be presented in the next theorem makes reference to two functions $f(s)$ and $g(s)$ which are derived from a generic scalar function $F(s) \in R L_{\infty}$ in the following way. Let $F_{a}(s), F_{s}(s), F_{\infty}$ be such that

$$
\begin{equation*}
F(s)=F_{a}(s)+F_{s}(s)+F_{\infty} \tag{5.162}
\end{equation*}
$$

where

$$
\begin{equation*}
F_{a}(s) \in R H_{2}^{\perp}, \quad F_{s}(s) \in R H_{2}, \quad F_{\infty}:=\lim _{s \rightarrow \infty} F(s) \tag{5.163}
\end{equation*}
$$

Further, let

$$
F_{a}(s):=\left[\begin{array}{c|c}
A & B  \tag{5.164}\\
\hline C & 0
\end{array}\right]
$$

and suppose that the triple $(A, B, C)$ constitutes a minimal realization of $F_{a}(s)$. Then, denote by $P_{r}>0$ and $P_{o}>0$ the unique solutions (recall Lemma C.1) of the two Lyapunov equations (in the unknown $P$ )

$$
\begin{align*}
& 0=-P A^{\prime}-A P+B B^{\prime}  \tag{5.165}\\
& 0=-P A-A^{\prime} P+C^{\prime} C \tag{5.166}
\end{align*}
$$

respectively. Moreover, let $\lambda_{M}^{2}$ be the maximum eigenvalue of the matrix $P_{o} P_{r}$ (recall Lemma B. 10 and the fact that all the eigenvalues of this matrix are different from 0 , since it is nonsingular) and $\beta$ a corresponding eigenvector, so that

$$
\begin{equation*}
P_{o} P_{r} \beta=\lambda_{M}^{2} \beta, \quad \beta \neq 0 \tag{5.167}
\end{equation*}
$$

and define

$$
\begin{equation*}
\chi:=\lambda_{M}^{-1} P_{r} \beta \tag{5.168}
\end{equation*}
$$

Now the following functions

$$
\begin{align*}
& f(s):=\left[\begin{array}{c|c}
-A^{\prime} & \beta \\
\hline B^{\prime} & 0
\end{array}\right] \in R H_{2}  \tag{5.169}\\
& g(s):=\left[\begin{array}{c|c}
A & \chi \\
\hline C & 0
\end{array}\right] \in R H_{2}^{\perp} \tag{5.170}
\end{align*}
$$

can be associated with the scalar function $F(s) \in R L_{\infty}$. The two functions defined by eqs. (5.169),(5.170) are endowed with the property stated in the following lemma, where $\Gamma_{F \sim}$ and $\Gamma_{F}^{*} \sim$ are the Hankel operator with symbol $F^{\sim}$ and its adjoint, respectively (recall Definition 2.34 and Lemma 2.27).

Lemma 5.2 Let $F(s) \in R L_{\infty}$ be a scalar function and $f(s) \in R H_{2}$ and $g(s) \in R H_{2}^{\perp}$ be two functions derived from $F(s)$ according to eqs. (5.162)-(5.170). Then

$$
\begin{aligned}
& \Gamma_{F \sim} g(s)=\lambda_{M} f(s) \\
& \Gamma_{F \sim}^{*} f(s)=\lambda_{M} g(s)
\end{aligned}
$$

Proof From eq. (5.165) it follows

$$
\begin{aligned}
& C(s I-A)^{-1}\left(P_{r} A^{\prime}+P_{r} s-s P_{r}+A P_{r}\right)\left(s I+A^{\prime}\right)^{-1} \beta= \\
& =C(s I-A)^{-1} B B^{\prime}\left(s I+A^{\prime}\right)^{-1} \beta
\end{aligned}
$$

Therefore, by taking into account eqs. (5.164)-(5.169), one obtains

$$
\begin{aligned}
& C(s I-A)^{-1}\left[P_{r}\left(s I+A^{\prime}\right)-(s I-A) P_{r}\right]\left(s I+A^{\prime}\right)^{-1} \beta= \\
& =F_{a}(s) f(s)
\end{aligned}
$$

from which it follows

$$
\begin{equation*}
C(s I-A)^{-1} P_{r} \beta-C P_{r}\left(s I+A^{\prime}\right)^{-1} \beta=F_{a}(s) f(s) \tag{5.171}
\end{equation*}
$$

On the other hand, eq. (5.168) is equivalent to

$$
P_{r} \beta=\lambda_{M} \chi
$$

so that, by recalling eqs. (5.170),(5.171), it results

$$
\lambda_{M} g(s)-C P_{r}\left(s I+A^{\prime}\right)^{-1} \beta=F_{a}(s) f(s)
$$

If the antistable orthogonal projection operator $\Pi_{a}$ is applied to both sides of this equation (recall Definitions 2.32-2.34, Remark 2.25 and Lemmas 2.26, 2.27) one obtains

$$
\lambda_{M} g(s)=\Gamma_{F_{a}^{\sim}}^{*} f(s)=\Gamma_{F^{\sim}}^{*} f(s)
$$

In a similar way, from eq. (5.166) it follows

$$
\begin{aligned}
& B^{\prime}\left(s I+A^{\prime}\right)^{-1}\left(A^{\prime} P_{o}+s P_{o}-P_{o} s+P_{o} A\right)(s I-A)^{-1} \chi= \\
& =B^{\prime}\left(s I+A^{\prime}\right)^{-1} C^{\prime} C(s I-A)^{-1} \chi
\end{aligned}
$$

so that, by taking into account eqs. (5.164), (5.167), (5.168), (5.170), one has

$$
B^{\prime} P_{o}(s I-A)^{-1} \chi-B^{\prime}\left(s I+A^{\prime}\right)^{-1} \beta \lambda_{M}=-F_{a}^{\sim}(s) g(s)
$$

In view of eq. (5.169), if the stable orthogonal projection operator $\Pi_{s}$ is applied to both sides of this equation (recall Definitions 2.32-2.34 and Remark 2.25), one obtains

$$
\lambda_{M} f(s)=\Gamma_{F_{a}} f(s)=\Gamma_{F \sim f}(s)
$$

Theorem 5.6 Let $F(s) \in R L_{\infty}$ be a given scalar function. Then, the function

$$
\begin{equation*}
X^{o}(s):=F(s)-\frac{\left\|\Gamma_{F^{\sim}}\right\| g(s)}{f(s)} \tag{5.172}
\end{equation*}
$$

is such that

$$
\left\|F(s)-X^{o}(s)\right\|_{\infty}=\inf _{X(s) \in R H_{\infty}}\|F(s)-X(s)\|_{\infty}
$$

In eq. (5.172) the two functions $f(s)$ and $g(s)$ are specified by eqs. (5.162)-(5.170).
Proof First observe that, thanks to Nehari's theorem (Theorem 2.19), there exists a function $X^{o}(s) \in R H_{\infty}$ such that $\left\|\Gamma_{F^{\sim}}\right\|=\left\|F(s)-X^{o}(s)\right\|_{\infty}$. Then, define $h(s):=$ $\left[F(s)-X^{o}(s)\right] f(s)$ and notice that $h(s) \in R L_{2}$, because $F(s)-X^{o}(s) \in R L_{\infty}$ and $f(s) \in R H_{2}$. It follows

$$
\begin{aligned}
\left\|h(s)-\Gamma_{F^{\sim}}^{*} f(s)\right\|_{2}^{2}= & <h(s)-\Gamma_{F \sim}^{*} f(s), h(s)-\Gamma_{F^{\sim}}^{*} f(s)> \\
= & <h(s), h(s)>+<\Gamma_{F \sim}^{*} f(s), \Gamma_{F \sim}^{*} f(s)>- \\
& -2<h(s), \Gamma_{F \sim}^{*} f(s)>
\end{aligned}
$$

By taking into account that $\Gamma_{F}^{*} \sim f(s) \in R H_{2}^{\perp}, X^{o}(s) f(s) \in R H_{2}$ and Lemma 2.27, one obtains

$$
\begin{aligned}
<h(s), \Gamma_{F^{\sim}}^{*} f(s)> & =<\Pi_{a} h(s), \Gamma_{F \sim}^{*} f(s)> \\
& =<\Pi_{a}\left[F(s)-X^{o}(s)\right] f(s), \Gamma_{F^{\sim}}^{*} f(s)> \\
& =<\Pi_{a} F(s) f(s), \Gamma_{F \sim}^{*} f(s)> \\
& =<\Gamma_{F \sim}^{*} f(s), \Gamma_{F \sim}^{*} f(s)>
\end{aligned}
$$

Therefore,

$$
\begin{align*}
\left\|h(s)-\Gamma_{F \sim}^{*} f(s)\right\|_{2}^{2}= & <h(s), h(s)>-<\Gamma_{F \sim}^{*} f(s), \Gamma_{F \sim}^{*} f(s)> \\
= & \left\|\left[F(s)-X^{o}(s)\right] f(s)\right\|_{2}^{2}- \\
& -<f(s), \Gamma_{F^{\sim}} \Gamma_{F \sim}^{*} f(s)> \\
= & \left\|\left[F(s)-X^{o}(s)\right] f(s)\right\|_{2}^{2}-\lambda_{M}^{2}\|f(s)\|_{2}^{2} \tag{5.173}
\end{align*}
$$

since, by Lemma 5.2, $\Gamma_{F \sim} \sim \Gamma_{F \sim}^{*} f(s)=\lambda_{M}^{2} f(s)$. In view of Theorem 2.12 it follows

$$
\left\|\left[F(s)-X^{o}(s)\right] f(s)\right\|_{2}^{2} \leq\left\|F(s)-X^{o}(s)\right\|_{\infty}^{2}\|f(s)\|_{2}^{2}
$$

and therefore from eq. (5.173) one has

$$
\left\|h(s)-\Gamma_{F^{\sim}}^{*} f(s)\right\|_{2}^{2} \leq\left[\left\|F(s)-X^{o}(s)\right\|_{\infty}^{2}-\lambda_{M}^{2}\right]\|f(s)\|_{2}^{2}=0
$$

because $\lambda_{M}=\left\|\Gamma_{F \sim}\right\|=\left\|F(s)-X^{o}(s)\right\|_{\infty}$ (recall Remark 2.27). Thus, $h(s)=$ $\Gamma_{F^{\sim}}^{*} f(s)$ and, thanks to Lemma 5.2, $\left[F(s)-X^{o}(s)\right] f(s)=h(s)=\lambda_{M} g(s)$, which implies, in turn,

$$
X^{o}(s)=F(s)-\frac{\lambda_{M} g(s)}{f(s)}=F(s)-\frac{\left\|\Gamma_{F \sim}\right\| g(s)}{f(s)}
$$

Example 5.13 Let

$$
F(s)=\frac{c}{s-a}, \quad a>0, \quad c \neq 0
$$

Then

$$
F_{a}(s)=\left[\begin{array}{c|c}
a & 1 \\
\hline c & 0
\end{array}\right], \quad F_{s}(s)=0, \quad F_{\infty}=0
$$

and the solutions of eqs. (5.165), (5.166) are $P_{r}=(2 a)^{-1}$ and $P_{o}=c^{2}(2 a)^{-1}$, respectively. Therefore, $\lambda_{M}=|c|(2 a)^{-1}, \beta=1, \chi=|c|^{-1}$. From eqs. (5.169) and (5.170) one obtains

$$
f(s)=\frac{1}{s+a}, g(s)=\frac{c}{s-a}|c|^{-1}
$$

so that

$$
X^{o}(s)=-\frac{c}{2 a}
$$

Theorem 5.7 Let $T_{i}(s) \in R H_{\infty}, i=1,2$, be two given scalar functions. If $T_{2}(s)$ does not have zeros on the extended imaginary axis, then there exists a function $Q^{\circ}(s) \in$ $R H_{\infty}$ such that

$$
\left\|T_{1}(s)-T_{2}(s) Q^{o}(s)\right\|_{\infty}=\inf _{Q(s) \in R H_{\infty}}\left\|T_{1}(s)-T_{2}(s) Q(s)\right\|_{\infty}
$$

Proof Let $T_{2 i}(s) \in R H_{\infty}$ and $T_{2 o}(s) \in R H_{\infty}$ be two functions, inner and outer, respectively, such that $T_{2}(s)=T_{2 i}(s) T_{2 o}(s)$ (recall Theorem 2.10) and notice that, being $T_{2}(j \omega) \neq 0,0 \leq \omega \leq \infty$, it is $T_{2 o}^{-1}(s) \in R H_{\infty}$. Then, taking into account the identity $T_{2 i}^{\sim}(s) T_{2 i}(s)=I$, one obtains

$$
\begin{aligned}
\left\|T_{1}(s)-T_{2}(s) Q(s)\right\|_{\infty} & =\left\|T_{2 i}(s)\left[T_{2 i}^{-1}(s) T_{1}(s)-T_{2 o}(s) Q(s)\right]\right\|_{\infty} \\
& =\left\|T_{2 i}^{-1}(s) T_{1}(s)-T_{2 o}(s) Q(s)\right\|_{\infty}
\end{aligned}
$$

Define $F(s):=T_{2 i}^{-1}(s) T_{1}(s)$ and $X(s):=T_{2 o}(s) Q(s)$. Then the problem of minimizing $\|F(s)-X(s)\|_{\infty}$ with respect to $X(s) \in R H_{\infty}$ is equivalent to that of minimizing $\left\|T_{2 i}^{-1}(s) T_{1}(s)-T_{2 o}(s) Q(s)\right\|_{\infty}$ with respect to $Q(s) \in R H_{\infty}$, since, once the optimal $X^{o}(s) \in R H_{\infty}$ has been found, it is

$$
Q^{o}(s)=T_{2 o}^{-1}(s) X^{o}(s) \in R H_{\infty}
$$

Therefore, the theorem follows from Theorem 5.6.
Remark 5.32 The request that the function $T_{2}(s)$ has no zeros on the extended imaginary axis is a sufficient condition only. Indeed, consider the case in which

$$
T_{2}(s)=\frac{1}{s+1}
$$

while $T_{1}(s)$ is an element of $R H_{\infty}$. By defining $\hat{T}_{1}:=\lim _{s \rightarrow \infty} T_{1}(s)$, one has

$$
\begin{aligned}
\left\|T_{1}(s)-T_{2}(s) Q(s)\right\|_{\infty} & \geq \lim _{s \rightarrow \infty}\left\|T_{1}(s)-T_{2}(s) Q(s)\right\| \\
& \geq\left|\hat{T}_{1}\right|, \quad \forall Q(s) \in R H_{\infty}
\end{aligned}
$$

so that

$$
\inf _{Q(s) \in R H_{\infty}}\left\|T_{1}(s)-T_{2}(s) Q(s)\right\|_{\infty}:=\alpha \geq\left|\hat{T}_{1}\right|
$$

Chosen

$$
\hat{Q}(s):=\frac{T_{1}(s)-\hat{T}_{1}}{T_{2}(s)}
$$

it is easy to verify that such a function belongs to $R H_{\infty}$ and

$$
\alpha \leq\left\|T_{1}(s)-T_{2}(s) \hat{Q}(s)\right\|_{\infty}=\left|\hat{T}_{1}\right| \leq \alpha
$$

which implies $\hat{Q}(s)=Q^{o}(s)$.
Remark 5.33 In view of the above results a procedure for the computation of $Q^{\circ}(s)$ can be established as follows.

1) Find a double coprime factorization in $R H_{\infty}$ of $P_{22}(s)$, thus getting $M(s), \hat{M}(s), \hat{Y}(s)$.
2) Compute $T_{i}(s), i=1,2,3$ from eqs.(5.156)-(5.158). If no zero of $T_{4}(s):=T_{2}(s) T_{3}(s)$ lies on the complete imaginary axis, go to point 3 ) below, otherwise stop.
3) Find an inner-outer factorization of $T_{4}(s)$, so that $T_{4}(s)=T_{4 i}(s) T_{4 o}(s)$.
4) Find the antistable and strictly proper part $F_{a}(s)$ of $F(s):=T_{4 i}^{-1}(s) T_{1}(s)$.
5) Find a minimal realization $(A, B, C)$ of $F_{a}(s)$.
6) Solve the Lyapunov equations (5.165) and (5.166).
7) Find the maximum eigenvalue $\lambda_{M}^{2}$ of $P_{o} P_{r}$ and $f(s)$ and $g(s)$ from eqs. (5.167)-(5.170).
8) Compute $X^{o}(s)$ (eq. (5.172)).
9) $\operatorname{Set} Q^{o}(s)=T_{4 o}^{-1}(s) X^{o}(s)$.

Example 5.14 Consider the control system shown in fig. 5.47, where

$$
G_{n}(s)=\frac{s-1}{(s-2)(s+1)}
$$

A controller $K(s)$ is sought which stabilizes the system corresponding to the widest possible class of perturbations $\Delta(s)$ of the form

$$
\Delta(s) \in\left\{\Delta(s) \mid \Delta(s) \in R H_{\infty}, \quad\|\Delta(s)\|_{\infty}<\alpha\right\}
$$



Figure 5.47: A control system with additive perturbations
that is, corresponding to the maximum possible value of the scalar $\alpha$. Letting

$$
P(s):=\left[\begin{array}{cc}
0 & I \\
I & G_{n}(s)
\end{array}\right]
$$

the problem can be solved, thanks to the here presented results, by following the procedure outlined in Remark 5.33. If the adopted (minimal) realization of $G_{n}(s)$ is

$$
G_{n}(s)=\left[\begin{array}{rr|r}
0 & 1 & 0 \\
2 & 1 & 1 \\
\hline-1 & 1 & 0
\end{array}\right]:=\Sigma\left(A_{c}, B_{c}, C_{c}, 0\right)
$$

and the chosen matrices $F$ and $H$ are

$$
F=\left[\begin{array}{ll}
-3 & -3
\end{array}\right], \quad H=\left[\begin{array}{ll}
-3 & -6
\end{array}\right]^{\prime}
$$

(observe that $A_{c}+B_{c} F$ and $A_{c}+H C_{c}$ are stable), one obtains

$$
M(s)=\hat{M}(s)=\frac{s-2}{s+1}, \quad \hat{Y}(s)=-\frac{27}{s+1}
$$

Further, from eqs. (5.156)-(5.158) it results

$$
T_{1}(s)=-27 \frac{s-2}{(s+1)^{2}}, \quad T_{2}(s)=T_{3}(s)=\frac{s-2}{s+1}
$$

so that

$$
T_{4}(s)=\frac{(s-2)^{2}}{(s+1)^{2}}
$$

The function $T_{4}(s)$ has no zeros on the complete imaginary axis. Therefore, starting from its inner-outer factorization given by

$$
T_{4 i}(s)=\frac{(s-2)^{2}}{(s+2)^{2}}, \quad T_{4 o}(s)=\frac{(s+2)^{2}}{(s+1)^{2}}
$$

it follows

$$
F(s)=-27 \frac{(s+2)^{2}}{(s+1)^{2}(s-2)}
$$

A realization of the antistable part $F_{a}(s)$ of $F(s)$ is

$$
F_{a}(s)=\left[\begin{array}{r|r}
2 & 8 \\
\hline-6 & 0
\end{array}\right]
$$

so that, from eqs. (5.165)-(5.170), one obtains

$$
\begin{gathered}
P_{r}=16, \quad P_{o}=9, \quad \lambda_{M}=12, \quad \beta=1, \quad \chi=\frac{4}{3} \\
f(s)=\frac{8}{s+2}, \quad g(s)=-\frac{8}{s-2}
\end{gathered}
$$

Then, from eq. (5.172) it follows

$$
X^{o}(s)=12 \frac{(s+2)(s+7 / 4)}{(s+1)^{2}}
$$

which implies

$$
Q^{o}(s)=12 \frac{s+7 / 4}{s+2}
$$

Finally, from eq. (5.155) and by taking into account what has been shown in the proof of Theorem 3.4, one obtains

$$
K^{o}(s)=-12 \frac{s+1}{s-7}
$$

Therefore, the widest class of perturbations corresponding to which stability is guaranteed is characterized by the value $\alpha=\lambda_{M}^{-1}=1 / 12$. Compare this conclusion with the one given in Example 5.7.

### 5.7 Notes and references

Since the pioneering work of Zames [64], a great deal of attention has been devoted to robust control and related topics. The material collected in Sections 5.1 and 5.2 partly concerns classical issues of control theory suitably revisited so as to account for the basic instances of robustness. In particular, reference has been made to the books of Francis [19], Doyle et al. [16] and Maciejowski [43]. Further reading on related topics are Ackermann [1], Doyle [13] and Stein and Doyle [57]. No mention has been done to the somehow more realistic case of structured plant perturbations. The interested readers are referred to the papers by Doyle et al. [18], Doyle [14] and [15]. Sections 5.3, 5.4 are mainly based on the paper by Doyle et al. [17]. Again, the results concerning the parametrization of the admissible controllers follow the lines traced in the paper by Mita et al. [44]. More on the connections between differential games and $H_{\infty}$ theory can be found in the book by Basar and Bernhard [2]. The robust stabilization problem of Remark 5.11 has been previously faced by Barmish [3], Khargonekar et al. [31] and Haddad and Bernstein [27]. The design of $\mathrm{RH}_{2}$ filters in the presence of uncertainties (as in Remark 5.19) has been studied by Petersen and McFarlane [51] and Bolzern et al. [8], whereas the same problem in the $R H_{\infty}$ context was tackled by De Souza et al. [12] and Fu et al. [20]. The paper by Safonov and Limebeer [56] has been exploited in writing down Remark 5.24 in Section 5.5. The rest of the section still relies on the paper by Doyle et al. [17]. A recent book on the $H_{\infty}$ control problem is that by Stoorvogel [59] who has also explored the so called singular problem [58]. Finally, Section 5.6 deeply exploits the content of the book by Francis [19]. The here neglected approaches which rely on the gap metric and the polynomial framework are exploited in the paper by Georgiou and Smith [21] and by Kwakernaak [36].

## Chapter 6

## Nonclassical Problems in $\mathrm{RH}_{2}$ and $R H_{\infty}$

### 6.1 Introduction

This chapter introduces nonclassical problems in $\mathrm{RH}_{2}$ and $\mathrm{RH} \mathrm{m}_{\infty}$ spaces which, roughly speaking, can not be solved in general by the machinery provided in the previous chapters. For a given stable transfer function the $\mathrm{RH}_{2}$ norm is a measure of the level of a fixed input disturbance present in the output. The minimization of the $R H_{2}$ norm preserves the system against external disturbances. On the other hand, the $R H_{\infty}$ norm of the same transfer function kept bounded above by a certain prescribed value $\gamma>0$, imposes to the system a stable behavior against unmodeled dynamics with $R H_{\infty}$ norm less than $1 / \gamma$. The key observation is that $R H_{2}$ and $R H_{\infty}$ norms compete one with other. Indeed, when $\gamma \rightarrow \infty$ the central $R H_{\infty}$ controller approaches the optimal $\mathrm{RH}_{2}$ controller, meaning that a level $\gamma<\infty$ can be imposed only at the expense of some performance level. Being so, it is in many cases important to determine a controller which, in some sense, expresses a desired tradeoff between both norms. This characterizes the so called mixed $R \mathrm{H}_{2} / R H_{\infty}$ optimal control problem. Its solution can not be addressed by means of classical methods based on Riccati equation solvers. Accordingly, our attention has to be moved to other numerical tools as for instance convex programming methods discussed in Appendix I.

### 6.2 Parameter Space Optimization

In this section, basic control synthesis problems involving stability, $R H_{2}$ and $R H_{\infty}$ optimization are analyzed in the parameter space generated by the free elements of the feedback law. The main idea is to convert such nonconvex problems into convex ones in order to determine their global solutions by means of very powerful numeric procedures. First of all we need to introduce some concepts and definitions which are also discussed in Appendix H.

Definition 6.1 (Convex sets) $A$ set $\Omega$ in $R^{n \times m}$ is convex if $\forall X_{1}, X_{2} \in \Omega$ the point $X=\alpha X_{1}+(1-\alpha) X_{2} \in \Omega$ for every $\alpha \in[0,1]$.

Definition 6.2 (Convex functions) A function $f(\cdot) \in R$ defined in a convex set $\Omega$ is convex if $\forall X_{1}, X_{2} \in \Omega$ and $X=\alpha X_{1}+(1-\alpha) X_{2}$ there holds $f(X) \leq \alpha f\left(X_{1}\right)+$ $(1-\alpha) f\left(X_{2}\right)$ for every $\alpha \in[0,1]$.

There is another characterization of convexity which is in many instances simpler to apply. A real valued function $f(\cdot)$ defined in a convex set $\Omega \subset R^{n \times m}$ is convex if and only if for each $X_{0} \in \Omega$ there exists a matrix $\Lambda_{0}$ of appropriate dimension such that

$$
\begin{equation*}
f(X) \geq f\left(X_{0}\right)+<\Lambda_{0}, X-X_{0}> \tag{6.1}
\end{equation*}
$$

for all $X \in \Omega$. In the above inequality (6.1) the inner product of matrices is defined by $\langle\Lambda, X\rangle$ : $=\operatorname{trace}\left[\Lambda^{\prime} X\right]$ which induces the Frobenius norm. The set $\partial f\left(X_{0}\right)$ of all matrices $\Lambda_{0}$ satisfying (6.1) is called the subdifferential of $f(\cdot)$ at $X=X_{0}$. If $f(\cdot)$ is differentiable at $X=X_{0}$ then $\Lambda_{0}=\nabla f\left(X_{0}\right) \in \partial f\left(X_{0}\right)$ is unique.

Convexity is a key concept because any relative minimum of a convex programming problem formulated as to minimize $f(X)$ over a closed convex subset of $\Omega$ is a global minimum. Additionally, given a convex set $\Omega$ and a point $X_{0} \notin \Omega$, it is always possible to determine an hyperplane which separates $X_{0}$ from $\Omega$. The same result can be generalized to cope with $X_{0}$ in the boundary of $\Omega$.

When dealing with linear system stability and optimization, one of the most important sets to be handled is $\mathcal{P}$, composed by all square, real, symmetric and positive semidefinite matrices with fixed and know dimension. Naturally, $\mathcal{P}$ is a subset of $\mathcal{S}$, the set of all real symmetric matrices with the same dimension and $\mathcal{P}$ is a convex set. To show this, let us recall that $X \in \mathcal{P}$ if and only if for any vector $x$ of appropriate dimension $x^{\prime} X x \geq 0$. Hence, for every $X_{1}, X_{2} \in \mathcal{P}$, the convex combination of them satisfies

$$
x^{\prime} X x=\alpha x^{\prime} X_{1} x+(1-\alpha) x^{\prime} X_{2} x
$$

for all $\alpha \in[0,1]$ implying that $X \in \mathcal{P}$. The convexity of $\mathcal{P}$ follows from Definition 6.1. This result is a particular case of a more general one given in the next lemma.

Lemma 6.1 Suppose the matrix function $\mathcal{A}(X)$ is affine and its range is contained in $\mathcal{S}$. The set of all matrices $X$ such that $\mathcal{A}(X)$ is positive semidefinite is convex.

Proof Since $\mathcal{A}(X)$ is affine, for all $\alpha \in[0,1]$ it follows that the convex combination $X$ of arbitrary matrices $X_{1}, X_{2}$ satisfies

$$
\begin{equation*}
\mathcal{A}(X)=\alpha \mathcal{A}\left(X_{1}\right)+(1-\alpha) \mathcal{A}\left(X_{2}\right) \tag{6.2}
\end{equation*}
$$

Using the fact that the range of $\mathcal{A}(\cdot)$ is a subset of $\mathcal{S}$, then $\mathcal{A}(X) \geq 0$ is equivalent to $x^{\prime} \mathcal{A}(X) x \geq 0$ for all vectors $x$ of compatible dimension. The above equality yields

$$
\begin{equation*}
x^{\prime} \mathcal{A}(X) x=\alpha x^{\prime} \mathcal{A}\left(X_{1}\right) x+(1-\alpha) x^{\prime} \mathcal{A}\left(X_{2}\right) x \tag{6.3}
\end{equation*}
$$

and so, $\mathcal{A}(X) \geq 0$ whenever $\mathcal{A}\left(X_{1}\right) \geq 0$ and $\mathcal{A}\left(X_{2}\right) \geq 0$, proving thus the lemma proposed.

Remark 6.1 Lemma 6.1 says that the set of all matrices $X$ such that $\mathcal{A}(X) \in \mathcal{P}$ is convex. In particular, the choice $\mathcal{A}(X)=X$ shows once again that $\mathcal{P}$ is convex. Under the same assumptions of Lemma 6.1, the set of all matrices $X$ such that $\mathcal{A}(X) \leq 0$ is also convex. The real valued function $f(X)$ defined as

$$
f(X)=\max _{i} \lambda_{i}[\mathcal{A}(X)]
$$

with $\mathcal{A}(X)$ affine and with range in $\mathcal{S}$ is convex as well. To show this, considering $X$ as a convex combination of arbitrary $X_{1}, X_{2}$, the conclusion is that

$$
\begin{aligned}
f(X) & =\max _{x \neq 0} \frac{x^{\prime} \mathcal{A}(X) x}{x^{\prime} x} \\
& \leq \alpha \max _{x \neq 0} \frac{x^{\prime} \mathcal{A}\left(X_{1}\right) x}{x^{\prime} x}+(1-\alpha) \max _{x \neq 0} \frac{x^{\prime} \mathcal{A}\left(X_{2}\right) x}{x^{\prime} x} \\
& \leq \alpha f\left(X_{1}\right)+(1-\alpha) f\left(X_{2}\right)
\end{aligned}
$$

holds for all $\alpha \in[0,1]$.
Remark 6.2 It is important to keep clear that the range of $\mathcal{A}(X)$ being a subset of $\mathcal{S}$ is essential to get the above result. In fact, consider $\mathcal{A}(X)$ an affine matrix function but with its range not included in $\mathcal{S}$. This means that $\mathcal{A}(X)$ is not necessarily symmetric and the set of all $n \times n$ matrices $X$ such that $\operatorname{Re}\left[\lambda_{i}(\mathcal{A}(X))\right]<0, \quad i=1, \cdots, n$ is not convex in general. This can be verified by means of a simple counterexample. For $\mathcal{A}(X)=X$ let us take

$$
X_{1}=\left[\begin{array}{rr}
-1 & 4 \\
0 & -1
\end{array}\right], \quad X_{2}=\left[\begin{array}{rr}
-1 & 0 \\
4 & -1
\end{array}\right]
$$

and

$$
X=0.5 X_{1}+0.5 X_{2}=\left[\begin{array}{rr}
-1 & 2 \\
2 & -1
\end{array}\right]
$$

It is clear that matrices $X_{1}$ and $X_{2}$ belong to the previously defined set. However, matrix $X$ which is a convex combination of $X_{1}$ and $X_{2}$ does not belong to it. The conclusion is that the set of all $2 \times 2$ matrices under consideration is not convex.

Remark 6.3 Notice that the convex function $f(X)$ introduced in Remark 6.1 is not affine. It is not difficult to making use of inequality (6.1) for the determination of matrix $\Lambda_{0}$. As an example, consider the simplest case where $\mathcal{A}(X)=X \in \mathcal{S}$. Given $X_{0} \in \mathcal{S}$ let $x_{0}$ be an unitary norm eigenvector associated to its largest eigenvalue. Obviously the equality $X_{0} x_{0}=f\left(X_{0}\right) x_{0}$ holds and for all $X \in \mathcal{S}$ we have

$$
\begin{aligned}
f(X) & =\max _{\|x\|=1} x^{\prime} X x \\
& \geq x_{0}^{\prime} X x_{0} \\
& \geq f\left(X_{0}\right)+<x_{0} x_{0}^{\prime}, X-X_{0}>
\end{aligned}
$$

Comparing this inequality with (6.1) it is apparent that $\Lambda_{0}=x_{0} x_{0}^{\prime}$. Function $f(\cdot)$ is not differentiable unless the multiplicity of the maximum eigenvalue of $X_{0}$ is one. In this case $\Lambda_{0}=x_{0} x_{0}^{\prime} \in \partial f\left(X_{0}\right)$ is the only matrix satisfying (6.1) at $X=X_{0}$.

Another important function is $f(X)=\operatorname{trace}\left[B^{\prime} X^{-1} B\right]$, defined for all nonsingular $X \in \mathcal{P}$. This function is differentiable and convex. Let us show this by means of inequality (6.1). Considering arbitrary nonsingular matrices $X, X_{0} \in \mathcal{P}$ we have

$$
\begin{aligned}
f(X)= & \operatorname{trace}\left[B^{\prime} X^{-1} B\right] \\
= & f\left(X_{0}\right)-<X_{0}^{-1} B B^{\prime} X_{0}^{-1}, X-X_{0}>+ \\
& +\operatorname{trace}\left[B^{\prime} X_{0}^{-1}\left(X-X_{0}\right) X^{-1}\left(X-X_{0}\right) X_{0}^{-1} B\right] \\
\geq & f\left(X_{0}\right)+<-X_{0}^{-1} B B^{\prime} X_{0}^{-1}, X-X_{0}>
\end{aligned}
$$

and again $\Lambda_{0}=-X_{0}^{-1} B B^{\prime} X_{0}^{-1} \in \partial f\left(X_{0}\right)$ is unique. These examples show that inequality (6.1) is an easy and useful way to characterize convexity of matrix functions.

Given matrices $A$ and $Q=Q^{\prime}$, the sets of all matrices $X \in \mathcal{S}$ such that

$$
\mathcal{A}(X)=A^{\prime} X+X A+Q \leq 0
$$

or

$$
\mathcal{A}(X)=A X+X A^{\prime}+Q \leq 0
$$

are both convex. This is an immediate consequence of Lemma 6.1. Furthermore, given matrices $A, B$ and $Q=Q^{\prime}$, the sets of all matrices $X=X^{\prime}>0$ and $Y$ such that

$$
\mathcal{B}(X, Y)=(A+B Y)^{\prime} X+X(A+B Y)+Q \leq 0
$$

or

$$
\mathcal{B}(X, Y)=(A+B Y) X+X(A+B Y)^{\prime}+Q \leq 0
$$

are both nonconvex. The second one however can be converted into an equivalent convex set. Indeed, since $X$ is symmetric and positive definite the change of variables $Y:=Z X^{-1}$ provides

$$
\mathcal{A}(X, Z)=A X+B Z+X A^{\prime}+Z^{\prime} B^{\prime}+Q \leq 0
$$

which is convex because $\mathcal{A}(X, Z)$ is jointly affine in the variables $X, Z$. Finally, the set of all $(X, Y, Z)$ matrices such that $X=X^{\prime}>0$ and $Y^{\prime} X^{-1} Y-Z \leq 0$ is convex as well. This follows immediately from the Schur complement formula (recall Lemma B.14) which states that this set is equivalent to the set of all $(X, Y, Z)$ matrices such that

$$
X>0, \quad\left[\begin{array}{ll}
X & Y \\
Y^{\prime} & Z
\end{array}\right] \geq 0
$$

This is one of the most important result to be exhaustively used in the sequel.

### 6.2.1 Stabilizing controllers

A convex parametrization of all stabilizing matrix gains of a linear system is provided. Consider the following linear and time-invariant dynamic system

$$
\begin{equation*}
\dot{x}=A x+B_{2} u \tag{6.4}
\end{equation*}
$$

with $n$ states, $m$ inputs and where the state variable is available for feedback. When dealing with optimal control problems in $R H_{2}$ and $R H_{\infty}$ spaces, we have to restrict our attention to those $m \times n$ matrix gains $F$ such that with $u=F x$ the closed loop system is stable, that is

$$
\begin{equation*}
F \in \mathcal{K}_{c}:=\left\{F \in R^{m \times n}: A+B_{2} F \text { stable }\right\} \tag{6.5}
\end{equation*}
$$

The sentence " $A+B_{2} F$ stable" means that the closed loop system is internally stable, that is $\operatorname{Re}\left[\lambda_{i}\left(A+B_{2} F\right)\right]<0, \quad i=1, \cdots, n$. Since $\mathcal{A}(F)=A+B_{2} F$ is affine with respect to $F$ but its range is not in $\mathcal{S}$ the conclusion, as discussed before (recall Remark 6.2) is that the set $\mathcal{K}_{c}$ is not convex in general. Any optimal control problem formulated in the parameter space generated by the (free) elements of $F$ and having $\mathcal{K}_{c}$ as the feasible set is not convex. At most only local optimal solutions can be numerically determined by the machinery available in the literature to date. The fact that the set $\mathcal{K}_{c}$ is not convex, introduces one of the major difficulties to be faced in this section.

To circumvent this difficulty let us proceed as follows. Define the extended $p \times p$ and $p \times m$ matrices $(p:=n+m)$

$$
M_{c}:=\left[\begin{array}{cc}
A & B_{2}  \tag{6.6}\\
0 & 0
\end{array}\right], \quad N_{c}:=\left[\begin{array}{l}
0 \\
I
\end{array}\right]
$$

where the null space of $N_{c}^{\prime}$ is spanned by all $v \in R^{p}$ such that $N_{c}^{\prime} v=0$ or equivalently by all vectors in $R^{p}$ having the form $v=\left[\begin{array}{ll}x^{\prime} & 0\end{array}\right]^{\prime}$ with $x \in R^{n}$ arbitrary. Then, it is
apparent that the set $\mathcal{N}_{c}$ of vectors $v$ in the null space of $N_{c}^{\prime}$ with $\|v\|=1$ equals the set of all $v=\left[\begin{array}{ll}x^{\prime} & 0\end{array}\right]^{\prime}$ with $x \in R^{n}$ and $\|x\|=1$. Define also the $p \times p$ symmetric matrices $W$ and $Q_{c}$ partitioned as

$$
W:=\left[\begin{array}{ll}
W_{1} & W_{2}  \tag{6.7}\\
W_{2}^{\prime} & W_{3}
\end{array}\right], \quad Q_{c}:=\left[\begin{array}{cc}
Q_{1 c} & 0 \\
0 & 0
\end{array}\right]
$$

where in both matrices, the $(1,1)$ block has dimension $n \times n$.
Theorem 6.1 Assume $Q_{1 c}$ is a positive definite matrix and consider the set

$$
\begin{equation*}
\mathcal{C}_{c}:=\left\{W: W \geq 0, \quad v^{\prime} \Theta_{c}(W) v \leq 0, \forall v \in \mathcal{N}_{c}\right\} \tag{6.8}
\end{equation*}
$$

where $\Theta_{c}(W):=M_{c} W+W M_{c}^{\prime}+Q_{c}$. The following hold
a) $\mathcal{C}_{c}$ is a convex set.
b) Each $W \in \mathcal{C}_{c}$ is such that $W_{1}>0$.
c) $\mathcal{K}_{c}=\left\{W_{2}^{\prime} W_{1}^{-1}: W \in \mathcal{C}_{c}\right\}$.

Proof Each part of the theorem is proved separately. Notice that $W$ is symmetric and the matrix valued function $\Theta_{c}(\cdot)$ is affine.

Point a) Since the empty set is convex, let us proceed by considering $X_{1}$ and $X_{2}$ two arbitrary matrices belonging to $\mathcal{C}_{c}$. Taking $X=\alpha X_{1}+(1-\alpha) X_{2}$ with $\alpha \in[0,1]$ we first notice that $X \geq 0$ since the set of all positive semidefinite matrices is convex. Furthermore, taking into account that the matrix valued function $\Theta_{c}(X)$ is affine we get

$$
v^{\prime} \Theta_{c}(X) v=\alpha v^{\prime} \Theta_{c}\left(X_{1}\right) v+(1-\alpha) v^{\prime} \Theta_{c}\left(X_{2}\right) v \leq 0
$$

for all $v \in \mathcal{N}_{c}$. Consequently $X \in \mathcal{C}_{c}$.
Point b) Let us prove this point by contradiction. Assume for some $W \in \mathcal{C}_{c}$ there exists $x \neq 0 \in R^{n}$ such that $W_{1} x=0$. Since $W \geq 0$, using the partitioning (6.7) we must have $W_{2}^{\prime} x=0$. On the other hand $W \in \mathcal{C}_{c}$ and $v=\left[\begin{array}{cc}x^{\prime} & 0\end{array}\right]^{\prime} \in \mathcal{N}_{c}$ imply that

$$
\begin{aligned}
0 & \geq v^{\prime} \Theta_{c}(W) v \\
& \geq x^{\prime}\left(A W_{1}+W_{1} A^{\prime}+B_{2} W_{2}^{\prime}+W_{2} B_{2}^{\prime}+Q_{1 c}\right) x \\
& \geq x^{\prime} Q_{1 c} x
\end{aligned}
$$

which is an impossibility because $Q_{1 c}>0$.
Point c) The proof of this part is done by construction. First assume $\mathcal{K}_{c} \neq \emptyset$, take $F \in \mathcal{K}_{c}$ and recall that $Q_{1 c}$ is positive definite. From the Extended Lyapunov lemma (see Appendix C), there exists a symmetric positive definite solution to the linear equation

$$
\left(A+B_{2} F\right) P+P\left(A+B_{2} F\right)^{\prime}+Q_{1 c}=0
$$

Consequently, choosing

$$
W=\left[\begin{array}{cc}
P & P F^{\prime} \\
F P & F P F^{\prime}
\end{array}\right]
$$

it is readily seen that such a matrix $W$ is feasible, that is $W \in \mathcal{C}_{c}$ and $W_{2}^{\prime} W_{1}^{-1}=$ $F P P^{-1}=F$. Conversely, assume $\mathcal{C}_{c} \neq \emptyset$ and take $W \in \mathcal{C}_{c}$. From point b) we already know that $W_{1}>0$ and $\forall v \in \mathcal{N}_{c}$ yields

$$
\begin{align*}
0 & \geq v^{\prime} \Theta_{c}(W) v \\
& \geq x^{\prime}\left(A W_{1}+W_{1} A^{\prime}+B_{2} W_{2}^{\prime}+W_{2} B_{2}^{\prime}+Q_{1 c}\right) x \\
& \geq x^{\prime}\left[\left(A+B_{2} W_{2}^{\prime} W_{1}^{-1}\right) W_{1}+W_{1}\left(A+B_{2} W_{2}^{\prime} W_{1}^{-1}\right)^{\prime}+Q_{1 c}\right] x \tag{6.9}
\end{align*}
$$



Figure 6.1: Relationship between sets $\mathcal{K}_{c}$ and $\mathcal{C}_{c}$

The vector $x$ in the third inequality of (6.9) being arbitrary implies that $F=W_{2}^{\prime} W_{1}^{-1} \in$ $\mathcal{K}_{c}$. In fact, assume by contradiction that there exist an eigenvalue $\lambda$ and an eigenvector $x$ of matrix $A+B_{2} W_{2}^{\prime} W_{1}^{-1}$ such that $\operatorname{Re}(\lambda) \geq 0$ and $\|x\|=1$. From inequality (6.9) it follows that $2 \operatorname{Re}(\lambda) x^{\sim} W_{1} x+x^{\sim} Q_{1 c} x \leq 0$ which contradicts the assumption that $Q_{1 c}$ is positive definite. The proof is concluded because, from the above, if one of the sets $\mathcal{K}_{c}$ or $\mathcal{C}_{c}$ is empty then both are empty.

Figure 6.1 gives an interpretation of this result. Point c) of Theorem 6.1 provides a nonlinear mapping namely $W_{2}^{\prime} W_{1}^{-1}$ which generates the set $\mathcal{K}_{c}$ from all matrices $W$ in the convex domain $\mathcal{C}_{c}$. In some sense, the nonconvexity involved has been isolated in the nonlinearity $W_{2}^{\prime} W_{1}^{-1}$. Furthermore, notice that for $\forall v \in \mathcal{N}_{c}$ the quantity $v^{\prime} \Theta_{c}(W) v$ depends only upon the blocks $W_{1}$ and $W_{2}$ of matrix $W$. Consequently for given $W_{1}$ and $W_{2}$ such that $v^{\prime} \Theta_{c}(W) v \leq 0, \forall v \in \mathcal{N}_{c}$, it is always possible to select $W_{3} \geq W_{2}^{\prime} W_{1}^{-1} W_{2}$ (recall Appendix B) in order to have $W \geq 0$ and $W \in \mathcal{C}_{c}$. This degree of freedom will be important for solving optimal control problems in $\mathrm{RH}_{2}$ and $R H_{\infty}$ spaces.

Remark 6.4 From Theorem 6.1 it is clear that the pair $\left(A, B_{2}\right)$ is stabilizable if and only if $\mathcal{C}_{c} \neq \emptyset$. This follows from the fact that when $W$ varies in $\mathcal{C}_{c}$ then all stabilizing state feedback gains are generated from $F=W_{2}^{\prime} W_{1}^{-1}$. If such a matrix does not exist, it is clear that $\mathcal{C}_{c}=\emptyset$.

Let us now move our attention to the output feedback case. The linear system to be dealt with is

$$
\begin{align*}
\dot{x} & =A x+B_{2} u  \tag{6.10}\\
y & =C_{2} x \tag{6.11}
\end{align*}
$$

with $n$ states, $m$ inputs and $r$ outputs. The controller structure is given by

$$
\begin{align*}
\dot{\xi} & =\left(A+B_{2} F\right) \xi+L\left(C_{2} \xi-y\right)  \tag{6.12}\\
u & =F \xi \tag{6.13}
\end{align*}
$$

where both gains $F$ and $L$ have to be determined in order to assure the closed loop connection of the system and controller is internally stable. Simple algebraic manipulations (recall Section 4.1 and the separation property) bring to light that the dynamic matrix $A_{F}$ of the resulting system satisfies

$$
\begin{equation*}
\operatorname{det}\left[s I-A_{F}\right]=\operatorname{det}\left[s I-\left(A+B_{2} F\right)\right] \operatorname{det}\left[s I-\left(A+L C_{2}\right)\right] \tag{6.14}
\end{equation*}
$$

Hence, as far as stability is concerned, the determination of both gains which defines completely the dynamic controller (6.12) - (6.13) have to be such that $F \in \mathcal{K}_{c}$ and

$$
\begin{equation*}
L \in \mathcal{K}_{f}:=\left\{L \in R^{n \times r}: A+L C_{2} \text { stable }\right\} \tag{6.15}
\end{equation*}
$$

However, keeping in mind that the eigenvalues of any matrix and its transpose are the same, all elements of $\mathcal{K}_{f}$ can be generated by the dual version of Theorem 6.1. Indeed, following (6.6) let us define the $q \times q$ and $q \times r$ extended matrices ( $q:=n+r$ ) namely,

$$
M_{f}:=\left[\begin{array}{cc}
A & 0  \tag{6.16}\\
C_{2} & 0
\end{array}\right], \quad N_{f}:=\left[\begin{array}{l}
0 \\
I
\end{array}\right]
$$

It is clear that the null space of $N_{f}^{\prime}$ spanned by all vectors $v \in R^{q}$ such that $N_{f}^{\prime} v=0$, exhibits again the structure pointed out before, that is $v=\left[x^{\prime} 0\right]^{\prime}$ with $x \in R^{n}$. The set $\mathcal{N}_{f}$ is composed by all vectors $v \in R^{q}$ in the null space of $N_{f}^{\prime}$ with unitary norm. Partitioned accordingly to the dimensions of the plant state and output vectors, the symmetric matrices

$$
V:=\left[\begin{array}{ll}
V_{1} & V_{2}  \tag{6.17}\\
V_{2}^{\prime} & V_{3}
\end{array}\right], \quad Q_{f}:=\left[\begin{array}{cc}
Q_{1 f} & 0 \\
0 & 0
\end{array}\right]
$$

are on the basis for the next result, dual of the one provided in Theorem 6.1.
Theorem 6.2 Assume $Q_{1 f}$ is a positive definite matrix and consider the set

$$
\begin{equation*}
\mathcal{C}_{f}:=\left\{V: V \geq 0, \quad v^{\prime} \Theta_{f}(V) v \leq 0, \forall v \in \mathcal{N}_{f}\right\} \tag{6.18}
\end{equation*}
$$

where $\Theta_{f}(V):=M_{f}^{\prime} V+V M_{f}+Q_{f}$. The following hold
a) $\mathcal{C}_{f}$ is a convex set.
b) Each $V \in \mathcal{C}_{f}$ is such that $V_{1}>0$.
c) $\mathcal{K}_{f}=\left\{V_{1}^{-1} V_{2}: V \in \mathcal{C}_{f}\right\}$.

Remark 6.5 Form Theorem 6.2 it is clear that the pair $\left(A, C_{2}\right)$ is detectable if and only if $\mathcal{C}_{f} \neq \emptyset$. This is the dual of the property discussed in Remark 6.4.

Matrices $F$ and $L$ which define a stabilizing controller are thus determined from completely decoupled convex sets. This is possible because the dimension of the dynamic controller (6.12) - (6.13) equals the system dimension. The controller structure is based on the internal model of the plant and as expected, the error dynamics $\varepsilon:=\xi-x$ is completely defined by the poles of matrix $A+L C_{2}$ and so independent upon the choice of $F \in \mathcal{K}_{c}$. For this reason, the controller (6.12) - (6.13) is called observer-based controller and, as will be seen in the next sections, it plays a central role in nonclassical control design.

### 6.2.2 $\mathrm{RH}_{2}$ control design

Control design problems in $\mathrm{RH}_{2}$ space have been studied in Chapter 4. The standard problem, divided in three related topics has been solved. The main purpose of this section is to analyze once again the same problems but in the context of convex analysis. Under the same assumptions introduced in Chapter 4 it is possible to show that


Figure 6.2: The state feedback control system
the standard problem in $\mathrm{RH}_{2}$ can be converted to an equivalent convex programming problem.

Before giving an alternate solution to the Full information, Output estimation and Partial information problems, let us consider the following linear system

$$
\begin{align*}
& \dot{x}=A x+B_{1} w+B_{2} u  \tag{6.19}\\
& z=C_{1} x+D_{12} u  \tag{6.20}\\
& y=x \tag{6.21}
\end{align*}
$$

where, as indicated, the whole state vector is available for feedback. The structure of controller is simply a static state feedback gain so that

$$
\begin{equation*}
u=F x \tag{6.22}
\end{equation*}
$$

Calling $P_{E}(s)$ the transfer function of (6.19) - (6.21), the feedback system is depicted in fig. 6.2. Any state feedback gain matrix $F \in \mathcal{K}_{c}$ is called feasible. For any feasible gain, the transfer function $T(z, w ; s)$ belongs to $R H_{2}$ and its norm can be evaluated as being

$$
\begin{aligned}
\|T(z, w ; s)\|_{2}^{2} & =\left\|\left(C_{1}+D_{12} F\right)\left[s I-\left(A+B_{2} F\right)\right]^{-1} B_{1}\right\|_{2}^{2} \\
& =\operatorname{trace}\left[B_{1}^{\prime} P_{o} B_{1}\right] \\
& =\operatorname{trace}\left[\left(C_{1}+D_{12} F\right) P_{r}\left(C_{1}+D_{12} F\right)^{\prime}\right]
\end{aligned}
$$

where $P_{o}$ and $P_{r}$ are the unique solutions of the Lyapunov equations

$$
\begin{align*}
& 0=A_{c c}^{\prime} P_{o}+P_{o} A_{c c}+C_{c c}^{\prime} C_{c c}  \tag{6.23}\\
& 0=A_{c c} P_{r}+P_{r} A_{c c}^{\prime}+B_{1} B_{1}^{\prime} \tag{6.24}
\end{align*}
$$

with $A_{c c}:=A+B_{2} F$ and $C_{c c}:=C_{1}+D_{12} F$. We are now in position to introduce the problem to be solved.

Problem 6.1 (State feedback problem in $\mathrm{RH}_{2}$ ) Find a feasible state feedback gain $F$ which minimizes $\|T(z, w ; s)\|_{2}$. In other terms, find the global optimal solution of the problem

$$
\min \left\{\|T(z, w ; s)\|_{2}^{2}: F \in \mathcal{K}_{c}\right\}
$$

Before solving this problem we observe that it is a particular case of the Full information problem treated in Chapter 4 and Assumptions 4.1-4.2 of the Full information problem are again made. That is, i) The pair $\left(A, B_{2}\right)$ is stabilizable and no eigenvalue of the unobservable part of the pair $\left[\left(A-B_{2} D_{12}^{\prime} C_{1}\right),\left(I-D_{12} D_{12}^{\prime}\right) C_{1}\right]$
lies on the imaginary axis of the complex plane and ii) $D_{12}^{\prime} D_{12}=I$. Under these assumptions, $\mathcal{K}_{c} \neq \emptyset$ and the Full information problem and the State feedback problem shares the same optimal solution because as we already know (recall Theorem 4.1) at the optimal solution of the former, the gain corresponding to $w$ is zero. The solution of Problem 6.1 is thus

$$
\begin{equation*}
F^{o}=F_{2}=-B_{2}^{\prime} P_{2}-D_{12}^{\prime} C_{1} \tag{6.25}
\end{equation*}
$$

where $P_{2}$ is the symmetric, positive semidefinite and stabilizing solution of the Riccati equation (in the unknown $P$ )

$$
\begin{equation*}
0=P A_{c}+A_{c}^{\prime} P-P B_{2} B_{2}^{\prime} P+C_{1 c}^{\prime} C_{1 c} \tag{6.26}
\end{equation*}
$$

with $A_{c}:=A-B_{2} D_{12}^{\prime} C_{1}$ and $C_{1 c}:=\left(I-D_{12} D_{12}^{\prime}\right) C_{1}$. With $u=F_{2} x$ the closed loop minimum performance is given by

$$
\begin{equation*}
\min \|T(z, w ; s)\|_{2}^{2}=\operatorname{trace}\left[B_{1}^{\prime} P_{2} B_{1}\right] \tag{6.27}
\end{equation*}
$$

Remark 6.6 The solution of the State feedback problem (as well as the Full information problem) does not depend upon matrix $B_{1}$. Indeed, for any matrix $B_{1}$ the optimal feedback gain $F_{2}$ is completely characterized by means of (6.25) and (6.26) which do not depend of the aforementioned matrix. This aspect will be important in the sequel. In order to determine a convex problem equivalent to the State feedback problem it will be necessary to introduce an additional assumption, namely $B_{1} B_{1}^{\prime}>0$. To see that this can be done with no loss of generality, consider $B_{1} B_{1}^{\prime} \geq 0$ and define $\bar{B}_{1}:=\left[\begin{array}{ll}B_{1} & \sqrt{\epsilon} I\end{array}\right]$ where $\epsilon$ is a positive constant. Obviously matrix $\bar{B}_{1}$ exhibits the desired property $\bar{B}_{1} \bar{B}_{1}^{\prime}>0$ for any $\epsilon>0$ chosen. On the other hand, solving the State feedback problem we get

$$
\min \|\bar{T}(z, w ; s)\|_{2}^{2}=\min \|T(z, w ; s)\|_{2}^{2}+\epsilon \operatorname{trace}\left[P_{2}\right]
$$

That is, the minimum value of the modified transfer function differs from the previous one by an amount of order $\epsilon$ which can be made arbitrarily small by a proper choice of this parameter.

Remark 6.7 It is important to keep in mind that the set of admissible controllers for the State feedback and Full information problems in $\mathrm{RH}_{2}$ are quite different. In the Full information problem we have considered as admissible, any dynamic controller $K(s)$ satisfying Definition 4.1 and we have proven (recall Theorem 4.1) that the optimal controller is given by

$$
K(s)=K_{F I}^{o}(s)=\left[\begin{array}{c|cc}
\emptyset & \emptyset & \emptyset \\
\hline \emptyset & F_{2} & 0
\end{array}\right]
$$

In the State feedback problem stated before, we have considered as admissible only the static and stabilizing state feedback gains, namely $u=F x$ such that $F \in \mathcal{K}_{c}$. Since this structure is a particular case of the former and $K_{F I}^{\circ}(s)$ given above implies that $u=F_{2} x$ then $F=F_{2} \in \mathcal{K}_{c}$ is the optimal solution of the State feedback problem. The proof of this fact could be done in a different way which as expected is simpler than the proof of Theorem 4.1. Actually, making use of (6.23) we can restate the State feedback problem as to determine $F \in \mathcal{K}_{c}$ and $P=P^{\prime} \geq 0$ such that

$$
\min \left\{\operatorname{trace}\left[B_{1}^{\prime} P B_{1}\right]: 0=A_{c c}^{\prime} P+P A_{c c}+C_{c c}^{\prime} C_{c c}\right\}
$$

The Lagrangian function associated to this problem ( $\Lambda=\Lambda^{\prime}$ being the matrix of Lagrange multipliers associated to the equality constraint) is

$$
\mathcal{L}(F, P, \Lambda)=\operatorname{trace}\left[B_{1}^{\prime} P B_{1}\right]+\operatorname{trace}\left[\Lambda\left(A_{c c}^{\prime} P+P A_{c c}+C_{c c}^{\prime} C_{c c}\right)\right]
$$

So, the necessary conditions for optimality are readily obtained by simple differentiation of $\mathcal{L}(\cdot)$ with respect to the unknown matrices $F, P$ and $\Lambda$, providing

$$
\begin{aligned}
& 0=A_{c c}^{\prime} P+P A_{c c}+C_{c c}^{\prime} C_{c c} \\
& 0=A_{c c} \Lambda+\Lambda A_{c c}^{\prime}+B_{1} B_{1}^{\prime} \\
& 0=\left(F+D_{12}^{\prime} C_{1}+B_{2}^{\prime} P\right) \Lambda
\end{aligned}
$$

A possible solution to this set of nonlinear equations is obtained by first setting $F=$ $-B_{2}^{\prime} P-D_{12}^{\prime} C_{1}$. This equality together with the first equation imply that $P$ must solve the Riccati equation

$$
0=P A_{c}+A_{c}^{\prime} P-P B_{2} B_{2}^{\prime} P+C_{1 c}^{\prime} C_{1 c}
$$

Under the previous assumptions, this equation admits a stabilizing solution $P=P_{2}$ yielding $F=F_{2}$. Noticing that $F_{2} \in \mathcal{K}_{\mathrm{c}}$, the second equation has a solution with respect to $\Lambda$ namely,

$$
\Lambda_{2}=\int_{0}^{\infty} e^{\left(A+B_{2} F_{2}\right) t} B_{1} B_{1}^{\prime} e^{\left(A+B_{2} F_{2}\right)^{\prime} t} d t
$$

consequently the conclusion is that the triple ( $F_{2}, P_{2}, \Lambda_{2}$ ) satisfies the necessary conditions for optimality. It remains to prove that this is in fact the global solution of the proposed problem. To this end, consider $F \in \mathcal{K}_{c}$ an arbitrary stabilizing state feedback gain. Simple algebraic manipulations show that

$$
\left(A+B_{2} F\right)^{\prime}\left(P-P_{2}\right)+\left(P-P_{2}\right)\left(A+B_{2} F\right)+\left(F-F_{2}\right)^{\prime}\left(F-F_{2}\right)=0
$$

and (recall Appendix C) so

$$
P=P_{2}+\int_{0}^{\infty} e^{\left(A+B_{2} F\right)^{\prime} t}\left(F-F_{2}\right)^{\prime}\left(F-F_{2}\right) e^{\left(A+B_{2} F\right) t} d t \geq P_{2}
$$

For $u=F x$ the closed loop transfer function satisfies

$$
\begin{aligned}
\|T(z, w ; s)\|_{2}^{2} & =\operatorname{trace}\left[B_{1}^{\prime} P B_{1}\right] \\
& \geq \operatorname{trace}\left[B_{1}^{\prime} P_{2} B_{1}\right], \quad \forall F \in \mathcal{K}_{c}
\end{aligned}
$$

proving thus that the static matrix gain $F_{2}$ generated by the stabilizing solution $P_{2}$ of the Riccati equation (6.26) is the global optimum of the State feedback problem indeed. Once again, it is noticed that its solution does not depend upon matrix $B_{1}$.

However, for $B_{1}$ given it is in principle possible to determine $F \neq F_{2}$ which equals the global minimum. In fact, $P$ is equal to $P_{2}$ only if $F=F_{2}$ but trace $\left[B_{1}^{\prime} P B_{1}\right]$ may be equal to trace $\left[B_{1}^{\prime} P_{2} B_{1}\right]$ even though $F \neq F_{2}$. Let $F \in \mathcal{K}_{c}$ satisfying the necessary conditions for optimality with $\Lambda \geq 0$ be such that $\left(F-F_{2}\right) \Lambda=0$. Taking into account that

$$
\Lambda=\int_{0}^{\infty} e^{\left(A+B_{2} F\right) t} B_{1} B_{1}^{\prime} e^{\left(A+B_{2} F\right)^{\prime} t} d t
$$

together with the relationship between $P$ and $P_{2}$ we get the desired result

$$
\begin{aligned}
0 & =\operatorname{trace}\left[\left(F-F_{2}\right)^{\prime}\left(F-F_{2}\right) \Lambda\right] \\
& =\operatorname{trace}\left[\int_{0}^{\infty} B_{1}^{\prime} e^{\left(A+B_{2} F\right)^{\prime} t}\left(F-F_{2}\right)^{\prime}\left(F-F_{2}\right) e^{\left(A+B_{2} F\right) t} B_{1} d t\right] \\
& =\operatorname{trace}\left[B_{1}^{\prime}\left(P-P_{2}\right) B_{1}\right]
\end{aligned}
$$

From this it is apparent that the way to prevent this pathological situation is to impose $\Lambda>0, \forall F \in \mathcal{K}_{c}$. Thanks to the Extended Lyapunov lemma (recall Appendix C), this is always verified whenever $B_{1} B_{1}^{\prime}>0$, a condition just discussed in the previous remark.

It is simple to verify that Problem 6.1 is not convex (recall Appendix H). This conclusion follows immediately from the fact that the set of admissible gains $\mathcal{K}_{c}$ is not convex. Furthermore, its objective function is not convex as well. Defining the symmetric matrices

$$
R_{c}:=\left[\begin{array}{c}
C_{1}^{\prime}  \tag{6.28}\\
D_{12}^{\prime}
\end{array}\right]\left[\begin{array}{ll}
C_{1} & D_{12}
\end{array}\right], \quad Q_{c}:=\left[\begin{array}{cc}
B_{1} B_{1}^{\prime} & 0 \\
0 & 0
\end{array}\right]
$$

the next theorem provides a convex problem equivalent to Problem 6.1.
Theorem 6.3 (State feedback) Assume $B_{1} B_{1}^{\prime}$ is a positive definite matrix and consider the following convex programming problem

$$
\begin{equation*}
\bar{W}:=\operatorname{argmin}\left\{\operatorname{trace}\left[R_{c} W\right]: W \in \mathcal{C}_{c}\right\} \tag{6.29}
\end{equation*}
$$

Then, $F=\bar{W}_{2}^{\prime} \bar{W}_{1}^{-1}$ solves Problem 6.1.
Proof It is done in two main steps. First, consider matrix $\bar{W}$ given by

$$
\bar{W}:=\left[\begin{array}{cc}
\Lambda_{2} & \Lambda_{2} F_{2}^{\prime}  \tag{6.30}\\
F_{2} \Lambda_{2} & F_{2} \Lambda_{2} F_{2}^{\prime}
\end{array}\right]
$$

where (recall Remark 6.7) $\bar{W}_{1}=\Lambda_{2}>0$ due to the assumption $B_{1} B_{1}^{\prime}>0$. On the other hand, it is simple to verify that $\bar{W}_{2}^{\prime} \bar{W}_{1}^{-1}=F_{2}, \bar{W} \geq 0$ and for all $v \in \mathcal{N}_{c}$

$$
v^{\prime} \Theta_{c}(\bar{W}) v=x^{\prime}\left[\left(A+B_{2} F_{2}\right) \Lambda_{2}+\Lambda_{2}\left(A+B_{2} F_{2}\right)^{\prime}+B_{1} B_{1}^{\prime}\right] x=0
$$

implying that $\bar{W} \in \mathcal{C}_{c}$. The same matrix, together with (6.28) also provides

$$
\begin{aligned}
\operatorname{trace}\left[R_{c} \bar{W}\right] & =\operatorname{trace}\left[\left[\begin{array}{ll}
C_{1} & D_{12}
\end{array}\right] \bar{W}\left[\begin{array}{c}
C_{1}^{\prime} \\
D_{12}^{\prime}
\end{array}\right]\right] \\
& =\operatorname{trace}\left[\left(C_{1}+D_{12} F_{2}\right) \Lambda_{2}\left(C_{1}+D_{12} F_{2}\right)^{\prime}\right] \\
& =\operatorname{trace}\left[B_{1}^{\prime} P_{2} B_{1}\right] \\
& =\min \|T(z, w ; s)\|_{2}^{2}
\end{aligned}
$$

which means that matrix $\bar{W}$ is feasible and generates the optimal solution of Problem 6.1. It remains to prove that matrix $W$ is actually the optimal solution of the convex programming problem (6.29). To this end, let us keep in mind (recall Theorem 6.1) that all $W \in \mathcal{C}_{c}$ are such that $F=W_{2}^{\prime} W_{1}^{-1} \in \mathcal{K}_{c}$ and

$$
\|T(z, w ; s)\|_{2}^{2}=\operatorname{trace}\left[\left(C_{1}+D_{12} W_{2}^{\prime} W_{1}^{-1}\right) P_{r}\left(C_{1}+D_{12} W_{2}^{\prime} W_{1}^{-1}\right)^{\prime}\right]
$$

where $P_{r}>0$ solves the linear matrix equation

$$
0=\left(A+B_{2} W_{2}^{\prime} W_{1}^{-1}\right) P_{r}+P_{r}\left(A+B_{2} W_{2}^{\prime} W_{1}^{-1}\right)^{\prime}+B_{1} B_{1}^{\prime}
$$

However, any $W \in \mathcal{C}_{c}$ is such that $W_{1}>0$ and satisfies the inequality

$$
\left(A+B_{2} W_{2}^{\prime} W_{1}^{-1}\right) W_{1}+W_{1}\left(A+B_{2} W_{2}^{\prime} W_{1}^{-1}\right)^{\prime}+B_{1} B_{1}^{\prime} \leq 0
$$

which enables us to conclude that $W_{1} \geq P_{r}$. Furthermore, $W \geq 0$ and $W_{1}>0$ is equivalent to $W_{3} \geq W_{2}^{\prime} W_{1}^{-1} W_{2}$ and we finally get

$$
\begin{aligned}
\operatorname{trace}\left[R_{c} \bar{W}\right] & =\min \|T(z, w ; s)\|_{2}^{2} \\
& \leq\|T(z, w ; s)\|_{2}^{2} \\
& \leq \operatorname{trace}\left[\left(C_{1}+D_{12} W_{2}^{\prime} W_{1}^{-1}\right) W_{1}\left(C_{1}+D_{12} W_{2}^{\prime} W_{1}^{-1}\right)^{\prime}\right] \\
& \leq \operatorname{trace}\left[R_{c} W\right]-\operatorname{trace}\left[D_{12}\left(W_{3}-W_{2}^{\prime} W_{1}^{-1} W_{2}\right) D_{12}^{\prime}\right] \\
& \leq \operatorname{trace}\left[R_{c} W\right]
\end{aligned}
$$



Figure 6.3: The state feedback control problem
which being true for all feasible $W \in \mathcal{C}_{c}$ completes the proof.
Remark 6.8 Theorem 6.3 opens a very attractive way to solve optimal control problems in $\mathrm{RH}_{2}$ spaces by means of many powerful methods available to date. The convexity property is the key issue to be sure the global solution is always attained. However, in order to guarantee a selected numerical convex programming method can effectively solve Problem 6.1 it is necessary to prove that the feasible convex set $\mathcal{C}_{c}$ is bounded. Unfortunately, this is not true. In fact, the set $\mathcal{C}_{c}$ is a convex and unbounded cone. To show this take any $W \in \mathcal{C}_{c}$ and $\lambda \geq 1$. Then $\lambda W \geq 0$ and $\forall v \in \mathcal{N}_{c}$

$$
\begin{aligned}
v^{\prime} \Theta_{c}(\lambda W) v & =\lambda v^{\prime} \Theta_{c}(W) v+(1-\lambda) v^{\prime} Q_{c} v \\
& \leq \lambda v^{\prime} \Theta_{c}(W) v \leq 0
\end{aligned}
$$

implying that $\lambda W \in \mathcal{C}_{c}$.
However, a weaker condition exists to guarantee Problem 6.1 is numerically solvable even though the feasible set is not bounded. The situation is illustrated in fig. 6.3 where the set $\mathcal{C}_{c}$ is an unbounded cone but the objective function trace $\left[R_{c} W\right]$ has a global and finite optimum at $W=\bar{W}$. This occurs because the convex set

$$
\mathcal{C}_{c \mu}:=\mathcal{C}_{c} \cap\left\{W: \operatorname{trace}\left[R_{c} W\right] \leq \mu\right\}
$$

is bounded for all $\mu$ such that trace $\left[R_{c} \bar{W}\right]<\mu<\infty$. To show that Problem 6.1 has this property we proceed as follows. By contradiction assume that for some finite $\mu$ as specified before, the set $\mathcal{C}_{c \mu}$ is unbounded. In this case, thanks to its convexity, there exist $W \in \mathcal{C}_{c \mu}$ and $\tilde{W} \neq 0$ such that $W+\lambda \tilde{W} \in \mathcal{C}_{c \mu}$ for $\lambda>0$ arbitrarily large. Then, $\tilde{W}$ must satisfy

$$
\begin{aligned}
\operatorname{trace}\left[R_{c} \tilde{W}\right] & =0 \\
\tilde{W} & \geq 0 \\
v^{\prime}\left(M_{c} \tilde{W}+\tilde{W} M_{c}^{\prime}\right) v & \leq 0, \quad \forall v \in \mathcal{N}_{c}
\end{aligned}
$$

Using the structure of the symmetric positive semidefinite matrix $R_{c}$, the first two equations are equivalent to $\left[\begin{array}{ll}C_{1} & D_{12}\end{array}\right] \tilde{W}^{1 / 2}=0$ and all its solutions are also solutions of

$$
\left[\begin{array}{ll}
C_{1} & D_{12}
\end{array}\right] \tilde{W}=\left[\begin{array}{ll}
C_{1} & D_{12}
\end{array}\right]\left[\begin{array}{ll}
\tilde{W}_{1} & \tilde{W}_{2} \\
\tilde{W}_{2}^{\prime} & \tilde{W}_{3}
\end{array}\right]=0
$$

Making reference to Appendix $B$, this equation is solvable provided $C_{1 c} \tilde{W}_{1}=0$ yielding

$$
\tilde{W}=\left[\begin{array}{c}
I \\
-D_{12}^{\prime} C_{1}
\end{array}\right] \tilde{W}_{1}\left[\begin{array}{ll}
I & -C_{1}^{\prime} D_{12}
\end{array}\right]
$$

which replaced in the previous conditions allows us to express them in terms of $\tilde{W}_{1}$ only, that is we have to get $\tilde{W}_{1} \neq 0$ such that

$$
\begin{aligned}
C_{1 c} \tilde{W}_{1} & =0 \\
\tilde{W}_{1} & \geq 0 \\
A_{c} \tilde{W}_{1}+\tilde{W}_{1} A_{c}^{\prime} & \leq 0
\end{aligned}
$$

Once again making use of a result included in Appendix C, there is no matrix $\tilde{W}_{1} \neq 0$ satisfying these conditions provided the pair $\left(-A_{c}, C_{1 c}\right)$ is detectable. Therefore, to be sure the set $\mathcal{C}_{c \mu}$ is bounded we need to change the assumption i) of the State feedback problem to i) the pair ( $A, B_{2}$ ) is stabilizable and no eigenvalue of the unobservable part of the pair $\left(A_{c}, C_{1 c}\right)=\left[\left(A-B_{2} D_{12}^{\prime} C_{1}\right),\left(I-D_{12} D_{12}^{\prime}\right) C_{1}\right]$ lies on left part (including the imaginary axis) of the complex plane. If the detectability assumption is violated then $\tilde{W}_{1} \neq 0$ satisfying the above conditions may exist. In the example

$$
A_{c}=\left[\begin{array}{rr}
1 & 0 \\
2 & -1
\end{array}\right], \quad C_{1 c}=\left[\begin{array}{ll}
1 & 0
\end{array}\right], \quad \tilde{W}_{1}=\left[\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right]
$$

this occurs because the pair ( $-A_{c}, C_{1 c}$ ) fails to be detectable.
Example 6.1 To illustrate the geometry of the convex problem introduced in Theorem 6.3 , let us consider the system $P_{E}(s)$ given by

$$
\begin{aligned}
& \dot{x}=x+\sqrt{2} w+u \\
& z=\left[\begin{array}{l}
1 \\
0
\end{array}\right] x+\left[\begin{array}{l}
0 \\
1
\end{array}\right] u \\
& y=x
\end{aligned}
$$

The optimal solution of the State feedback problem is characterized by

$$
P_{2}=1+\sqrt{2}, \quad F_{2}=-(1+\sqrt{2}), \quad \Lambda_{2}=\frac{1}{\sqrt{2}}
$$

and $\|T(z, w ; s)\|_{2}^{2}=2(1+\sqrt{2})=4.82$. For this simple example, Problem 6.29 is written as

$$
\min \left\{W_{1}+W_{3}: W_{1}+W_{2}+1 \leq 0,\left[\begin{array}{ll}
W_{1} & W_{2} \\
W_{2} & W_{3}
\end{array}\right] \geq 0\right\}
$$

or in other terms, after using Schur complements

$$
\min \left\{W_{1}+\frac{W_{2}^{2}}{W_{1}}: W_{1}+W_{2}+1 \leq 0, \quad W_{1}>0\right\}
$$

The optimal solution is given by (recall the proof of Theorem 6.3)

$$
\bar{W}=\left[\begin{array}{cc}
\Lambda_{2} & \Lambda_{2} F_{2}^{\prime} \\
F_{2} \Lambda_{2} & F_{2} \Lambda_{2} F_{2}^{\prime}
\end{array}\right]=\left[\begin{array}{rr}
0.70 & -1.70 \\
-1.70 & 4.12
\end{array}\right]
$$

Figure 6.4 illustrates the feasible set which lies below the line indicated and the set of points such that the objective function is equal to $\mu=1,2,3$ and 4.82 . For the last value of $\mu=4.82$ the curve is tangent to the feasible region yielding the global minimum $W=W$. This example illustrates also that it is possible to eliminate matrix $W_{3}$ from the set of variables by setting $W_{3}=W_{2}^{\prime} W_{1}^{-1} W_{2}$. Of course this can only be done at the expense of changing a linear objective function to a nonlinear (although convex) objective function to be minimized.


Figure 6.4: Feasible region and objective function

Remark 6.9 (Linear matrix inequalities - LMI) The State feedback problem can also be stated only in terms of Linear matrix inequalities, that is the feasible set is defined by affine constraints only. The importance of this approach is mainly the standard formulation of many optimal control problems.

As we already know, the solution of the State feedback problem depends upon the determination of matrix $P_{2}$, the positive semidefinite stabilizing solution of the Riccati equation (6.26). To this purpose, we assume that the pair $\left(A_{c}, C_{1 c}\right)$ is observable to guarantee that $P_{2}>0$ and notice that all $P=P^{\prime}$ solving

$$
0 \geq P A_{c}+A_{c}^{\prime} P-P B_{2} B_{2}^{\prime} P+C_{1 c}^{\prime} C_{1 c}
$$

are such that $0<P_{2} \leq P$. To show this notice that for any $P$ satisfying this inequality, there exists $\bar{C}_{1 c}$ such that $\bar{C}_{1 c}^{\prime} \bar{C}_{1 c} \geq C_{1 c}^{\prime} C_{1 c}$ for which the equality holds. Consequently $F=-B_{2}^{\prime} P$ entails $A_{c}+B_{2} F$ stable and for all $x_{0} \in R^{n}$ we get (recall that $C_{1 c}^{\prime} D_{12}=0$ )

$$
\begin{aligned}
x_{0}^{\prime}\left(P-P_{2}\right) x_{0} & =x_{0}^{\prime} P x_{0}-\min _{F \in \mathcal{K}_{c}}\left\|\left(C_{1 c}+D_{12} F\right)\left[s I-\left(A_{c}+B_{2} F\right)\right]^{-1} x_{0}\right\|_{2}^{2} \\
& \geq x_{0}^{\prime} P x_{0}-\left\|\left(C_{1 c}-D_{12} B_{2}^{\prime} P\right)\left[s I-\left(A_{c}-B_{2} B_{2}^{\prime} P\right)\right]^{-1} x_{0}\right\|_{2}^{2} \\
& \geq x_{0}^{\prime}\left[\int_{0}^{\infty} e^{\left(A_{c}-B_{2} B_{2}^{\prime} P\right)^{\prime} t}\left(\bar{C}_{1 c}^{\prime} \bar{C}_{1 c}-C_{1 c}^{\prime} C_{1 c}\right) e^{\left(A_{c}-B_{2} B_{2}^{\prime} P\right) t} d t\right] x_{0} \\
& \geq 0
\end{aligned}
$$

and $P=P_{2}$ is the minimum of trace $\left[B_{1}^{\prime} P B_{1}\right]$ over all feasible $P$. Defining $X=P^{-1}$, the previous inequality can be rewritten as

$$
A_{c} X+X A_{c}^{\prime}-B_{2} B_{2}^{\prime}+X C_{1 c}^{\prime} C_{1 c} X \leq 0
$$

which after use of the Schur complements is equivalent to

$$
\mathcal{A}(X)=\left[\begin{array}{cc}
A_{c} X+X A_{c}^{\prime}-B_{2} B_{2}^{\prime} & X C_{1 c}^{\prime} \\
C_{1 c} X & -I
\end{array}\right] \leq 0
$$

Finally, from these calculations we have shown that matrix $P_{2}$ is the optimal global solution of the problem

$$
\min \left\{\operatorname{trace}\left[B_{1}^{\prime} X^{-1} B_{1}\right]: X>0, \mathcal{A}(X) \leq 0\right\}
$$



Figure 6.5: The output feedback control system

Clearly, the feasible set as well as the objective function are convex (recall Remark 6.3). Moreover, the feasible set is completely defined by affine functions only.

Remark 6.10 (Full information) As discussed before, the solution of the Full information problem is the same as the solution of the State feedback problem. For the sake of completeness the solution of the former is now restated. Consider the Problem 4.1 relative to system (4.10)-(4.14). Under Assumptions 4.1 and 4.2, it has the optimal solution:
a)

$$
\min \|T(z, w ; s)\|_{2}=\sqrt{\operatorname{trace}\left[R_{\mathrm{c}} \bar{W}\right]}
$$

b)

$$
K_{F I}^{o}(s)=\left[\begin{array}{c|cc}
\emptyset & \emptyset & \emptyset \\
\hline \emptyset & \bar{W}_{2}^{\prime} \bar{W}_{1}^{-1} & 0
\end{array}\right]
$$

where matrix $\bar{W}$ partitioned as indicated in (6.7) solves the convex programming problem (6.29).

Now, our attention is moved to the following situation. Consider the linear system described by

$$
\begin{align*}
& \dot{x}=A x+B_{1} w+B_{2} u  \tag{6.31}\\
& z=C_{1} x+D_{12} u  \tag{6.32}\\
& y=C_{2} x+D_{21} w \tag{6.33}
\end{align*}
$$

where only the partial information of the system state provided by the measured output $y$ is available for feedback. Accordingly to the results of Chapter 4, the controller structure is fixed as being

$$
K(s):=\left[\begin{array}{c|c}
A+B_{2} F+L C_{2} & -L  \tag{6.34}\\
\hline F & 0
\end{array}\right]
$$

which as indicated is completely parametrized by matrices $F$ and $L$. The rationale behind the choice of this structure is that the optimal solution of the Partial information problem can be obtained by a proper choice of the unknown matrices. With $P_{O}(s)$ being the transfer function of (6.31)-(6.33), the situation is illustrated by the block-scheme of fig. 6.5. The feedback connection has a state space realization of the form

$$
\Sigma_{F}:=\left[\begin{array}{c|c}
\tilde{A} & \tilde{B} \\
\hline \tilde{C} & 0
\end{array}\right]
$$

where the indicated matrices are

$$
\tilde{A}=\left[\begin{array}{cc}
A+B_{2} F & B_{2} F \\
0 & A+L C_{2}
\end{array}\right]
$$

and

$$
\tilde{B}=\left[\begin{array}{c}
B_{1} \\
-B_{1}-L D_{21}
\end{array}\right], \quad \tilde{C}=\left[\begin{array}{cc}
C_{1}+D_{12} F & D_{12} F
\end{array}\right]
$$

being thus clear that to preserve internal stability we have to consider only matrices $(F, L)$ such that $F \in \mathcal{K}_{c}$ and $L \in \mathcal{K}_{f}$. All pairs of matrices with this property are called feasible and for any feasible pair the transfer function $T(z, w ; s)$ belongs to $\mathrm{RH}_{2}$ and its norm is easily determined by

$$
\begin{aligned}
\|T(z, w ; s)\|_{2}^{2} & =\| \tilde{C}\left[s I-\tilde{A}^{-1} \tilde{B} \|_{2}^{2}\right. \\
& =\operatorname{trace}\left[\tilde{B}^{\prime} \tilde{P}_{o} \tilde{B}\right] \\
& =\operatorname{trace}\left[\tilde{C} \tilde{P}_{r} \tilde{C}^{\prime}\right]
\end{aligned}
$$

where $\tilde{P}_{o}$ and $\tilde{P}_{r}$ are the unique solutions of the Lyapunov equations associated with the closed-loop system, that is

$$
\begin{align*}
& 0=\tilde{A}^{\prime} \tilde{P}_{o}+\tilde{P}_{o} \tilde{A}+\tilde{C}^{\prime} \tilde{C}  \tag{6.35}\\
& 0=\tilde{A} \tilde{P}_{r}+\tilde{P}_{r} \tilde{A}^{\prime}+\tilde{B} \tilde{B}^{\prime} \tag{6.36}
\end{align*}
$$

For all feasible pairs $(F, L)$ the above equations are always solvable thanks to the stability of matrix $\tilde{A}$. As far as stability is concerned, the constraints on matrices $F$ and $L$ are completely decoupled. So, the control problem to be dealt with can be formulated as follows:

Problem 6.2 (Output feedback problem in $\mathrm{RH}_{2}$ ) Find a pair of feasible matrices $(F, L)$ which minimizes $\|T(z, w ; s)\|_{2}$. In other terms, find the global optimal solution of problem

$$
\min \left\{\|T(z, w ; s)\|_{2}^{2}: F \in \mathcal{K}_{c}, \quad L \in \mathcal{K}_{f}\right\}
$$

The difficulty to solve this problem stems from the fact that the objective function depends in a very unusual way on the unknown matrices. Fortunately, taking into account that the constraints are decoupled the optimality conditions can be expressed in terms of two separable subproblems, yielding each one, the optimal matrices $F$ and $L$. It should be clear that the Output feedback problem is a particular case of the Partial information problem introduced in Chapter 4 since the structure of the controller (6.34) has been fixed and completely described by means of only two unknown matrices. As a consequence, the Assumptions 4.6-4.9 are again considered, that is i) The pair ( $A, B_{2}$ ) is stabilizable and the pair ( $A, C_{2}$ ) is detectable, ii) $D_{12}^{\prime} D_{12}=I$, iii) The eigenvalues of the unobservable and unreachable part of the pairs $\left[\left(A-B_{2} D_{12}^{\prime} C_{1}\right)\left(I-D_{12} D_{12}^{\prime}\right) C_{1}\right]$ and $\left[\left(A-B_{1} D_{21}^{\prime} C_{2}\right), B_{1}\left(I-D_{21}^{\prime} D_{21}\right)\right]$ respectively, do not lie on the imaginary axis and iv) $D_{21} D_{21}^{\prime}=I$. Under these assumptions, $\mathcal{K}_{c} \neq \emptyset, \mathcal{K}_{f} \neq \emptyset$ and the optimal solution of the proposed problem is provided by Theorem 4.3, that is

$$
\begin{equation*}
F^{o}=F_{2}, \quad L^{o}=L_{2} \tag{6.37}
\end{equation*}
$$

which defines the optimal controller $K^{o}(s)$. With the optimal controller, the minimum cost is given by

$$
\min \|T(z, w ; s)\|_{2}^{2}=\left\|P_{c}(s) L_{2}\right\|_{2}^{2}+\left\|C_{1} P_{f}(s)\right\|_{2}^{2}
$$

where the transfer functions $P_{c}(s)$ and $P_{f}(s)$ are provided by Theorems 4.1 and 4.2 respectively. At this point, the key observation is that these two transfer functions defines completely the optimal solution of the problem under consideration and can be calculated by means of two decoupled problems. Indeed, simple calculations show that

$$
\begin{aligned}
P_{c}(s) L_{2} & =\left[\begin{array}{c|c}
A_{c}-B_{2} B_{2}^{\prime} P_{2} & L_{2} \\
\hline C_{1 c}-D_{12} B_{2}^{\prime} P_{2} & 0
\end{array}\right] \\
& =\left[\begin{array}{c|c}
A+B_{2} F_{2} & L_{2} \\
\hline C_{1}+D_{12} F_{2} & 0
\end{array}\right]
\end{aligned}
$$

and

$$
\begin{aligned}
C_{1} P_{f}(s) & =\left[\begin{array}{c|c}
A_{f}-\Pi_{2} C_{2}^{\prime} C_{2} & B_{1 f}-\Pi_{2} C_{2}^{\prime} D_{21} \\
\hline C_{1} & 0
\end{array}\right] \\
& =\left[\begin{array}{c|c}
A+L_{2} C_{2} & B_{1}+L_{2} D_{21} \\
\hline C_{1} & 0
\end{array}\right]
\end{aligned}
$$

which together with the solution of the State feedback problem, allows us to say that matrix $L^{o}=L_{2}$ solves the auxiliary state feedback problem

$$
\begin{equation*}
\left\|C_{1} P_{f}(s)\right\|_{2}^{2}=\min _{L \in \mathcal{K}_{f}}\left\|C_{1}\left[s I-\left(A+L C_{2}\right)\right]^{-1}\left(B_{1}+L D_{21}\right)\right\|_{2}^{2} \tag{6.38}
\end{equation*}
$$

while having obtained matrix $L_{2}$, it is apparent that the remaining unknown matrix $F_{2}$ is the optimal solution of another auxiliary state feedback problem

$$
\begin{equation*}
\left\|P_{c}(s) L_{2}\right\|_{2}^{2}=\min _{F \in \mathcal{K}_{c}}\left\|\left(C_{1}+D_{12} F\right)\left[s I-\left(A+B_{2} F\right)\right]^{-1} L_{2}\right\|_{2}^{2} \tag{6.39}
\end{equation*}
$$

The above solution to the Output feedback problem admits a very important interpretation. First it can be decomposed in two decoupled problems. The first one depends only on the system data and can be readily solved. Its solution provides the gain matrix $L_{2}$ which is used to define the objective function of the other optimization problem yielding the matrix gain $F_{2}$. This decomposition of the Output feedback problem is on the basis of what is called the Separation Principle which is valid for many optimal control problems with the controller structure given by (6.34).
Remark 6.11 Following the same lines of Remark 6.7, it is worth noticing that the solution of the Output feedback problem can also be obtained by means of mathematical programming arguments. The key observation is that all feasible controllers are constrained to have the structure (6.34) and the goal is to determine matrices $F \in \mathcal{K}_{c}, L \in \mathcal{K}_{f}$ and $\tilde{P} \geq 0$, such that

$$
\min \left\{\operatorname{trace}\left[\tilde{B}^{\prime} \tilde{P} \tilde{B}\right]: 0=\tilde{A}^{\prime} \tilde{P}+\tilde{P} \tilde{A}+\tilde{C}^{\prime} \tilde{C}\right\}
$$

Using $\tilde{\Lambda}=\tilde{\Lambda}^{\prime}$ as the matrix of Lagrange multipliers associated to the equality constraint, the partial derivatives of

$$
\mathcal{L}(F, L, \tilde{P}, \tilde{\Lambda}):=\operatorname{trace}\left[\tilde{B}^{\prime} \tilde{P} \tilde{B}\right]+\operatorname{trace}\left[\tilde{\Lambda}\left(\tilde{A}^{\prime} \tilde{P}+\tilde{P} \tilde{A}+\tilde{C}^{\prime} \tilde{C}\right)\right]
$$

give the necessary conditions for optimality

$$
\begin{aligned}
& 0=\tilde{A}^{\prime} \tilde{P}+\tilde{P} \tilde{A}+\tilde{C}^{\prime} \tilde{C} \\
& 0=\tilde{A} \tilde{\Lambda}+\tilde{\Lambda} \tilde{A}^{\prime}+\tilde{B} \tilde{B}^{\prime} \\
& 0=\frac{\partial \mathcal{L}}{\partial F} \\
& 0=\frac{\partial \mathcal{L}}{\partial L}
\end{aligned}
$$

which are simple to be solved from the observation that for $F=F_{2}$ and $L=L_{2}$ the solution of the two first conditions are

$$
\tilde{P}=\left[\begin{array}{cc}
P_{2} & 0 \\
0 & Y_{2}
\end{array}\right], \quad \tilde{\Lambda}=\left[\begin{array}{cc}
\Pi_{2}+X_{2} & -\Pi_{2} \\
-\Pi_{2} & \Pi_{2}
\end{array}\right]
$$

where matrices $X_{2}$ and $Y_{2}$ are solutions to the Lyapunov equations (in the unknown $X$ and $Y$ respectively)

$$
\begin{aligned}
& 0=\left(A+B_{2} F_{2}\right) X+X\left(A+B_{2} F_{2}\right)^{\prime}+L_{2} L_{2}^{\prime} \\
& 0=\left(A+L_{2} C_{2}\right)^{\prime} Y+Y\left(A+L_{2} C_{2}\right)+F_{2} F_{2}^{\prime}
\end{aligned}
$$

At the same point under consideration, the last two partial derivatives of the Lagrangian with respect to $F$ and $L$ are given by

$$
\frac{\partial \mathcal{L}}{\partial F}=\left(F-F_{2}\right) X_{2}, \quad \frac{\partial \mathcal{L}}{\partial L}=Y_{2}\left(L-L_{2}\right)
$$

showing that the pair ( $F_{2}, L_{2}$ ) satisfies the necessary conditions for optimality indeed. The importance of the calculations introduced here is to make explicit the structure of the matrix variables $\tilde{P}$ and $\tilde{\Lambda}$ at the optimal solution of the Output feedback control problem. This result will be of great importance in the sequel.

We already know that the Output feedback problem can be converted in two decoupled State feedback problems, hence it is not convex. Defining the extended matrices

$$
R_{f}:=\left[\begin{array}{c}
B_{1}  \tag{6.40}\\
D_{21}
\end{array}\right]\left[\begin{array}{ll}
B_{1}^{\prime} & D_{21}^{\prime}
\end{array}\right], Q_{f}:=\left[\begin{array}{cc}
C_{1}^{\prime} C_{1} & 0 \\
0 & 0
\end{array}\right]
$$

both problems (6.38) and (6.39) can be converted in two convex programming problems as is indicated in the next theorem.

Theorem 6.4 (Output feedback) The global optimal solution of the Output feedback problem can be calculated by means of the following procedure involving convex programming problems only.
a) Solve the convex programming problem

$$
\begin{equation*}
\bar{V}:=\operatorname{argmin}\left\{\operatorname{trace}\left[R_{f} V\right]: V \in \mathcal{C}_{f}\right\} \tag{6.41}
\end{equation*}
$$

b) Redefine the extended matrix

$$
Q_{c}:=\left[\begin{array}{cc}
\bar{V}_{1}^{-1} \bar{V}_{2} \bar{V}_{2}^{\prime} \bar{V}_{1}^{-1} & 0  \tag{6.42}\\
0 & 0
\end{array}\right]
$$

and solve the convex programming problem

$$
\begin{equation*}
\bar{W}:=\operatorname{argmin}\left\{\operatorname{trace}\left[R_{c} W\right]: W \in \mathcal{C}_{c}\right\} \tag{6.43}
\end{equation*}
$$

Then, matrices $F=\bar{W}_{2}^{\prime} \bar{W}_{1}^{-1}$ and $L=\bar{V}_{1}^{-1} \bar{V}_{2}$ solve Problem 6.2.
Proof It follows from the decomposition of Problem 6.2 in terms of (6.38) and (6.39). The result of Theorem 6.3 is used to convert them to equivalent convex programming problems.

Remark 6.12 Remark 6.6 applies to both optimization problems defined in Theorem 6.4. If matrices $C_{1}^{\prime} C_{1}$ and $L L^{\prime}$ are not strictly positive definite then $C_{1}^{\prime}$ and $L=\bar{V}_{1}^{-1} \bar{V}_{2}$ have to be replaced, with no loss of generality, by $\left[C_{1}^{\prime} \sqrt{\epsilon} I\right]$ and $[L \sqrt{\epsilon} I]$ respectively where $\epsilon$ is an arbitrarily small positive parameter.

Remark 6.13 It is important to discuss under which condition problem (6.41) can be effectively solved by a chosen convex programming procedure. As discussed in Remark 6.8, we have to impose conditions assuring that the set

$$
\mathcal{C}_{f \mu}:=\mathcal{C}_{f} \cap\left\{V: \operatorname{trace}\left[R_{f} V\right] \leq \mu\right\}
$$

is bounded for all $\mu$ such that trace $\left[R_{f} \bar{V}\right]<\mu<\infty$. Following the same reasoning, this occurs provided the pair ( $-A_{f}, B_{1 f}$ ) is stabilizable. Therefore to have both sets $\mathcal{C}_{c \mu}$ and $\mathcal{C}_{f \mu}$ bounded, the assumption iii) of the Output feedback problem has to be changed to iii) The eigenvalues of the unobservable and unreachable part of the pairs $\left[\left(A-B_{2} D_{12}^{\prime} C_{1}\right)(I-\right.$ $\left.\left.D_{12} D_{12}^{\prime}\right) C_{1}\right]$ and $\left[\left(A-B_{1} D_{21}^{\prime} C_{2}\right), B_{1}\left(I-D_{21}^{\prime} D_{21}\right)\right]$ respectively, do not lie on the left part (including the imaginary axis) of the complex plane.

Remark 6.14 It is simple to be verified that in the State feedback problem, the set of all $F \in \mathcal{K}_{c}$ such that $\|T(z, w ; s)\|_{2}$ is bounded above by a given positive scalar $\gamma$ can be generated by (recall the definition of the convex set $\mathcal{C}_{c \mu}$ in Remark 6.8)

$$
F=W_{2}^{\prime} W_{1}^{-1}, W \in \mathcal{C}_{c} \cap\left\{W: \operatorname{trace}\left[R_{c} W\right] \leq \gamma^{2}\right\}
$$

Our purpose now is to generalize this result to the Output feedback problem in the case that the matrix gain $L$ is fixed and equals the optimal value $L=L_{2}$. For any $F \in \mathcal{K}_{c}$ the transfer function $T(z, w ; s)$ belongs to $\mathrm{RH}_{2}$ and

$$
\|T(z, w ; s)\|_{2}^{2}=\operatorname{trace}\left[\tilde{C} \tilde{P}_{r} \tilde{C}^{\prime}\right]
$$

where $\tilde{P}_{r}$ is the solution of the linear equation (6.36). Simple algebraic manipulations show that

$$
\tilde{P}_{r}=\left[\begin{array}{cc}
\Pi_{2}+X & -\Pi_{2} \\
-\Pi_{2} & \Pi_{2}
\end{array}\right]
$$

where matrix $X$ is the solution of

$$
0=\left(A+B_{2} F\right) X+X\left(A+B_{2} F\right)^{\prime}+L_{2} L_{2}^{\prime}
$$

It is very interesting to compare this result to the ones given in Remark 6.11. Matrix $\tilde{P}_{r}$ has the same structure as matrix $\tilde{\Lambda}$ and the only difference, due to have considered now $F \in \mathcal{K}_{c}$ arbitrary, is restricted to the definition of matrix $X$ above which equals matrix $X_{2}$ provided $F=F_{2}$. Using this we also have

$$
\begin{aligned}
\|T(z, w ; s)\|_{2}^{2} & =\operatorname{trace}\left[C_{1} \Pi_{2} C_{1}^{\prime}\right]+\operatorname{trace}\left[\left(C_{1}+D_{12} F\right) X\left(C_{1}+D_{12} F\right)^{\prime}\right] \\
& =\operatorname{trace}\left[C_{1} \Pi_{2} C_{1}^{\prime}\right]+\left\|\left(C_{1}+D_{12} F\right)\left[s I-\left(A+B_{2} F\right)\right]^{-1} L_{2} D_{21}\right\|_{2}^{2}
\end{aligned}
$$

Finally, taking into account Theorem 6.4 it is clear that

$$
F=W_{2}^{\prime} W_{1}^{-1}, W \in \mathcal{C}_{c} \cap\left\{W: \operatorname{trace}\left[R_{c} W\right] \leq \gamma^{2}-\operatorname{trace}\left[C_{1} \Pi_{2} C_{1}^{\prime}\right]\right\}
$$

generates all matrices $F$ such that, in the Output feedback problem, the internal stability is assured while $\|T(z, w ; s)\|_{2}$ is bounded above by a given positive scalar $\gamma$.

Remark 6.15 (Output estimation) The Output estimation problem as defined in the Chapter 4 can also be solved by means of convex programming tools. The main observation is that Assumptions 4.3-4.5 of the Output estimation problem imply the assumptions of the Output feedback problem are all verified with $D_{12}=I$. Furthermore, the assumption $A_{c}=A-B_{2} C_{1}$ stable yields $P_{2}=0$ and the associated stabilizing matrix $F_{2}=-C_{1}$. Then, Problem 4.1 relative to system (4.58)-(4.60) has the optimal solution :
a)

$$
\min \|T(z, w ; s)\|_{2}=\sqrt{\operatorname{trace}\left[R_{f} \bar{V}\right]}
$$

b)

$$
K_{O E}^{o}(s)=\left[\begin{array}{c|c}
A-B_{2} C_{1}+\bar{V}_{1}^{-1} \bar{V}_{2} C_{2} & -\bar{V}_{1}^{-1} \bar{V}_{2} \\
\hline-C_{1} & 0
\end{array}\right]
$$

where matrix $\bar{V}$ partitioned as indicated in (6.17) solves the convex programming problem (6.41).

Remark 6.16 (Partial information) For completeness we give here the solution of the Partial information problem. The additional result is the value of the objective function written in terms of matrices provided in Theorem 6.4. Consider the system (4.92)-(4.94) and Assumptions 4.6-4.9 then Problem 4.1 has the following solution :
a)

$$
\min \|T(z, w ; s)\|_{2}^{2}=\operatorname{trace}\left[R_{f} \bar{V}\right]+\operatorname{trace}\left[R_{c} \bar{W}\right]
$$

b)

$$
K^{o}(s)=\left[\begin{array}{c|c}
A+B_{2} \bar{W}_{2}^{\prime} \bar{W}_{1}^{-1}+\bar{V}_{1}^{-1} \bar{V}_{2} C_{2} & -\bar{V}_{1}^{-1} \bar{V}_{2} \\
\hline \bar{W}_{2}^{\prime} \bar{W}_{1}^{-1} & 0
\end{array}\right]
$$

where matrices $\bar{W}$ and $\bar{V}$ solve the convex optimization problems introduced in Theorem 6.4. Finally, it is important to notice again that these matrices can be determined separately.

### 6.2.3 $R H_{\infty}$ control design

This section presents the convex analysis counterpart of Chapter 5. The State feedback and Output feedback control design problems are again addressed but in a slight different setting. Indeed, not only the set of all controllers which impose to the plant a certain $R H_{\infty}$ norm level $\gamma$ is obtained but we also address the problem of determining the controller such that $\gamma$ is minimized. Particular attention must be payed to the fact that in many instances, on the contrary of what has been done in Chapter 5, the constraints involving $R H_{\infty}$ norms are not taken strictly:

Let the system under consideration be defined as

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{6.44}\\
z & =C_{1} x+D_{12} u  \tag{6.45}\\
y & =x \tag{6.46}
\end{align*}
$$

and illustrated in fig. 6.2 for which the controller is simply given by $u=F x$ where matrix $F \in \mathcal{K}_{c}$ is such that $T(z, w ; s)$ belongs to $R H_{\infty}$. For all $F \in \mathcal{K}_{c}$, those that additionally satisfies the constraint $\|T(z, w ; s)\|_{\infty}<\gamma$ with $\gamma$ being a positive scalar are characterized from the result of Theorem 2.14 which states that there exists a symmetric and positive semidefinite stabilizing solution of the algebraic Riccati equation (in the unknown $S$ )

$$
\begin{equation*}
0=S A_{c \infty}+A_{c \infty}^{\prime} S+\gamma^{-2} S B_{1} B_{1}^{\prime} S+C_{c \infty}^{\prime} C_{c \infty} \tag{6.47}
\end{equation*}
$$

with $A_{c \infty}:=A+B_{2} F$ and $C_{c \infty}:=C_{1}+D_{12} F$. An equivalent condition is obtained from Theorem 2.15 which requires the existence of a symmetric and positive definite feasible solution to the Riccati inequality

$$
\begin{equation*}
S A_{c \infty}+A_{c \infty}^{\prime} S+\gamma^{-2} S B_{1} B_{1}^{\prime} S+C_{c \infty}^{\prime} C_{c \infty}<0 \tag{6.48}
\end{equation*}
$$

in this case no further stabilizing condition is needed. However, it is a simple matter to see (recall Remark 2.22) that $\left\|T^{\prime}(z, w ; s)\right\|_{\infty}=\|T(z, w ; s)\|_{\infty}$ meaning that the same set of controllers can also be characterized by the existence of a symmetric and positive semidefinite stabilizing solution of the algebraic Riccati equation (in the unknown $P$ )

$$
\begin{equation*}
0=A_{c \infty} P+P A_{c \infty}^{\prime}+\gamma^{-2} P C_{c \infty}^{\prime} C_{c \infty} P+B_{1} B_{1}^{\prime} \tag{6.49}
\end{equation*}
$$

or equivalently from Theorem 2.15 the existence of a symmetric and positive definite feasible solution to the Riccati inequality

$$
\begin{equation*}
A_{c \infty} P+P A_{c \infty}^{\prime}+\gamma^{-2} P C_{c \infty}^{\prime} C_{c \infty} P+B_{1} B_{1}^{\prime}<0 \tag{6.50}
\end{equation*}
$$

As well as, the set of all $F \in \mathcal{K}_{c}$ such that additionally $\|T(z, w ; s)\|_{\infty} \leq \gamma$ is completely characterized from the results of Theorem 2.16 which requires the existence of symmetric and positive definite feasible solutions to the nonstrict versions of the Riccati inequalities (6.48) and (6.50). Moreover, in the above Riccati inequalities each feasible solution are related one to the other by $P=\gamma^{2} S^{-1}$. Although equivalent, to our purpose, inequality ( 6.50 ) or its nonstrict version are more convenient to handle.

With $\gamma$ being an arbitrary and positive scalar (not fixed a priori), the pair ( $F, \gamma$ ) is called feasible whenever $(F, \gamma) \in \mathcal{K}_{\gamma c}$ where

$$
\begin{equation*}
\mathcal{K}_{\gamma c}:=\left\{(F, \gamma): F \in \mathcal{K}_{c}, \gamma>0,\|T(z, w ; s)\|_{\infty} \leq \gamma\right\} \tag{6.51}
\end{equation*}
$$

which makes clear that the set $\mathcal{K}_{\gamma c}$ is a subset of $\mathcal{K}_{c}$. Based on our previous discussion, it is now a fact that the above defined set is not convex. Fortunately it can be converted to an equivalent convex set as indicated in the next theorem.

Theorem 6.5 Assume $B_{1} B_{1}^{\prime}$ is a positive definite matrix and consider the set

$$
\begin{equation*}
\mathcal{C}_{\gamma c}:=\left\{(W, \mu): W \geq 0, \quad \mu>0, \quad v^{\prime} \Theta_{\gamma c}(W, \mu) v \leq 0, \forall v \in \mathcal{N}_{c}\right\} \tag{6.52}
\end{equation*}
$$

where $\Theta_{\gamma c}(W, \mu):=\Theta_{c}(W)+\mu^{-1} W R_{c} W$. The following hold
a) $\mathcal{C}_{\gamma c}$ is a convex set.
b) Each $(W, \mu) \in \mathcal{C}_{\gamma c}$ is such that $W_{1}>0$.
c) $\mathcal{K}_{\gamma c}=\left\{\left(W_{2}^{\prime} W_{1}^{-1}, \sqrt{\mu}\right):(W, \mu) \in \mathcal{C}_{\gamma c}\right\}$.

Proof Notice that function $\Theta_{\gamma c}(W, \mu)$ is not affine but it can be converted to an affine one by using Schur complements.

Point a) It is clear that we only need to prove that the set of all $(W, \mu)$ such that $v^{\prime} \Theta_{\gamma c}(W, \mu) v \leq 0, \forall v \in \mathcal{N}_{c}$ is convex. Since $R_{c} \geq 0$ and all $v \in \mathcal{N}_{c}$ can be obtained from $v=U_{c} x$ where $U_{c}=\left[\begin{array}{ll}I & 0\end{array}\right]^{\prime}$ and $\|x\|=1$, then this set is equivalently described as

$$
0 \geq U_{c}^{\prime} \Theta_{c}(W) U_{c}+\mu^{-1} U_{c}^{\prime} W R_{c}^{1 / 2} R_{c}^{1 / 2} W U_{c}
$$

which for $\mu>0$, after using of the Schur complement formula becomes

$$
\mathcal{A}(W, \mu)=\left[\begin{array}{cc}
U_{c}^{\prime} \Theta_{c}(W) U_{c} & U_{c}^{\prime} W R_{c}^{1 / 2} \\
R_{c}^{1 / 2} W U_{c} & -\mu I
\end{array}\right] \leq 0
$$

Since $\mathcal{A}(W, \mu)$ is affine, the convexity of $\mathcal{C}_{\gamma c}$ follows directly from Lemma 6.1.

Point b) Since for each $(W, \mu) \in \mathcal{C}_{\gamma c}$ the matrix $W$ is positive semidefinite, if there exists a nonzero vector $x \in R^{n}$ such that $W_{1} x=0$ then for the same vector $W_{2}^{\prime} x=0$, yielding $W v=0$ and obviously

$$
0 \geq v^{\prime} \Theta_{\gamma c}(W, \mu) v=\left\|B_{1}^{\prime} x\right\|^{2}
$$

which is an impossibility due to the fact that, by assumption, $B_{1} B_{1}^{\prime}>0$.
Point c) Suppose $\mathcal{K}_{\gamma c} \neq \emptyset$ and the pair $(F, \gamma)$ is feasible for the set $\mathcal{K}_{\gamma c}$. In this case, there exists a symmetric and positive definite matrix $P$ satisfying the Riccati inequality (recall Theorem 2.16)

$$
A_{c \infty} P+P A_{c \infty}^{\prime}+\gamma^{-2} P C_{c \infty}^{\prime} C_{c \infty} P+B_{1} B_{1}^{\prime} \leq 0
$$

Choosing $\mu=\gamma^{2}$ and

$$
W=\left[\begin{array}{cc}
P & P F^{\prime} \\
F P & F P F^{\prime}
\end{array}\right]
$$

it is a simple matter to verify that $(W, \mu) \in \mathcal{C}_{\gamma c}$ and $W_{2}^{\prime} W_{1}^{-1}=F$. Conversely, assume $\mathcal{C}_{\gamma c} \neq \emptyset$ and consider any $(W, \mu) \in \mathcal{C}_{\gamma c}$, for all $v \in \mathcal{N}_{c}$ we get

$$
\begin{align*}
0 \geq & v^{\prime} \Theta_{\gamma c}(W, \mu) v \\
\geq & x^{\prime}\left[\left(A+B_{2} W_{2}^{\prime} W_{1}^{-1}\right) W_{1}+W_{1}\left(A+B_{2} W_{2}^{\prime} W_{1}^{-1}\right)^{\prime}+\right. \\
& \left.+\mu^{-1} W_{1}\left(C_{1}+D_{12} W_{2}^{\prime} W_{1}^{-1}\right)^{\prime}\left(C_{1}+D_{12} W_{2}^{\prime} W_{1}^{-1}\right) W_{1}+B_{1} B_{1}^{\prime}\right] x \tag{6.53}
\end{align*}
$$

where the factorization is possible since from point b) it has been already proven that $W_{1}>0$. Choosing $P=W_{1}, F=W_{2}^{\prime} W_{1}^{-1}$ and $\gamma=\sqrt{\mu}$, this inequality assures that $F \in \mathcal{K}_{c}$ and the existence of a symmetric and positive definite matrix $P=W_{1}>0$ satisfying the Riccati inequality

$$
A_{c \infty} P+P A_{c \infty}^{\prime}+\gamma^{-2} P C_{c \infty}^{\prime} C_{c \infty} P+B_{1} B_{1}^{\prime} \leq 0
$$

Hence, using Theorem 2.16 it is verified that $\|T(z, w ; s)\|_{\infty} \leq \gamma$ which, from (6.51) implies that the pair $(F, \gamma)=\left(W_{2}^{\prime} W_{1}^{-1}, \sqrt{\mu}\right) \in \mathcal{K}_{\gamma c}$. From the above, if one set $\mathcal{K}_{\gamma c}$ or $\mathcal{C}_{\gamma c}$ is empty both are empty. The proof is then complete.

The joint convexity of the set $\mathcal{C}_{\gamma c}$ with respect to both variables $(W, \mu)$ is of great importance. To get some insight on this fact let us consider the problem of determining the feasible pair $(F, \gamma)$ which solves

$$
\begin{equation*}
\min \left\{\gamma:(F, \gamma) \in \mathcal{K}_{\gamma c}\right\} \tag{6.54}
\end{equation*}
$$

This nonconvex problem has not been directly addressed in Chapter 5 although it is possible, in principle, to get its solution iteratively with the results provided in Theorem 5.2 (recall the forthcoming Remark 6.20). It is equivalent to the convex programming problem

$$
\begin{equation*}
\min \left\{\mu:(W, \mu) \in \mathcal{C}_{\gamma c}\right\} \tag{6.55}
\end{equation*}
$$

Indeed, at the optimal solution the value of $\mu$ is the minimum which preserves feasibility implying that for $\gamma=\sqrt{\mu}$ there exists a matrix $F$ such that $(F, \gamma) \in \mathcal{K}_{\gamma c}$.

Let us now introduce the problem to be dealt with in the sequel. It resembles the Full information problem treated in Chapter 5. To this end it is first needed to put in evidence only the feasible pairs $(F, \gamma)$ such that $\|T(z, w ; s)\|_{\infty}<\gamma$, that is all $(F, \gamma) \in \operatorname{int} \mathcal{K}_{\gamma c}$, where

$$
\text { int } \mathcal{K}_{\gamma c}:=\left\{(F, \gamma): F \in \mathcal{K}_{c}, \gamma>0,\|T(z, w ; s)\|_{\infty}<\gamma\right\}
$$

Problem 6.3 (State feedback problem in $R H_{\infty}$ ) Given a scalar $\gamma>0$, determine the conditions for the existence of a state feedback matrix $F$ such that the pair $(F, \gamma)$ is strictly feasible, that is $(F, \gamma) \in$ int $\mathcal{K}_{\gamma c}$.

The Full information problem as stated and solved in Theorem 5.2 is much more general than its State feedback version. However, the last captures the essential features of the former in the sense that the existence of a solution to the Full information problem, under Assumptions 5.1-5.2, is assured provided there exists a symmetric, positive semidefinite and stabilizing solution $P_{\infty}$ of the Riccati equation (in the unknown $P$ )

$$
\begin{equation*}
0=P A_{c}+A_{c}^{\prime} P-P\left(B_{2} B_{2}^{\prime}-\gamma^{-2} B_{1} B_{1}^{\prime}\right) P+C_{1 c}^{\prime} C_{1 c} \tag{6.56}
\end{equation*}
$$

As in the proof of Theorem 5.2 let us proceed by assuming that i) The pair [( $A-$ $\left.\left.B_{2} D_{12}^{\prime} C_{1}\right),\left(I-D_{12} D_{12}^{\prime}\right) C_{1}\right]$ is observable and the pair $\left(A, B_{2}\right)$ is stabilizable and ii) $D_{12}^{\prime} D_{12}=I$.

Theorem 6.6 (State feedback) Assume matrix $B_{1} B_{1}^{\prime}$ is positive definite and let $\gamma$ a positive scalar be given. There exists a strictly feasible pair $(F, \gamma) \in \operatorname{int} \mathcal{K}_{\gamma c}$ if and only if there exists $W$ such that $\left(W, \gamma^{2}\right) \in \operatorname{int} \mathcal{C}_{\gamma c}$, where

$$
\text { int } \mathcal{C}_{\gamma c}:=\left\{(W, \mu): W \geq 0, \quad \mu>0, v^{\prime} \Theta_{\gamma c}(W, \mu) v<0, \forall v \in \mathcal{N}_{c}\right\}
$$

Proof Under the assumptions made, the State feedback problem is solvable if and only if there exists $P=P_{\infty}>0$ solution of the Riccati equation (6.56). However, the existence of a stabilizing state feedback control such that the transfer function norm $\|T(z, w ; s)\|_{\infty}$ is strictly less than $\gamma$ implies that there exists $0<\bar{\gamma}<\gamma$ for which the State feedback problem is also solvable. Consequently (recall Theorem 5.2) we can say that there exists a strictly feasible pair $(F, \gamma)$ if and only if there exists a positive definite and stabilizing solution $\bar{P}_{\infty}$ of the Riccati equation (in the unknown $P$ )

$$
0=P A_{c}+A_{c}^{\prime} P-P\left(B_{2} B_{2}^{\prime}-\bar{\gamma}^{-2} B_{1} B_{1}^{\prime}\right) P+C_{1 c}^{\prime} C_{1 c}
$$

Hence, the necessity is proved if we are able to calculate a matrix $W$ such that $\left(W, \gamma^{2}\right) \in \operatorname{int} \mathcal{C}_{\gamma c}$. Such a matrix is

$$
W_{\infty}:=\gamma^{2}\left[\begin{array}{cc}
\bar{P}_{\infty}^{-1} & \bar{P}_{\infty}^{-1} \bar{F}_{\infty}^{\prime} \\
\bar{F}_{\infty} \bar{P}_{\infty}^{-1} & \bar{F}_{\infty} \bar{P}_{\infty}^{-1} \bar{F}_{\infty}^{\prime}
\end{array}\right] \geq 0
$$

where $\bar{F}_{\infty}=-B_{2}^{\prime} \bar{P}_{\infty}-D_{12}^{\prime} C_{1}$. Actually, it suffices to observe that, from the previous Riccati equation, the inequality

$$
\begin{aligned}
v^{\prime} \Theta_{\gamma c}\left(W_{\infty}, \gamma^{2}\right) v= & \left(\gamma \bar{P}_{\infty}^{-1} x\right)^{\prime}\left[A_{c}^{\prime} \bar{P}_{\infty}+\bar{P}_{\infty} A_{c}-\right. \\
& \left.-\bar{P}_{\infty}\left(B_{2} B_{2}^{\prime}-\gamma^{-2} B_{1} B_{1}^{\prime}\right) \bar{P}_{\infty}+C_{1 c}^{\prime} C_{1 c}\right]\left(\gamma \bar{P}_{\infty}^{-1} x\right) \\
== & {\left[1-(\gamma / \bar{\gamma})^{2}\right]\left\|B_{1}^{\prime} x\right\|^{2}<0 }
\end{aligned}
$$

holds for all $v \in \mathcal{N}_{c}$. The converse follows readily from Theorem 2.15.
It is interesting to analyze the optimal solution of problem (6.55) on the light of this result. Suppose the global optimum of (6.55) has been calculated yielding $\mu=\mu^{o}$. Then, $F=W_{2}^{\prime} W_{1}^{-1}$ and $\gamma^{o}=\sqrt{\mu^{o}}$ are such that $\left(F, \gamma^{o}\right) \in \mathcal{K}_{\gamma c}$. On the other hand, Theorem 6.6 assures that the Riccati equation (6.56) is solvable for any $\gamma>\gamma^{o}$ providing thus $F_{\infty}=-B_{2}^{\prime} P_{\infty}-D_{12}^{\prime} C_{1}$ which is also strictly feasible, that is $\left(F_{\infty}, \gamma\right) \in \operatorname{int} \mathcal{K}_{\gamma c}$.

Remark 6.17 The discussion made in Remark 6.6 concerning the assumption on matrix $B_{1}$, namely, $B_{1} B_{1}^{\prime}>0$ is still valid in the present case. Indeed, assume $B_{1} B_{1}^{\prime} \geq 0$ and define $\bar{B}_{1}:=\left\{B_{1} \sqrt{\epsilon} I\right\}$ with $\epsilon>0$ arbitrarily small. Clearly we now have $\bar{B}_{1} \bar{B}_{1}^{\prime}>0$ and for all $F \in \mathcal{K}_{c}$, it is true that the transfer function $\bar{T}(z, w ; s):=C_{c \infty}\left(s I-A_{c \infty}\right)^{-1} \bar{B}_{1}$ can be factorized as

$$
\bar{T}(z, w ;-j \omega) \bar{T}^{\prime}(z, w ; j \omega)=T(z, w ;-j \omega) T^{\prime}(z, w ; j \omega)+\epsilon G(-j \omega) G^{\prime}(j \omega), \quad \forall \omega \in R
$$

where $G(s):=C_{c \infty}\left(s I-A_{c \infty}\right)^{-1}$. The consequence is that

$$
\|T(z, w ; s)\|_{\infty} \leq\|\bar{T}(z, w ; s)\|_{\infty}
$$

and also

$$
\begin{aligned}
\|\bar{T}(z, w ; s)\|_{\infty}^{2} & =\sup _{\omega}\left\|\bar{T}^{\prime}(z, w ; j \omega)\right\|^{2} \\
& \leq \sup _{\omega}\left\|T^{\prime}(z, w ; j \omega)\right\|^{2}+\sup _{\omega}\|G(j \omega)\|^{2} \\
& \leq\|T(z, w ; s)\|_{\infty}^{2}+\epsilon\left\|C_{c \infty}\left(s I-A_{c \infty}\right)^{-1}\right\|_{\infty}^{2}
\end{aligned}
$$

which means that $\|\bar{T}(z, w ; s)\|_{\infty} \leq \gamma$ implies $\|T(z, w ; s)\|_{\infty} \leq \gamma$ and both norms differs one from the other by an amount of order $\epsilon$.

Remark 6.18 Based on the result of the previous theorem, it is readily seen that if we want to involve the set int $\mathcal{C}_{\gamma c}$ in an optimization problem them it is possible to work with the approximation

$$
\operatorname{int} \mathcal{C}_{\gamma c}=\left\{(W, \mu): W \geq 0, \mu>0, v^{\prime} \Theta_{\gamma c}(W, \mu) v \leq-\epsilon, \forall v \in \mathcal{N}_{c}\right\}
$$

where the scalar $\epsilon>0$ must be taken sufficiently small. Notice that this is exactly equivalent to replace matrix $B_{1}$ by matrix $\left\{B_{1} \sqrt{\epsilon} I\right\}$. The obvious advantage is that in doing this the above approximation is always a closed convex set.

Remark 6.19 Theorem 5.2 makes possible the existence of a state feedback controller $u=F_{\infty} x$ called central controller which solves the Full information problem under the assumptions i) The pair $\left[\left(A-B_{2} D_{12}^{\prime} C_{1}\right),\left(I-D_{12} D_{12}^{\prime}\right) C_{1}\right]$ is detectable and the pair $\left(A, B_{2}\right)$ is stabilizable and ii) $D_{12}^{\prime} D_{12}=I$. Using mathematical programming arguments only, it can be obtained as follows. Given $\gamma>0$, determine $F \in \mathcal{K}_{c}$ and $P=P^{\prime} \geq 0$ such that

$$
\min \left\{\operatorname{trace}\left[B_{1}^{\prime} P B_{1}\right]: 0=A_{c \infty}^{\prime} P+P A_{c \infty}+\gamma^{-2} P B_{1} B_{1}^{\prime} P+C_{c \infty}^{\prime} C_{c \infty}\right\}
$$

Writing the associated Lagrangian (recall Remark 6.7) with $\Lambda=\Lambda^{\prime}$ being the matrix of Lagrange multipliers associated to the equality constraint, the necessary conditions for optimality are

$$
\begin{aligned}
& 0=A_{c \infty}^{\prime} P+P A_{c \infty}+\gamma^{-2} P B_{1} B_{1}^{\prime} P+C_{c \infty}^{\prime} C_{c \infty} \\
& 0=\left(A_{c \infty}+\gamma^{-2} B_{1} B_{1}^{\prime} P\right) \Lambda+\Lambda\left(A_{c \infty}+\gamma^{-2} B_{1} B_{1}^{\prime} P\right)^{\prime}+B_{1} B_{1}^{\prime} \\
& 0=\left(F+D_{12}^{\prime} C_{1}+B_{2}^{\prime} P\right) \Lambda
\end{aligned}
$$

Restricting ourselves to the solutions such that matrix $A_{c \infty}+\gamma^{-2} B_{1} B_{1}^{\prime} P$ is stable then with $B_{1} B_{1}^{\prime}>0$ we necessarily have $\Lambda>0$. Solving the last equation for $F$, the first one yields $F=F_{\infty}=-B_{2}^{\prime} P_{\infty}-D_{12}^{\prime} C_{1}$ where, under the previous assumptions, $P_{\infty}$ is the symmetric, positive semidefinite and stabilizing solution of the Riccati equation (in the unknown $P$ )

$$
0=A_{c}^{\prime} P+P A_{c}-P\left(B_{2} B_{2}^{\prime}-\gamma^{-2} B_{1} B_{1}^{\prime}\right) P+C_{1 c}^{\prime} C_{1 c}
$$

Moreover, since $A_{c c}=A+B_{2} F+\gamma^{-2} B_{1} B_{1}^{\prime} P$ calculated for $F=F_{\infty}$ and $P=P_{\infty}$ is stable then $B_{1} B_{1}^{\prime}>0$ assures that

$$
\Lambda_{\infty}:=\int_{0}^{\infty} e^{A_{c c} t} B_{1} B_{1}^{\prime} e^{A_{c c}^{\prime} t} d t>0
$$

and consequently ( $F_{\infty}, P_{\infty}, \Lambda_{\infty}$ ) is the unique solution of the necessary conditions for optimality. It remains to verify that it is actually the global solution of the proposed problem . To show this, let $F \in \mathcal{K}_{\gamma c}$ be arbitrary which is the same to say that only the first necessary condition is satisfied for some positive semidefinite and stabilizing matrix $P$. Some simple algebraic manipulations gives

$$
\begin{aligned}
0= & A_{c c}^{\prime}\left(P-P_{\infty}\right)+\left(P-P_{\infty}\right) A_{c c}- \\
& -\gamma^{-2}\left(P-P_{\infty}\right) B_{1} B_{1}^{\prime}\left(P-P_{\infty}\right)+\left(F-F_{\infty}\right)^{\prime}\left(F-F_{\infty}\right)
\end{aligned}
$$

The stability of matrix $A_{c c}$ together with the result of Lemma C. 4 imply that there exists $P-P_{\infty} \geq 0$ solution to this Riccati equation so that the objective function attains the global minimum at $P=P_{\infty}$. The problem introduced here deserves an additional interpretation to be used in the analysis of mixed $R H_{2} / R H_{\infty}$ optimal control problems. For any $F \in \mathcal{K}_{\gamma c}$ it is simple to verify that

$$
\begin{aligned}
\operatorname{trace}\left[B_{1}^{\prime} P B_{1}\right] & =\operatorname{trace}\left[\int_{0}^{\infty} B_{1}^{\prime} e^{A_{c \infty}^{\prime} t}\left(\gamma^{-2} P B_{1} B_{1}^{\prime} P+C_{c \infty}^{\prime} C_{c \infty}\right) e^{A_{c \infty} t} B_{1} d t\right] \\
& \geq \operatorname{trace}\left[\int_{0}^{\infty} B_{1}^{\prime} e^{A_{c \infty}^{\prime} t} C_{c \infty}^{\prime} C_{c \infty} e^{A_{c \infty} t} B_{1} d t\right] \\
& \geq\|T(z, w ; s)\|_{2}^{2}
\end{aligned}
$$

which enables us to say that the central controller actually minimizes an upper bound (and so provides a sub-optimal solution) to the problem

$$
\inf \left\{\|T(z, w ; s)\|_{2}^{2}: F \in \mathcal{K}_{c},\|T(z, w ; s)\|_{\infty}<\gamma\right\}
$$

which is a mixed $\mathrm{RH}_{2} / R H_{\infty}$ optimal control design problem to be deeply studied in the forthcoming section. For the moment notice that when $\gamma \rightarrow \infty$ the mixed problem tends to the State feedback problem in $\mathrm{RH}_{2}$. The necessary optimality conditions reduce to those of Remark 6.7.

Remark 6.20 The central controller can also be used to determine an approximate solution to the problem

$$
\gamma_{o p t}=\inf \left\{\gamma:(F, \gamma) \in \operatorname{int} \mathcal{K}_{\gamma c}\right\}
$$

The basic algorithm can be stated as follows. For a given $\gamma_{k}>0$ suppose that $P_{\infty} \geq 0$ is stabilizing. To make clear the dependence on $\gamma_{k}$ of the transfer function $T(z, w ; s)$ when the central controller is used, it is denoted as $T_{k}(z, w ; s)$.

1) Choose $\gamma_{0}>\gamma_{o p t}$ (possibly $\gamma_{0}=\infty$ ).
2) Iterate until convergence $\gamma_{k+1}=\left\|T_{k}(z, w ; s)\right\|_{\infty}$.

It is immediate to see that this procedure approaches to $\gamma_{o p t}$ as $k$ goes to infinite since $\gamma_{k+1}=\left\|T_{k}(z, w ; s)\right\|_{\infty}<\gamma_{k}$, however the speed of convergence may be very slow when $\gamma_{k}$ becomes close to $\gamma_{o p t}$. This feature is numerically illustrated with the system of Example 5.2 for $\Omega=0$. The system is

$$
\begin{aligned}
\dot{x} & =A x+B_{1} w_{1}+B_{2} u \\
z & =C_{1} x+u \\
y & =x
\end{aligned}
$$

where

$$
A=\left[\begin{array}{rr}
0 & 1 \\
-1 & -1
\end{array}\right], \quad B_{1}=\left[\begin{array}{l}
1 \\
0
\end{array}\right], \quad B_{2}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad C_{1}=\left[\begin{array}{ll}
0 & -2
\end{array}\right]
$$

It has been calculated that $1.61<\gamma_{o p t}<1.62$ while the first six iterations of the previous algorithm furnish

$$
\left\{\gamma_{k}\right\}_{k=0}^{6}=\{\infty, 2.30,2.00,1.88,1.82,1.78,1.76\}
$$

The importance of problem (6.55) is now apparent. It is jointly convex on the variables ( $W, \mu$ ) and so can be solved more efficiently.

Remark 6.21 It is numerically attractive to get a priori bounds such that $\gamma_{\min } \leq \gamma_{o p t} \leq$ $\gamma_{\max }$. Indeed, a possible value to $\gamma_{\max }$ is simply determined from any $\left(F, \gamma_{\max }\right) \in \mathcal{K}_{\gamma c}$, as for instance $F=F_{2}$ yields

$$
\gamma_{\max }:=\|T(z, w ; s)\|_{\infty}
$$

The determination of $\gamma_{\text {min }}$ is much more involved. It follows from the factorization

$$
\begin{aligned}
M_{c} W+W M_{c}^{\prime}+\mu^{-1} W R_{c \epsilon} W+Q_{c}= & Q_{c}-\mu M_{c} R_{c \epsilon}^{-1} M_{c}^{\prime}+ \\
& +\mu^{-1}\left(\mu M_{c} R_{c \epsilon}^{-1}+W\right) R_{c \epsilon}\left(\mu M_{c} R_{c \epsilon}^{-1}+W\right)^{\prime} \\
\geq & Q_{c}-\mu M_{c} R_{c \epsilon}^{-1} M_{c}^{\prime}
\end{aligned}
$$

where $R_{c \epsilon}:=R_{c}+\epsilon I$ and $\epsilon>0$. The left hand side of this inequality goes to $\Theta_{\gamma c}(W, \mu)$ as $\epsilon$ goes to zero then if we set

$$
\gamma_{m}^{2}:=\lim _{\epsilon \rightarrow 0} \min \left\{\mu: v^{\prime}\left(Q_{c}-\mu M_{c} R_{c \epsilon}^{-1} M_{c}^{\prime}\right) v \leq 0, \quad \forall v \in \mathcal{N}_{c}\right\}
$$

for any $\mu<\gamma_{m}^{2}$ there exists a vector $\bar{v} \in \mathcal{N}_{c}$ such that $\bar{v}^{\prime}\left(Q_{c}-\mu M_{c} R_{c \varepsilon}^{-1} M_{c}^{\prime}\right) \bar{v}>0$ and so $\bar{v}^{\prime} \Theta_{\gamma c}(W, \mu) \bar{v}>0$ enabling the final conclusion that $(W, \mu) \notin \mathcal{C}_{\gamma c}$ which from Theorem 6.5 is the same to say that there is no $F$ matrix such that the pair $(F, \gamma) \in \mathcal{K}_{\gamma c}$ for all $\gamma<\gamma_{m}$. Unfortunately, the calculation involved to get $\gamma_{m}$ is prohibitive. However, assuming $B_{1}$ of full column rank, restricting $v=\left[x^{\prime} 0\right]^{\prime} \in \mathcal{N}_{c}$ to $x=B_{1}\left(B_{1}^{\prime} B_{1}\right)^{-1} \xi$ and defining the matrix $V:=\left(B_{1}^{\prime} B_{1}\right)^{-1} B_{1}^{\prime}\left[A B_{2}\right]$ we get

$$
\begin{aligned}
\gamma_{m}^{2} & \geq \lim _{\epsilon \rightarrow 0} \min \left\{\mu: \xi^{\prime}\left(I-\mu V R_{c \epsilon}^{-1} V^{\prime}\right) \xi \leq 0, \quad \forall\|\xi\|=1\right\} \\
& \geq \lim _{\epsilon \rightarrow 0} \min \left\{\mu: \mu^{-1} \leq \xi^{\prime} V R_{c \epsilon}^{-1} V^{\prime} \xi, \quad \forall\|\xi\|=1\right\} \\
& \geq \lim _{\epsilon \rightarrow 0} \lambda_{\min }^{-1}\left(V R_{c \epsilon}^{-1} V^{\prime}\right):=\gamma_{\min }^{2}
\end{aligned}
$$

This limit can be easily calculated numerically and is a valid lower bound to $\gamma_{o p t}$ since as we have shown $\mu=\gamma_{m i n}^{2}<\gamma_{m}^{2}$ is always infeasible.

Remark 6.22 The numerical solution of problem (6.55) by means of convex programming algorithms depends on the boundedness of the feasible set $\mathcal{C}_{\gamma c}$. Following the same lines of Remark 6.8 , we have to analyze the conditions under which, for some $\mu>0$, there exists $(W, \mu) \in \mathcal{C}_{\gamma c}$ and $\tilde{W} \neq 0$ such that $W+\lambda \tilde{W} \in \mathcal{C}_{\gamma c}$ with $\lambda>0$ arbitrarily large. Unfortunately, the only constraint involving the sub-block $W_{3}$ imposed by $W \in \mathcal{C}_{\gamma c}$ is $W_{3} \geq W_{2}^{\prime} W_{1}^{-1} W_{2}$ and so the above conditions are always satisfied for the matrix (other possibilities exist)

$$
\tilde{W}=\left[\begin{array}{ll}
0 & 0 \\
0 & I
\end{array}\right]
$$

consequently, the set $\mathcal{C}_{\gamma c}$ is not bounded even though $\mu<\infty$. In order to circumvent this difficulty we first notice that

$$
\mathcal{C}_{\gamma c} \subseteq\{\mu: \mu>0\} \times \mathcal{C}_{c}
$$

where with no loss of generality (recall Remark 6.21) the unbounded set $\mu>0$ can be replaced by the bounded one $\gamma_{\text {min }}^{2} \leq \mu \leq \gamma_{\text {max }}^{2}$. Second, we already know that the set $\mathcal{C}_{c}$ is not bounded but the set

$$
\mathcal{C}_{c \beta}:=\mathcal{C}_{c} \cap\left\{W: \operatorname{trace}\left[R_{c} W\right] \leq \beta\right\}
$$

is bounded for all finite $\beta$ (recall Remark 6.8) provided the pair ( $-A_{c}, C_{1 c}$ ) is detectable. Under this assumption, the conclusion is that we can always solve an approximate version of problem (6.55), namely

$$
\min \left\{\mu:(W, \mu) \in \mathcal{C}_{\gamma c}, \quad \operatorname{trace}\left[R_{c} W\right] \leq \beta\right\}
$$

which has a bounded feasible domain and the optimal solution is arbitrarily close to the optimal solution of problem (6.55) provided the parameter $\beta$ is chosen appropriately large. This problem has another interesting and important interpretation for robust control design to be discussed in the next section.

Remark 6.23 (Linear matrix inequalities - LMI) As in the $R H_{2}$ design (recall Remark 6.9), the State feedback problem in $R H_{\infty}$ can also be stated only in terms of Linear matrix inequalities. Assuming the pair $\left(A_{c}, C_{1 c}\right)$ is observable, then any $P>0$ satisfying the Riccati inequality

$$
0 \geq A_{c}^{\prime} P+P A_{c}-P\left(B_{2} B_{2}^{\prime}-\gamma^{-2} B_{1} B_{1}^{\prime}\right) P+C_{1 c}^{\prime} C_{1 c}
$$

assures that with $F=-B_{2}^{\prime} P-D_{12}^{\prime} C_{1}$, the closed-loop transfer function is such that $\|T(z, w ; s)\|_{\infty} \leq \gamma$. Defining $X=P^{-1}$ and applying the Schur complement formula, the above inequality turns out to be equivalent to the linear matrix inequality

$$
\mathcal{A}(X)=\left[\begin{array}{cc}
A_{c} X+X A_{c}^{\prime}-B_{2} B_{2}^{\prime}+\gamma^{-2} B_{1} B_{1}^{\prime} & X C_{1 c}^{\prime} \\
C_{1 c} X & -I
\end{array}\right] \leq 0
$$

On the other hand, from Remark 6.19 it is clear that $X=P_{\infty}^{-1}$ is the global solution of

$$
\min \left\{\operatorname{trace}\left[B_{1}^{\prime} X^{-1} B_{1}\right]: X>0, \mathcal{A}(X) \leq 0\right\}
$$

which is a convex programming problem. Furthermore, it reduces to that of Remark 6.9 as $\gamma$ goes to infinity.

We consider now the Partial information problem. Let the linear system be defined by the standard state space representation

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{6.57}\\
z & =C_{1} x+D_{12} u  \tag{6.58}\\
y & =C_{2} x+D_{21} w \tag{6.59}
\end{align*}
$$

where the output variable $y$ is available for feedback. The situation is illustrated in fig. 6.5 where $K(s)$ is the controller transfer function. We call the controller $K(s)$ strictly feasible if it assures that $T(z, w ; s)$ belongs to $R H_{\infty}$ and $\|T(z, w ; s)\|_{\infty}<\gamma$. The so called central controller is then obtained from part b) of Theorem 5.4 with $Q(s)=0$. Notice that the central controller does not exhibits the classical observerbased structure. Indeed, it is given by

$$
K(s):=\left[\begin{array}{c|c}
A+B_{2} F_{\infty}+\gamma^{-2} B_{1} B_{1}^{\prime} P_{\infty}+Z_{\infty} L_{\infty} C_{2 \infty} & -Z_{\infty} L_{\infty}  \tag{6.60}\\
\hline F_{\infty} & 0
\end{array}\right]
$$

where $C_{2 \infty}:=C_{2}+\gamma^{-2} D_{21} B_{1}^{\prime} P_{\infty}$. However, the existence of a solution to the Partial information problem depends only upon the existence of a symmetric, positive semidefinite and stabilizing matrix $P_{\infty}$ together with the existence of a symmetric, positive semidefinite and stabilizing matrix $\Pi_{\infty}$ of the Riccati equation (in the unknown II)

$$
\begin{equation*}
0=\Pi A_{f}^{\prime}+A_{f} \Pi-\Pi\left(C_{2}^{\prime} C_{2}-\gamma^{-2} C_{1}^{\prime} C_{1}\right) \Pi+B_{1 f} B_{1 f}^{\prime} \tag{6.61}
\end{equation*}
$$

satisfying the additional constraint $r_{s}\left(P_{\infty} \Pi_{\infty}\right)<\gamma^{2}$. Based on this, the following problem is stated and solved.

Problem 6.4 (Output feedback problem in $R H_{\infty}$ ) Given a scalar $\gamma>0$, determine the conditions for the existence of a strictly feasible output feedback controller $K(s)$.

Recall Theorem 5.4 where the complete solution of the Partial information problem is provided. Here, we consider the following assumptions i) The pair $[A-$ $\left.\left.B_{2} D_{12}^{\prime} C_{1}\right),\left(I-D_{12} D_{12}^{\prime}\right) C_{1}\right]$ is observable and the pair $\left[\left(A-B_{1} D_{21}^{\prime} C_{2}\right), B_{1}\left(I-D_{21}^{\prime} D_{21}\right)\right]$ is reachable, ii) The pair $\left(A, B_{2}\right)$ is stabilizable and the pair $\left(A, C_{2}\right)$ is detectable and iii) $D_{12}^{\prime} D_{12}=I$ and $D_{21} D_{21}^{\prime}=I$. The first assumption assures that matrices $P_{\infty}$ and $\Pi_{\infty}$ are both symmetric, positive definite and stabilizing. Furthermore, it is also assumed that $B_{1} B_{1}^{\prime}$ and $C_{1}^{\prime} C_{1}$ are positive definite matrices.

Before proceed we need to define the dual version of the convex set $\mathcal{C}_{\gamma c}$, it is

$$
\begin{equation*}
\mathcal{C}_{\gamma f}:=\left\{(V, \mu): V \geq 0, \mu>0, \quad v^{\prime} \Theta_{\gamma f}(V, \mu) v \leq 0, \forall v \in \mathcal{N}_{f}\right\} \tag{6.62}
\end{equation*}
$$

where $\Theta_{\gamma f}(V, \mu):=\Theta_{f}(V)+\mu^{-1} V R_{f} V$. Clearly this set is convex as well (recall the proof of Theorem 6.5).

Theorem 6.7 (Output feedback) Consider the previous assumptions and let $\gamma$ a positive scalar be given. There exists a strictly feasible controller with transfer function $K(s)$ if and only if there exist $W$ such that $\left(W, \gamma^{2}\right) \in$ int $\mathcal{C}_{\gamma c}, V$ such that $\left(V, \gamma^{2}\right) \in$ int $\mathcal{C}_{\gamma f}$ satisfying the convex coupling constraint $(W, V, \gamma) \in$ int $\mathcal{Z}_{c f}$, where

$$
\mathcal{Z}_{c f}:=\left\{(W, V, \gamma):\left[\begin{array}{cc}
W_{1} & \gamma I  \tag{6.63}\\
\gamma I & V_{1}
\end{array}\right] \geq 0\right\}
$$

Proof Under the assumptions made, the Output feedback problem is solvable if and only if $\bar{P}_{\infty}>0$ solves the Riccati equation ( 6.56 ), $\bar{\Pi}_{\infty}>0$ solves the Riccati equation (6.61), both with $\gamma$ replaced by a suitable $0<\bar{\gamma}<\gamma$ and $r_{s}\left(\bar{P}_{\infty} \bar{\Pi}_{\infty}\right)<\bar{\gamma}^{2}$. We have already shown (recall Theorem 6.6) that $\left(W_{\infty}, \gamma^{2}\right) \in \mathcal{C}_{\gamma c}$. On the other hand, it can be verified that matrix

$$
V_{\infty}:=\gamma^{2}\left[\begin{array}{cc}
\bar{\Pi}_{\infty}^{-1} & \bar{\Pi}_{\infty}^{-1} \bar{L}_{\infty} \\
\bar{L}_{\infty}^{\prime} \bar{\Pi}_{\infty}^{-1} & \bar{L}_{\infty}^{\prime} \bar{\Pi}_{\infty}^{-1} \bar{L}_{\infty}
\end{array}\right] \geq 0
$$

with $\bar{L}_{\infty}=-\bar{\Pi}_{\infty} C_{2}^{\prime}-B_{1} D_{21}^{\prime}$ is such that

$$
\begin{aligned}
v^{\prime} \Theta_{\gamma f}\left(V_{\infty}, \gamma^{2}\right) v= & \left(\gamma \bar{\Pi}_{\infty}^{-1} x\right)^{\prime}\left[A_{f} \bar{\Pi}_{\infty}+\bar{\Pi}_{\infty} A_{f}^{\prime}-\right. \\
& \left.-\bar{\Pi}_{\infty}\left(C_{2}^{\prime} C_{2}-\gamma^{-2} C_{1}^{\prime} C_{1}\right) \bar{\Pi}_{\infty}+B_{1 f} B_{1 f}^{\prime}\right]\left(\gamma \bar{\Pi}_{\infty}^{-1} x\right) \\
= & {\left[1-(\gamma / \bar{\gamma})^{2}\right]\left\|C_{1} x\right\|^{2}<0 }
\end{aligned}
$$

holds for all $v \in \mathcal{N}_{f}$, proving that $\left(V_{\infty}, \gamma^{2}\right) \in \operatorname{int} \mathcal{C}_{\gamma f}$. Finally, using the Schur complement formula it is also verified that the coupling constraint is satisfied for
$W_{1}=\gamma^{2} \bar{P}_{\infty}^{-1}$ and $V_{1}=\gamma^{2} \bar{\Pi}_{\infty}^{-1}$ since $r_{s}\left(\bar{P}_{\infty} \bar{\Pi}_{\infty}\right)<\bar{\gamma}^{2}<\gamma^{2}$. The converse is immediate from the previous result of Theorem 6.6.

We have interest to reduce $\gamma$ to the minimum value keeping feasibility of all constraints introduced in Theorem 6.7. The main goal is to generate a convex problem, similar to (6.55), valid in the output feedback case. Unfortunately, the change of variable defined previously, namely $\mu=\gamma^{2}$ destroys the convexity of the set $\mathcal{Z}_{c f}$. The other possibility comes to light from the simple observation that $\left(\gamma^{-1} W, \gamma^{2}\right) \in \mathcal{C}_{\gamma c}$ defines a convex set with respect to the new pair of variables $(\hat{W}, \gamma)$ where $\hat{W}:=\gamma^{-1} W$. The same occurs for the set $\mathcal{C}_{\gamma f}$, that is $\left(\gamma^{-1} V, \gamma^{2}\right) \in \mathcal{C}_{\gamma f}$ defines a convex set with respect to the new pair of variables $(\hat{V}, \gamma)$ where $\hat{V}:=\gamma^{-1} V$. Furthermore, the set $\mathcal{Z}_{c f}$, rewritten as

$$
\mathcal{Z}_{c f}:=\left\{(\hat{W}, \hat{V}):\left[\begin{array}{cc}
\hat{W}_{1} & I \\
I & \hat{V}_{1}
\end{array}\right] \geq 0\right\}
$$

makes clear that convexity is once again preserved. Based on these facts we are able to conclude that Theorem 6.7 can equivalently be stated in terms of the existence of a triple $(\hat{W}, \hat{V}, \gamma) \in \operatorname{int} \mathcal{O}_{c f}$ with $\mathcal{O}_{c f}$ being a convex set. Naturally, the output feedback counterpart of (6.55) is the convex programming problem

$$
\begin{equation*}
\min \left\{\gamma:(\hat{W}, \hat{V}, \gamma) \in \mathcal{O}_{c f}\right\} \tag{6.64}
\end{equation*}
$$

which provides the minimum value $\gamma^{o}$ keeping the Output feedback problem solvable for any $\gamma>\gamma^{\circ}$. In other words, for $\gamma>\gamma^{o}$ the associated Riccati equations admit stabilizing solutions $P_{\infty}$ and $\Pi_{\infty}$ respectively. These solutions may be used to calculate from (6.60), the (central) feasible controller.

Remark 6.24 Remark 6.17 applies to the set $\mathcal{C}_{\gamma f}$ also. If matrix $C_{1}^{\prime} C_{1}$ is not strictly positive definite then $C_{1}^{\prime}$ must be replaced by $\left[C_{1}^{\prime} \sqrt{\epsilon} I\right]$ where $\epsilon>0$ is arbitrarily small. This is done with no loss of generality and is important to have $V_{1}>0$ for all $V \in \mathcal{C}_{\gamma f}$.

Remark 6.25 The set $C_{\gamma f}$ is not bounded but fortunately, the same reasoning of Remarks 6.21 and 6.22 applies to it. First the set $\mu>0$ can be replaced by the bounded one $\gamma_{\text {min }}^{2} \leq \mu \leq \gamma_{\text {max }}^{2}$ and second, the convex set

$$
\mathcal{C}_{\gamma f} \cap\left\{V: \operatorname{trace}\left[R_{f} V\right] \leq \beta\right\}
$$

has a bounded domain for all finite $\beta$ provided the pair $\left(-A_{f}, B_{1 f}\right)$ is stabilizable.
The results of this section put in evidence the potentialities of the convex programming approach to deal with optimal control problems formulated in $R H_{2}$ and $R H_{\infty}$ spaces. The same manipulations will be used in the next sections to handle more involved problems. To ease the presentation, the technical assumptions involving the positive definiteness of, for instance, matrices $B_{1} B_{1}^{\prime}$ and $C_{1}^{\prime} C_{1}$ are assumed throughout. Actually, as discussed before they can be enforced with no loss of generality.

### 6.3 Mixed $\mathrm{RH}_{2} / R H_{\infty}$ control

The main goal of the so called mixed $R \mathrm{H}_{2} / R \mathrm{H}_{\infty}$ control problem is to take into account the two major features of any control system design. First it is desirable to optimize performance. Second it is important to be aware that the model at hand


Figure 6.6: The mixed control design structure
always represents a nominal system while the true system is subject to uncertainties. In order to improve practical needs, these two aspects have to be accommodated in the same design problem. It consists on the optimization of the $R H_{2}$ norm of a transfer function while the $R H_{\infty}$ norm of another transfer function is constrained to be less than a certain prescribed level.

The situation is represented in fig. 6.6 where $\Sigma$ is a linear time-invariant system and $\Sigma_{c}$ is the controller. The mixed $R H_{2} / R H_{\infty}$ optimal control problem is to find (if one exists) a controller $\Sigma_{c}$ such that, for $\gamma>0$ given

$$
\begin{equation*}
\inf \left\{\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}:\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty}<\gamma\right\} \tag{6.65}
\end{equation*}
$$

A controller solving this problem imposes to the closed loop system an optimal performance against exogenous perturbation in the channel $w$ to $z_{0}$ while robust stability is guaranteed to all model uncertainty expressed as $z_{1}=\Delta w$ such that $\|\Delta\|_{\infty} \leq 1 / \gamma$. The solution of this problem can not be obtained from matrix Riccati manipulations. We show here the difficulties we have to face and the manipulations and approximations we have to introduce in order to solve it by means of convex programming methods. Before all, based on the results of the previous section, it is to be noticed that if instead of Problem (6.65) we solve

$$
\begin{equation*}
\min \left\{\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}:\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty} \leq \gamma\right\} \tag{6.66}
\end{equation*}
$$

then its optimal solution (if any) provides a suboptimal solution to Problem (6.65) with $\gamma$ replaced by $\bar{\gamma}>\gamma$ and the degree of suboptimality becomes arbitrarily small provided $\bar{\gamma}$ is chosen arbitrarily close to $\gamma$.

### 6.3.1 State feedback design

Consider $\gamma>0$ be a fixed scalar and let the systems $\Sigma$ be defined as follows

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{6.67}\\
z_{0} & =C_{0} x+D_{0} u  \tag{6.68}\\
z_{1} & =C_{1} x+D_{12} u  \tag{6.69}\\
y & =x \tag{6.70}
\end{align*}
$$

where the whole state vector is available for feedback. It is assumed that i) the pair $\left(A, B_{2}\right)$ is stabilizable and ii) $D_{12}^{\prime} D_{12}=I$. Furthermore, the controller $\Sigma_{c}$ is assumed to be given by

$$
\begin{equation*}
u=F x \tag{6.71}
\end{equation*}
$$

where the gain matrix $F$ is to be determined. To guarantee that the constraint $\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty} \leq \gamma$ is satisfied, we have to consider $F \in \mathcal{K}_{\gamma c}$ which is the same to
say that $F \in \mathcal{K}_{c}$ and there exists $P>0$ satisfying the Riccati inequality (recall that Theorem 2.16 applies since $B_{1} B_{1}^{\prime}>0$ ).

$$
\begin{equation*}
0 \geq A_{c \infty} P+P A_{c \infty}^{\prime}+\gamma^{-2} P C_{c \infty}^{\prime} C_{c \infty} P+B_{1} B_{1}^{\prime} \tag{6.72}
\end{equation*}
$$

with $A_{c \infty}:=A+B_{2} F$ and $C_{c \infty}:=C_{1}+D_{12} F$. On the other hand, for the same control gain $F \in \mathcal{K}_{c}$, the transfer function from the input $w$ to the output $z_{0}$ belongs to $\mathrm{RH}_{2}$ and

$$
\begin{equation*}
\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}=\operatorname{trace}\left[\left(C_{0}+D_{0} F\right) P_{r}\left(C_{0}+D_{0} F\right)^{\prime}\right] \tag{6.73}
\end{equation*}
$$

where $P_{r}>0$ solves the Lyapunov equation

$$
\begin{equation*}
0=A_{c \infty} P_{r}+P_{r} A_{c \infty}^{\prime}+B_{1} B_{1}^{\prime} \tag{6.74}
\end{equation*}
$$

Simple comparison of inequality (6.72) with equality (6.74) puts in evidence that $P \geq P_{r}$. Consequently, any feasible $P$ for inequality (6.72) provides an upper bound for the norm of the transfer function under consideration, that is

$$
\begin{equation*}
\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2} \leq \operatorname{trace}\left[\left(C_{0}+D_{0} F\right) P\left(C_{0}+D_{0} F\right)^{\prime}\right] \tag{6.75}
\end{equation*}
$$

In the optimal control problem (6.66), the quantity $\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}$ is then replaced by the quantity in right hand side of (6.75). It must be clear that this simplifies the problem to be dealt with at the expense of optimality. In fact, the minimum of the proposed upper bound may not coincide with the optimum of the true design problem (6.66). The next Theorem shows that this problem is convex which opens again the possibility of solving it by means of efficient numerical methods.

Theorem 6.8 Let $\gamma>0$ be given and define the symmetric and positive semidefinite extended matrix

$$
R_{0}:=\left[\begin{array}{c}
C_{0}^{\prime} \\
D_{0}^{\prime}
\end{array}\right]\left[\begin{array}{ll}
C_{0} & D_{0}
\end{array}\right]
$$

Let $\bar{W}$ be the optimal solution of the convex programming problem

$$
\begin{equation*}
J_{s u b}:=\min \left\{\operatorname{trace}\left[R_{0} W\right]:\left(W, \gamma^{2}\right) \in \mathcal{C}_{\gamma c}\right\} \tag{6.76}
\end{equation*}
$$

Then, $F=\bar{W}_{2}^{\prime} \bar{W}_{1}^{-1} \in \mathcal{K}_{\gamma c}$ minimizes an upper bound of the objective function of problem (6.66) in the sense that $\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2} \leq \operatorname{trace}\left[R_{0} W\right]$, for all $\left(W, \gamma^{2}\right) \in \mathcal{C}_{\gamma c}$.

Proof From Theorem 6.5 it follows that $\left(\bar{W}, \gamma^{2}\right) \in \mathcal{C}_{\gamma c}$ implies $F=\bar{W}_{2}^{\prime} \bar{W}_{1}^{-1} \in \mathcal{K}_{\gamma C}$. On the other hand, for any $\left(W, \gamma^{2}\right) \in \mathcal{C}_{\gamma c}$ we have $W_{1} \geq P_{r}$, consequently

$$
\begin{aligned}
\operatorname{trace}\left[R_{0} W\right]= & \operatorname{trace}\left[\left(C_{0}+D_{0} W_{2}^{\prime} W_{1}^{-1}\right) W_{1}\left(C_{0}+D_{0} W_{2}^{\prime} W_{1}^{-1}\right)^{\prime}\right]+ \\
& +\operatorname{trace}\left[D_{0}\left(W_{3}-W_{2}^{\prime} W_{1}^{-1} W_{2}\right) D_{0}^{\prime}\right] \\
\geq & \operatorname{trace}\left[\left(C_{0}+D_{0} W_{2}^{\prime} W_{1}^{-1}\right) P_{r}\left(C_{0}+D_{0} W_{2}^{\prime} W_{1}^{-1}\right)^{\prime}\right] \\
\geq & \left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}
\end{aligned}
$$

and the proof is complete.
The joint convexity of the set $\mathcal{C}_{\gamma c}$ is crucial for the introduction of another mixed problem, called the inverse mixed $R H_{2} / R H_{\infty}$ design problem, it is

$$
\begin{equation*}
\min \left\{\mu: \quad \operatorname{trace}\left[R_{0} W\right] \leq \beta,(W, \mu) \in \mathcal{C}_{\gamma c}\right\} \tag{6.77}
\end{equation*}
$$

where $\beta$ is a fixed positive scalar. For a given $\beta>0$, the constraint $\left\|T\left(z_{0}, w ; s\right)\right\|_{2} \leq$ $\sqrt{\beta}$ is satisfied while the upper bound on the admissible perturbation $\|\Delta\|_{\infty}<1 / \sqrt{\mu}$ is maximized. Needless to say that as $\beta$ goes to infinity, its optimal solution goes to the solution of problem (6.55) already treated in the last section. Another problem of practical interest is

$$
\begin{equation*}
\min \left\{\mu+\theta \operatorname{trace}\left[R_{0} W\right],(W, \mu) \in \mathcal{C}_{\gamma c}\right\} \tag{6.78}
\end{equation*}
$$

where $\theta$ is a fixed positive scalar. In this case, the scalar $\theta>0$ can be interpreted as the dual variable associated to the inequality constraint of problem (6.77), that is $\theta$ is the tradeoff, to be fixed by the designer, between the norms of the two transfer functions of system (6.67)-(6.70). Finally, it is important to stress that (recall Remark 6.18) if matrix $B_{1}$ is replaced by $\left[B_{1} \sqrt{\epsilon} I\right]$ with $\epsilon>0$ arbitrarily small them the optimal solutions of all previous mixed problems are strictly feasible, that is they belong to int $\mathcal{C}_{\gamma c}$.

Example 6.2 To get some feelings on the previous results, let us consider the following numerical example. The system is defined as

$$
\begin{aligned}
\dot{x} & =A x+B_{1} w+B_{2} u \\
z_{0} & =C_{0} x+u \\
z_{1} & =C_{1} x+u \\
y & =x
\end{aligned}
$$

where

$$
A=\left[\begin{array}{rr}
0 & 1 \\
-1 & -1
\end{array}\right], \quad B_{1}=\left[\begin{array}{l}
1 \\
0
\end{array}\right], \quad B_{2}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad C_{1}=\left[\begin{array}{ll}
0 & -2
\end{array}\right]
$$

Two different situations have been considerated :

1) With $C_{0}=[10]$ the mixed optimization problem (6.76) has been solved for $1.65 \leq \gamma \leq$ 3.00. The optimal solution provided the controller M. For comparison purpose we have also determined the central controller C associated to the constraint $\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty}<$ $\gamma$. Fig. 6.7 shows the optimal upper bound $J_{s u b}$ as well as the quantity $\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}$ for each controller. It is useful to observe that, in this case, the controller M obtained by the mixed problem is always better than the central controller. This behavior is expected for large values of $\gamma$ since in this situation, the gap between $J_{s u b}$ and $\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}$ becomes arbitrarily small. On the contrary, for small values of $\gamma$ the aforementioned gap is important and it is possible to build examples for which the central controller performs better than the one provided by the optimal solution of the mixed problem.
2) With $C_{0}=C_{1}$ the mixed optimization problem (6.76) has been solved again for the $\gamma$ in the same interval. As predicted in Remark 6.19, the optimal mixed controller M coincides with the central controller. For comparison, the performances of these controllers together with the apper bound $J_{s u b}$ are depicted in fig. 6.8.
Based on this, we can say that in many instances, the mixed problem (6.76) provides a controller that performs well when compared to the central controller. This suggestion becomes a fact for moderate and large values of $\gamma$. The main point to be retained however is that the controller M is only suboptimal for the true mixed $R H_{2} / R H_{\infty}$ problem. Numerically speaking, one of the main attractive features of problem (6.76) is that it is convex and so easy to solve.

Remark 6.26 (Post-optimization procedure) The optimal solution of the mixed optimization problem (6.76) is, generally, a suboptimal solution to the true mixed problem


Figure 6.7: The performances of the controllers $M$ and $C$
(6.66) since for $F=F_{s u b}$ provided by (6.76) the constraint $\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty} \leq \gamma$ is not satisfied with equality. Hence we are now concerned to determine a new state feedback gain, still suboptimal, but with a smaller value of the objective function $\left\|T\left(z_{0}, w ; s\right)\right\|_{2}$. To ease the algebraic manipulations involved we assume the pair $\left(A, B_{2}\right)$ is stabilizable, the pair $\left(A-B_{2} D_{0}^{\prime} C_{0},\left(I-D_{0} D_{0}^{\prime}\right) C_{0}\right)$ is detectable and $D_{0}^{\prime} D_{0}=I$. For $F=F_{\text {sub }}$ we get $\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}=\operatorname{trace}\left[B_{1}^{\prime} P_{s u b} B_{1}\right]$ where $P_{s u b} \geq 0$ solves the Lyapunov equation

$$
\begin{equation*}
0=\left(A+B_{2} F_{s u b}\right)^{\prime} P_{s u b}+P_{s u b}\left(A+B_{2} F_{s u b}\right)+\left(C_{0}+D_{0} F_{s u b}\right)^{\prime}\left(C_{0}+D_{0} F_{s u b}\right) \tag{6.79}
\end{equation*}
$$

Defining the state feedback gain

$$
\begin{equation*}
F_{\alpha}:=(1-\alpha) F_{s u b}-\alpha\left(B_{2}^{\prime} P_{\alpha}+D_{0}^{\prime} C_{0}\right) \tag{6.80}
\end{equation*}
$$

where $\alpha$ is a scalar to be determined such that $A+B_{2} F_{\alpha}$ is stable and $P_{\alpha}$ solves the Lyapunov equation

$$
\begin{equation*}
0=\left(A+B_{2} F_{\alpha}\right)^{\prime} P_{\alpha}+P_{\alpha}\left(A+B_{2} F_{\alpha}\right)+\left(C_{0}+D_{0} F_{\alpha}\right)^{\prime}\left(C_{0}+D_{0} F_{\alpha}\right) \tag{6.81}
\end{equation*}
$$

the following conclusions can be drawn. First, it is clear that for $F=F_{\alpha}$ the associated cost function is written as $\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}=\operatorname{trace}\left[B_{1}^{\prime} P_{\alpha} B_{1}\right]$ and, the equality

$$
0=\left(A+B_{2} F_{s u b}\right)^{\prime}\left(P_{s u b}-P_{\alpha}\right)+\left(P_{s u b}-P_{\alpha}\right)\left(A+B_{2} F_{s u b}\right)+\alpha(2-\alpha) \bar{F}_{s u b}^{\prime} \bar{F}_{s u b}
$$

valid for $\bar{F}_{\text {sub }}:=F_{\text {sub }}+B_{2} P_{\alpha}+D_{0}^{\prime} C_{0}$ guarantees an improvement in the $R H_{2}$ norm provided $\alpha \in[0,2]$ since under this condition $P_{\text {sub }} \geq P_{\alpha}$. Second, from (6.80) and (6.81) simple algebraic manipulations put in evidence that matrix $P_{\alpha}$ satisfies the Riccati equation (in the unknown $P$ )

$$
\begin{equation*}
0=A_{\alpha}^{\prime} P+P A_{\alpha}-P B_{2 \alpha} B_{2 \alpha}^{\prime} P+C_{0 \alpha}^{\prime} C_{0 \alpha} \tag{6.82}
\end{equation*}
$$

where

$$
\begin{aligned}
A_{\alpha} & :=A+B_{2}\left[(1-\alpha)^{2} F_{\text {sub }}+\alpha(\alpha-2) D_{0}^{\prime} C_{0}\right] \\
B_{2 \alpha} & :=B_{2} \sqrt{2 \alpha-\alpha^{2}} \\
C_{0 \alpha} & :=C_{0}+D_{0}\left[(1-\alpha) F_{\text {sub }}-\alpha D_{0}^{\prime} C_{0}\right]
\end{aligned}
$$



Figure 6.8: Performances of controllers M and C

The pairs $\left(A_{\alpha}, B_{2 \alpha}\right)$ and ( $A_{\alpha}, C_{0 \alpha}$ ) are stabilizable and detectable, respectively. In fact, stabilizability easily follows from stabilizability of $\left(A, B_{2}\right)$, whereas detectability is implied by detectability of the pair $\left(A-B_{2} D_{0}^{\prime} C_{0},\left(I-D_{0} D_{0}^{\prime}\right) C_{0}\right)$ and by noticing that

$$
\begin{aligned}
A_{\alpha} & =A-B_{2} D_{0}^{\prime} C_{0}+B_{2}(1-\alpha)\left[(1-\alpha)\left(F_{\text {sub }}+D_{0}^{\prime} C_{0}\right)\right] \\
C_{0 \alpha} & =\left(I-D_{0} D_{0}^{\prime}\right) C_{0}+D_{0}\left[(1-\alpha)\left(F_{\text {sub }}+D_{0}^{\prime} C_{0}\right)\right] \\
0 & =D_{0}^{\prime}\left(I-D_{0} D_{0}^{\prime}\right) C_{0}
\end{aligned}
$$

Such an equation (see Lemma C. 4 of Appendix C) admits a stabilizing solution whenever $\alpha \in[0,2]$. Actually, under the stated assumptions, this solution is also the unique positive semidefinite one. In order to check the stability of $A+B_{2} F_{\alpha}$, with $F_{\alpha}$ given by (6.80), consider equation (6.81) and notice that

$$
\begin{aligned}
A+B_{2} F_{\alpha} & =A-B_{2} D_{0}^{\prime} C_{0}+B_{2}\left[(1-\alpha)\left(F_{\text {sub }}+D_{0}^{\prime} C_{0}\right)-\alpha B_{2}^{\prime} P_{\alpha}\right] \\
C_{0}+D_{0} F_{\alpha} & =\left(I-D_{0} D_{0}^{\prime}\right) C_{0}+D_{0}\left[(1-\alpha)\left(F_{s u b}+D_{0}^{\prime} C_{0}\right)-\alpha B_{2}^{\prime} P_{\alpha}\right] \\
0 & =D_{0}^{\prime}\left(I-D_{0} D_{0}^{\prime}\right) C_{0}
\end{aligned}
$$

Thus, the detectability of the pair $\left(A-B_{2} D_{0}^{\prime} C_{0},\left(I-D_{0} D_{0}^{\prime}\right) C_{0}\right)$ together with the above equations imply that the pair $\left(A+B_{2} F_{\alpha}, C_{0}+D_{0} F_{\alpha}\right)$ is detectable as well. This fact and the existence of a solution $P_{\alpha} \geq 0$, entails, (see Lemma C. 1 of Appendix C) that $A+B_{2} F_{\alpha}$ is a stable matrix.

The above results can be exploited in the following way. With the state feedback gain $F_{\text {sub }}$ such that $A+B_{2} F_{\text {sub }}$ is stable and $\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty}<\gamma$ an one dimensional search in the interval $[0,2]$ for $\alpha$ with equation (6.82) taken into account allows to determine the value of $\alpha^{o}$ corresponding to which the control law $u=F_{\alpha} x$ minimizes the $\mathrm{RH}_{2}$ norm while keeping the $R H_{\infty}$ norm not greater than $\gamma$. Incidentally, notice that the choice $\alpha=1$ corresponds to the optimal unconstrained $\mathrm{RH}_{2}$ control law.

The solution $P_{\alpha}$ of equation (6.82) is a function of the parameter $\alpha$ which enjoys a symmetry property, namely $P_{\alpha}=P_{2-\alpha}$. Denoting with $\Gamma_{\alpha}$ the matrix which is the derivative of $P_{\alpha}$ with respect to $\alpha$ it is possible to verify that it satisfies the following Lyapunov equation

$$
\begin{equation*}
0=\bar{A}_{\alpha}^{\prime} \Gamma_{\alpha}+\Gamma_{\alpha} \bar{A}_{\alpha}-2(1-\alpha) \bar{F}_{s u b}^{\prime} \bar{F}_{s u b} \tag{6.83}
\end{equation*}
$$

where $\bar{A}_{\alpha}:=A_{\alpha}-B_{2 \alpha} B_{2 \alpha}^{\prime} P_{\alpha}$ is a stable matrix since $P_{\alpha}$ is a stabilizing solution of (6.82). Hence it follows from (6.83) that $\Gamma_{\alpha} \leq 0$ for $\alpha \in[0,1]$ and $\Gamma_{\alpha} \geq 0$ for $\alpha \in[1,2]$. Therefore
$\alpha^{\circ}$ is given by

$$
\begin{aligned}
& \alpha^{o}:=\min \left\{1-\alpha_{1}, \alpha_{2}-1\right\} \\
& \alpha_{1}:=\max _{\alpha \in[0,1]}\left\{\alpha:\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty} \leq \gamma\right\} \\
& \alpha_{2}:=\min _{\alpha \in[1,2]}\left\{\alpha:\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty} \leq \gamma\right\}
\end{aligned}
$$

Finally, it is important to stress that the above developments remain valid for any state feedback gain $F_{\text {sub }}$ such that $A+B_{2} F_{\text {sub }}$ is stable and $\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty}<\gamma$ and not just for the one provided by the solution of the mixed optimization problem (6.76). For instance, in view of the previous discussion (recall Example 6.2), we also have interest to adopt $F_{\text {sub }}$ as the matrix gain defined by the central $R H_{\infty}$ controller. Obviously, $\alpha^{\circ}$ and the relevant value for the $\mathrm{RH}_{2}$ norm depend on the chosen gain $F_{s u b}$.

To put in evidence the improvement obtained by this post-optimization procedure, let us consider the system defined in Example 6.2 with $C_{0}=[10]$ and $\gamma=2$. First take $F_{s u b}$ as the feedback gain which is the solution of the convex programming problem (6.76), that is $F_{s u b}=[-1.1879-0.5965]$ and the associated cost $\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}=0.1556$. Corresponding to the best choice $\alpha^{o}$, we obtain

$$
F_{\alpha}=\left[\begin{array}{cc}
-1.1829 & -0.5829
\end{array}\right],\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}=0.1504,\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty}=2
$$

If, on the contrary, we adopt $F_{s u b}$ as the state feedback gain corresponding to the central controller, that is $F_{s u b}=\left[\begin{array}{ll}-1.6 & -1.2\end{array}\right]$, the best choice of $\alpha$ yields

$$
F_{\alpha}=\left[\begin{array}{ll}
-1.2786 & -0.6076
\end{array}\right],\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}=0.1439,\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty}=2
$$

In this case, it is apparent that the post-optimization procedure supplies a better result when starting from a worse gain.

Remark 6.27 (Nash game approach) The mixed state feedback problem can also be approached by the following Nash game when the system (6.67) - (6.70) satisfies the additional assumptions $C_{0}=C_{1}, D_{0}=D_{12}$ and $D_{12}^{\prime} C_{1}=0$. Defining the criteria

$$
\begin{aligned}
& J_{1}(u, w)=\int_{0}^{\infty}\left[\gamma^{2} w^{\prime}(t) w(t)-z_{1}^{\prime}(t) z_{1}(t)\right] d t \\
& J_{0}(u, w)=\int_{0}^{\infty} z_{0}^{\prime}(t) z_{0}(t) d t
\end{aligned}
$$

the aim is to find the equilibrium strategies $\left(u^{\star}, w^{\star}\right)$ which satisfy the Nash equilibria conditions

$$
\begin{array}{ll}
J_{1}\left(u^{\star}, w^{\star}\right) \leq J_{1}\left(u^{\star}, w\right), \quad \forall w \in R H_{2} \\
J_{0}\left(u^{\star}, w^{\star}\right) \leq J_{0}\left(u, w^{\star}\right), \quad \forall u \in R H_{2}
\end{array}
$$

It can be proven that the optimal strategies are given by

$$
\begin{aligned}
u^{\star}(t) & =F_{2} x(t), & F_{2}:=-B_{2}^{\prime} P_{2} \\
w^{\star}(t) & =F_{1} x(t), & F_{1}:=-\gamma^{-2} B_{1}^{\prime} P_{1}
\end{aligned}
$$

provided there exist matrices $P_{1} \leq 0$ and $P_{2} \geq 0$ solutions to the coupled Riccati equations

$$
\begin{aligned}
& 0=A^{\prime} P_{1}+P_{1} A-C_{1}^{\prime} C_{1}-\left[\begin{array}{ll}
P_{1} & P_{2}
\end{array}\right]\left[\begin{array}{cc}
\gamma^{-2} B_{1} B_{1}^{\prime} & B_{2} B_{2}^{\prime} \\
B_{2} B_{2}^{\prime} & B_{2} B_{2}^{\prime}
\end{array}\right]\left[\begin{array}{l}
P_{1} \\
P_{2}
\end{array}\right] \\
& 0=A^{\prime} P_{2}+P_{2} A+C_{1}^{\prime} C_{1}-\left[\begin{array}{ll}
P_{1} & P_{2}
\end{array}\right]\left[\begin{array}{cc}
0 & \gamma^{-2} B_{1} B_{1}^{\prime} \\
\gamma^{-2} B_{1} B_{1}^{\prime} & B_{2} B_{2}^{\prime}
\end{array}\right]\left[\begin{array}{l}
P_{1} \\
P_{2}
\end{array}\right]
\end{aligned}
$$

Based on these optimality conditions, the following interpretation can be drawn. The first Riccati equation rewritten as

$$
0=\left(A+B_{2} F_{2}\right)^{\prime} P_{1}+P_{1}\left(A+B_{2} F_{2}\right)-\gamma^{-2} P_{1} B_{1} B_{1}^{\prime} P_{1}-\left(C_{1}+D_{12} F_{2}\right)^{\prime}\left(C_{1}+D_{12} F_{2}\right)
$$

implies that $\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty}<\gamma$ since $P_{1} \leq 0$ and stabilizing. On the other hand, the second Riccati equation, factorized as

$$
0=\left(A+B_{1} F_{1}\right)^{\prime} P_{2}+P_{2}\left(A+B_{1} F_{1}\right)-P_{2} B_{2} B_{2}^{\prime} P_{2}+C_{1}^{\prime} C_{1}
$$

shows that the criterion $J_{0}$ for $w=w^{\star}$ attains its minimum value at $u=u^{\star}$ as required by the Nash game. It is then clear that for $w \neq w^{\star}$ the control $u=F_{2} x$ is merely a suboptimal policy for the mixed design problem. Unfortunately this is frequently the case because for $\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty}<\gamma$ the worst input $w=w^{\star}$ is not the one which produces the output $z_{1}$ such that $\left\|z_{1}\right\|_{2}=\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty}\|w\|_{2}$. For comparison purpose let us consider the following numerical example

$$
\begin{aligned}
\dot{x} & =2 x+w+3 u \\
z_{1} & =\left[\begin{array}{l}
3 \\
0
\end{array}\right] x+\left[\begin{array}{l}
0 \\
1
\end{array}\right] u \\
z_{0} & =z_{1} \\
y & =x
\end{aligned}
$$

and $\gamma=0.4$. The pair $\left(P_{1}, P_{2}\right)=(-6.84,10.05)$ is a solution of the coupled Riccati equations which provides the optimal Nash gain $F_{2}=-30.16$. The closed-loop system exhibits the performances

$$
\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}=5.19, \quad\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty}=0.34<\gamma
$$

Then, problem (6.76) has also been solved. It provides the feedback gain $F=-8.05$ which imposes to the closed-loop system

$$
\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}=1.66, \quad\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty}=0.38<\gamma
$$

From these results it is clear that, in this case, the mixed design introduced in Theorem 6.8 is much better than the Nash game approach.

Remark 6.28 (Structured robust stability and performance) Consider a linear system depending on uncertain parameters, more precisely

$$
\begin{aligned}
\dot{x} & =\left(A+\Delta_{A}\right) x+B_{0} w+\left(B_{2}+\Delta_{B_{2}}\right) u \\
z_{0} & =C_{0} x+D_{0} u \\
y & =x
\end{aligned}
$$

with

$$
\Delta_{A}:=B_{11} \Omega_{1} C_{11}, \quad \Delta_{B_{2}}:=B_{12} \Omega_{2} C_{12}
$$

where the only information available for matrices $\Omega_{1}$ and $\Omega_{2}$ is that

$$
\left\|\Omega_{1}\right\| \leq 1, \quad\left\|\Omega_{2}\right\| \leq 1
$$

Defining the matrices

$$
B_{1}=\left[\begin{array}{ll}
B_{11} & B_{12}
\end{array}\right], \quad C_{1}=\left[\begin{array}{c}
C_{11} \\
0
\end{array}\right], \quad D_{12}=\left[\begin{array}{c}
0 \\
C_{12}
\end{array}\right],
$$

the above system can be rewritten as

$$
\begin{aligned}
\dot{x} & =A x+B_{0} w+B_{1} w_{1}+B_{2} u \\
z_{1} & =C_{1} x+D_{12} u \\
w_{1} & =\Omega z_{1} \\
z_{0} & =C_{0} x+D_{0} u \\
y & =x
\end{aligned}
$$

where matrix $\Omega:=\operatorname{diag}\left[\Omega_{1}, \Omega_{2}\right]$ satisfies $\|\Omega\| \leq 1$. Our first goal is to determine (if one exists) a matrix $F$ such that with $u=F x$ the closed-loop system stability is assured for all feasible parametric perturbations. A way we already know to solve this problem is to introduce the constraint $\left\|T\left(z_{1}, w_{1} ; s\right)\right\|_{\infty}<1$. Indeed, a controller which satisfies this constraint solve the stated problem. However, this approach produces frequently very conservative results since the block-diagonal structure of the uncertainty (represented by matrix $\Omega$ ) is not taken into account. In a slightly more general setting, we consider matrix $\Omega$ composed by $N$ square blocks, that is

$$
\Omega:=\operatorname{diag}\left[\Omega_{1}, \cdots, \Omega_{N}\right], \quad\|\Omega\| \leq 1
$$

which together with the matrix

$$
\Lambda:=\operatorname{diag}\left[\lambda_{1} I, \cdots, \lambda_{N} I\right]
$$

where each sub-block has the same dimension as the corresponding sub-block of matrix $\Omega$ and $\lambda_{1}>0, \cdots, \lambda_{N}>0$, enable us to say that $\Omega \Lambda \Omega^{\prime} \leq \Lambda, \quad \forall\|\Omega\| \leq 1$.

The closed loop system stability depends on the stability of matrix

$$
A_{\Omega}:=A_{c \infty}+B_{1} \Omega C_{c \infty}
$$

where as before, $A_{\mathrm{c} \infty}=A+B_{2} F$ and $C_{c \infty}=C_{1}+D_{12} F$. This property follows from the following inequality which holds for any symmetric matrix $P$

$$
\begin{aligned}
B_{1} \Omega C_{c \infty} P+P C_{c \infty}^{\prime} \Omega^{\prime} B_{1}^{\prime}= & B_{1} \Omega \Lambda \Omega^{\prime} B_{1}^{\prime}+P C_{c \infty}^{\prime} \Lambda^{-1} C_{c \infty} P- \\
& -\left(B_{1} \Omega \Lambda^{1 / 2}-P C_{c \infty}^{\prime} \Lambda^{-1 / 2}\right)\left(B_{1} \Omega \Lambda^{1 / 2}-P C_{c \infty}^{\prime} \Lambda^{-1 / 2}\right)^{\prime} \\
\leq & B_{1} \Lambda B_{1}^{\prime}+P C_{c \infty}^{\prime} \Lambda^{-1} C_{c \infty} P
\end{aligned}
$$

Indeed, if there exist a symmetric and positive definite matrix $P$, a positive definite matrix $\Lambda$ with the above structure and a matrix $F$ such that the inequality

$$
A_{c \infty} P+P A_{c \infty}^{\prime}+P C_{c \infty}^{\prime} \Lambda^{-1} C_{c \infty} P+B_{1} \Lambda B_{1}^{\prime}<0
$$

holds then (recall Theorem 2.15) the transfer function

$$
T_{\Lambda}\left(z_{1}, w_{1} ; s\right):=\Lambda^{-1 / 2} T\left(z_{1}, w_{1} ; s\right) \Lambda^{1 / 2}
$$

belongs to $R H_{\infty}$ and $\left\|T_{\Lambda}\left(z_{1}, w_{1} ; s\right)\right\|_{\infty}<1$. Since the scaling matrix $\Lambda$ can also be used to redefine the parametric perturbation as $\Omega_{\Lambda}:=\Lambda^{-1 / 2} \Omega \Lambda^{1 / 2}$ such that $\left\|\Omega_{\Lambda}\right\|_{\infty} \leq 1$, from Theorem 5.1 the stability of $A_{\Omega}$ follows. In addition, the above property is not lost if the previous Riccati inequality is replaced by (assuming again that $B_{0} B_{0}^{\prime}>0$ )

$$
A_{c \infty} P+P A_{c \infty}^{\prime}+P C_{c \infty}^{\prime} \Lambda^{-1} C_{c \infty} P+B_{1} \Lambda B_{1}^{\prime}+B_{0} B_{0}^{\prime} \leq 0
$$

However, doing this we now have

$$
A_{\Omega} P+P A_{\Omega}^{\prime}+B_{0} B_{0}^{\prime} \leq 0, \quad \forall\|\Omega\| \leq 1
$$

yielding an upper bound to the $\mathrm{RH}_{2}$ norm of the transfer function from the input $w$ to the output $z_{0}$, that is

$$
\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2} \leq \operatorname{trace}\left[\left(C_{0}+D_{0} F\right) P\left(C_{0}+D_{0} F\right)^{\prime}\right]
$$

which is valid for all $\|\Omega\| \leq 1$. The minimization of this upper bound keeping stability is a mixed $R H_{2} / R H_{\infty}$ problem with an additional matrix variable, namely the scaling matrix $\Lambda$. It is somewhat surprising that this problem can be converted into a convex one. Actually, let us define the set

$$
\mathcal{C}_{\Lambda c}:=\left\{(W, \Lambda): W \geq 0, \Lambda>0, v^{\prime} \Theta_{\Lambda c}(W, \Lambda) v \leq 0 \forall v \in \mathcal{N}_{c}\right\}
$$

where

$$
\begin{aligned}
\Theta_{\Lambda c}(W, \Lambda):= & M_{c} W+W M_{c}^{\prime}+W\left[\begin{array}{c}
C_{1}^{\prime} \\
D_{12}^{\prime}
\end{array}\right] \Lambda^{-1}\left[\begin{array}{ll}
C_{1} & D_{12}
\end{array}\right] W+ \\
& +\left[\begin{array}{c}
B_{1} \\
0
\end{array}\right] \Lambda\left[\begin{array}{ll}
B_{1}^{\prime} & 0
\end{array}\right]+\left[\begin{array}{c}
B_{0} \\
0
\end{array}\right]\left[\begin{array}{ll}
B_{0}^{\prime} & 0
\end{array}\right]
\end{aligned}
$$

Following the proof of Theorem 6.6, we see that it remains true in the present case. The set $\mathcal{C}_{\Lambda c}$ is convex and any feasible pair $(W, \Lambda) \in \mathcal{C}_{\Lambda c}$ provides $F=W_{2}^{\prime} W_{1}^{-1}, P=W_{1}$ and $\Lambda$ which satisfy the previous Riccati inequality. Finally, in the present context, the associated mixed $R H_{2} / R H_{\infty}$ control design problem is (recall Theorem 6.8)

$$
J_{s u b}:=\min \left\{\operatorname{trace}\left[R_{0} W\right]:(W, \Lambda) \in \mathcal{C}_{\Lambda c}\right\}
$$

which is jointly convex on both variables $(W, \Lambda)$. Its optimal solution provides a robust control gain $F=W_{2}^{\prime} W_{1}^{-1}$ imposing in addition $\left\|T\left(z_{0}, w ; s\right)\right\|_{2} \leq \sqrt{J_{s u b}}$ for all structured parametric perturbations $\|\Omega\| \leq 1$.

Remark 6.29 Following the same lines of Remark 6.22, any convex programming method is effective to solve problem (6.76) provided the convex set

$$
\mathcal{C}_{\gamma c} \cap\left\{W: \operatorname{trace}\left[R_{0} W\right] \leq \beta\right\}
$$

is bounded for any finite $\beta>0$. Based on our previous discussion, this occurs whenever the pair $\left[-\left(A-B_{2} D_{0}^{\prime} C_{0}\right),\left(I-D_{0} D_{0}^{\prime}\right) C_{0}\right]$ is detectable.

### 6.3.2 Output feedback design

In this case, with $\gamma>0$ being a fixed scalar the system $\Sigma$ is defined as follows

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{6.84}\\
z_{0} & =C_{0} x+D_{0} u  \tag{6.85}\\
z_{1} & =C_{1} x+D_{12} u  \tag{6.86}\\
y & =C_{2} x+D_{21} w \tag{6.87}
\end{align*}
$$

and the controller $\Sigma_{c}$ has to be determined from the solution of problem (6.65). As in the State feedback case, the complete solution to this problem is not known up to now. So, we search for a suboptimal and easy to calculate solution. The main idea to be pursued is to propose a structure to the controller $\Sigma_{c}$ depending on only one unknown matrix in such way the Output feedback design problem reduces to the State feedback design problem already solved.

To this end, we make the following assumptions i) the pair $\left(A, B_{2}\right)$ is stabilizable and the pair $\left(A, C_{2}\right)$ is detectable and ii) $D_{12}^{\prime} D_{12}=I$ and $D_{21} D_{21}^{\prime}=I$. We also assume that there exists $\Pi_{\infty}$ a positive semidefinite and stabilizing solution of the Riccati equation (in the unknown $\Pi$ )

$$
\begin{equation*}
0=\Pi A_{f}^{\prime}+A_{f} \Pi-\Pi\left(C_{2}^{\prime} C_{2}-\gamma^{-2} C_{1}^{\prime} C_{1}\right) \Pi+B_{1 f} B_{1 f}^{\prime} \tag{6.88}
\end{equation*}
$$

which, once it has been calculated, enables us to define matrices

$$
\begin{aligned}
A_{\infty} & :=A+\gamma^{-2} \Pi_{\infty} C_{1}^{\prime} C_{1} \\
B_{2 \infty} & :=B_{2}+\gamma^{-2} \Pi_{\infty} C_{1}^{\prime} D_{12}
\end{aligned}
$$

and $L_{\infty}=-\Pi_{\infty} C_{2}^{\prime}-B_{1} D_{21}^{\prime}$. The controller $\Sigma_{c}$ with transfer function $K(s)$ given by

$$
K(s):=\left[\begin{array}{c|c}
A_{\infty}+B_{2 \infty} F+L_{\infty} C_{2} & -L_{\infty}  \tag{6.89}\\
\hline F & 0
\end{array}\right]
$$

has the important properties provided in the following lemma.
Lemma 6.2 For all $F$ such that there exists a symmetric, positive semidefinite and stabilizing solution to the Riccati equation

$$
\begin{align*}
0= & \left(A_{\infty}+B_{2 \infty} F\right) X+X\left(A_{\infty}+B_{2 \infty} F\right)^{\prime}+ \\
& +\gamma^{-2} X\left(C_{1}+D_{12} F\right)^{\prime}\left(C_{1}+D_{12} F\right) X+L_{\infty} L_{\infty}^{\prime} \tag{6.90}
\end{align*}
$$

the controller $\Sigma_{c}$ with transfer function $K(s)$ given in (6.89) imposes to the closed loop system the following properties:
a) It is stable and $\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty}<\gamma$
b) $\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2} \leq \operatorname{trace}\left[C_{0} \Pi_{\infty} C_{0}^{\prime}\right]+\operatorname{trace}\left[\left(C_{0}+D_{0} F\right) X\left(C_{0}+D_{0} F\right)^{\prime}\right]$

Proof The feedback connection indicated in fig. 6.6 has the state space representation

$$
\Sigma_{F}:=\left[\begin{array}{c|c}
\tilde{A} & \tilde{B} \\
\hline \tilde{C}_{0} & 0 \\
\tilde{C}_{1} & 0
\end{array}\right]
$$

where the indicated matrices are

$$
\begin{aligned}
\tilde{A} & =\left[\begin{array}{cc}
A+B_{2} F & B_{2} F \\
\gamma^{-2} \Pi_{\infty} C_{1}^{\prime}\left(C_{1}+D_{12} F\right) & A+L_{\infty} C_{2}+\gamma^{-2} \Pi_{\infty} C_{1}^{\prime}\left(C_{1}+D_{12} F\right)
\end{array}\right] \\
\tilde{B} & =\left[\begin{array}{c}
B_{1} \\
-B_{1}-L_{\infty} D_{21}
\end{array}\right]
\end{aligned}
$$

and

$$
\tilde{C}_{0}=\left[\begin{array}{ll}
C_{0}+D_{0} F & D_{0} F
\end{array}\right], \quad \tilde{C}_{1}=\left[\begin{array}{ll}
C_{1}+D_{12} F & D_{12} F
\end{array}\right]
$$

Point a) Assuming there exists $X_{\infty}$ satisfying (6.90), simple although tedious algebraic calculations show that the Riccati equation

$$
\begin{equation*}
0=\tilde{A} \tilde{P}+\tilde{P} \tilde{A}^{\prime}+\gamma^{-2} \tilde{P} \tilde{C}_{1}^{\prime} \tilde{C}_{1} \tilde{P}+\tilde{B} \tilde{B}^{\prime} \tag{6.91}
\end{equation*}
$$

has a positive semidefinite stabilizing solution given by

$$
\tilde{P}=\left[\begin{array}{cc}
\Pi_{\infty}+X_{\infty} & -\Pi_{\infty}  \tag{6.92}\\
-\Pi_{\infty} & \Pi_{\infty}
\end{array}\right]
$$

which yields the conclusion from Theorem 2.14 that the closed loop system is stable and $\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty}<\gamma$.

Point b) From the above, it is know that matrix $\tilde{A}$ is stable then, we proceed by making the calculation of the indicated norm, that is

$$
\begin{aligned}
\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2} & =\operatorname{trace}\left[\tilde{C}_{0} \int_{0}^{\infty} e^{\tilde{A} t} \tilde{B} \tilde{B}^{\prime} e^{\tilde{A}^{\prime} t} d t \tilde{C}_{0}^{\prime}\right] \\
& \leq \operatorname{trace}\left[\tilde{C}_{0} \tilde{P} \tilde{C}_{0}^{\prime}\right] \\
& \leq \operatorname{trace}\left[C_{0} \Pi_{\infty} C_{0}^{\prime}\right]+\operatorname{trace}\left[\left(C_{0}+D_{0} F\right) X\left(C_{0}+D_{0} F\right)^{\prime}\right]
\end{aligned}
$$

concluding thus the proof of the lemma proposed.
No major difficulty has to be faced to state and prove similar results using Riccati inequalities. Furthermore, it is to be noticed that the assumption on the existence of a solution for the Riccati equation (6.90) is always verified whenever the unknown matrix gain $F$ is such that $A_{\infty}+B_{2 \infty} F$ is stable and

$$
\left\|\left(C_{1}+D_{12} F\right)\left[s I-\left(A_{\infty}+B_{2 \infty} F\right)\right]^{-1} L_{\infty} D_{21}\right\|_{\infty}<\gamma
$$

Therefore, from the first part of Lemma 6.2 these values of $F$ preserve admissibility of the controller $\Sigma_{c}$, that is $\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty}<\gamma$ and in the context of the mixed design we have to calculate one among them which minimizes the upper bound on $\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}$ provided in the second part of the same lemma. This is accomplished if the State feedback design is applied to the auxiliary plant

$$
\begin{align*}
\dot{x} & =A_{\infty} x+L_{\infty} D_{21} w+B_{2 \infty} u  \tag{6.93}\\
z_{0} & =C_{0} x+D_{0} u  \tag{6.94}\\
z_{1} & =C_{1} x+D_{12} u  \tag{6.95}\\
y & =x \tag{6.96}
\end{align*}
$$

that is, problem (6.76) should be solved with the feasible set $\mathcal{C}_{\gamma c}$ being defined for matrices $A, B_{2}$ and $B_{1}$ replaced by matrices $A_{\infty}, B_{2 \infty}$ and $L_{\infty} D_{21}$ respectively.
Remark 6.30 The particular structure of matrix $\tilde{P}$ in (6.92) does not means that some conservativeness has been introduced in our calculations. To show this assume the Partial information problem is solvable and $P_{\infty}>0$. It is possible to verify that for

$$
F=F_{c}:=F_{\infty} Z_{\infty}
$$

the Riccati equation (6.90) is solvable in $X$ providing

$$
X_{\infty}=\gamma^{2} Z_{\infty}^{-1} P_{\infty}^{-1}=\gamma^{2} P_{\infty}^{-1}-\Pi_{\infty}
$$

which is positive definite since $r_{s}\left(P_{\infty} \Pi_{\infty}\right)<\gamma^{2}$. Moreover, for $F=F_{c}$, the transfer function of the controller $K(s)$ turns out to be

$$
\begin{aligned}
K_{c}(s) & =\left[\begin{array}{c|c}
A_{\infty}+B_{2 \infty} F_{c}+L_{\infty} C_{2} & -L_{\infty} \\
\hline F_{c} & 0
\end{array}\right] \\
& =\left[\begin{array}{c|c}
Z_{\infty}\left(A_{\infty}+B_{2 \infty} F_{\infty} Z_{\infty}+L_{\infty} C_{2}\right) Z_{\infty}^{-1} & -Z_{\infty} L_{\infty} \\
\hline F_{\infty} & 0
\end{array}\right] \\
& =\left[\begin{array}{c|c}
A_{c c}+Z_{\infty} L_{\infty}\left(C_{2}+\gamma^{-2} D_{21} B_{1}^{\prime} P_{\infty}\right) & -Z_{\infty} L_{\infty} \\
\hline F_{\infty} & 0
\end{array}\right]
\end{aligned}
$$

showing that the transfer function $K_{c}(s)$ meets exactly the transfer function of the central $R H_{\infty}$ controller.


Figure 6.9: The performances of the controllers M and C

Example 6.3 We solve again the Example 5.7 but making use of the results of this section. Consider the system (6.84) - (6.87) with

$$
\begin{gathered}
A=\left[\begin{array}{ll}
0 & 1 \\
2 & 1
\end{array}\right], \quad B_{1}=\left[\begin{array}{l}
0 \\
0
\end{array}\right], \quad B_{2}=\left[\begin{array}{l}
0 \\
1
\end{array}\right] \\
C_{1}=\left[\begin{array}{ll}
0 & 0
\end{array}\right], \quad C_{2}=\left[\begin{array}{ll}
-1 & 1
\end{array}\right], \quad D_{12}=D_{21}=1 \\
C_{0}=\left[\begin{array}{ll}
-1 & 0
\end{array}\right], \quad D_{0}=1
\end{gathered}
$$

The Partial information problem is solvable for $\gamma>12.0$ hence we consider the interval $12.5 \leq \gamma \leq 100.0$. In this particular case the auxiliary plant does not depend on $\gamma$, however the Riccati equation (6.90) admits a positive solution only if $\gamma>12.0$. The State feedback design has been applied to the auxiliary plant and the mixed problem provided the controller $M$ which minimizes the upper bound introduced in part b) of Lemma 6.2. Fig. 6.9 shows the actual value of $\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}$ as well as the minimum upper bound $J_{s u b}$ as $\gamma$ varies in the given interval. The same figure shows also the performance in terms of $\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}$ when the $R H_{\infty}$ central controller is used (it is indicated by $C$ ). It is interesting to observe that for moderate values of $\gamma$ (approximately in the interval $12.5 \leq \gamma \leq 35.0$ ) the central controller is even better than the mixed controller as far as the $R H_{2}$ norm of the transfer function $T\left(z_{0}, w ; s\right)$ is concerned. As expected, this behavior is reversed for large values of $\gamma$ when so the central controller performs worse than the mixed controller. This example is important because it shows practically the existence of systems for which the mixed design does not furnish a good solution to the problem under consideration. However, it is important to stress that the mixed problem as introduced here, is convex and due to this fact it can be solved very efficiently.

## 6.4 $\mathrm{RH}_{2}$ control with regional pole placement

The performance of a system can be expressed in terms of $R H_{2}$ and $R H_{\infty}$ of certain closed loop transfer functions. However, as it is simple to notice from the results of the last sections, the pole locations of the resulting controlled system are naturally defined by the optimality conditions of the associated optimal control problem. For


Figure 6.10: The pole placement region
instance, once a mixed $R \mathrm{H}_{2} / R H_{\infty}$ design problem is solved, the optimal controller imposes the closed-loop poles to be in certain places of the open left complex plane which may not be changed by the designer. The present section aims to combine the $R \mathrm{H}_{2}$ design with regional pole placement.

In practice, a very popular design specification is expressed as

$$
\begin{equation*}
\xi \geq \xi_{\min } \tag{6.97}
\end{equation*}
$$

and

$$
\begin{equation*}
\xi \omega_{n} \geq \alpha \tag{6.98}
\end{equation*}
$$

where the constraint (6.97) imposes a minimum damping ratio and the constraint (6.98) imposes a minimum decay of the time response of the closed-loop system. This is illustrated in Fig. 6.10 where it is also represented the circular region

$$
\begin{equation*}
\mathcal{R}:=\{s:|s+(\alpha+r)|<r\} \tag{6.99}
\end{equation*}
$$

with $\alpha>0$ and radius $r>0$. The radius $r$ can be calculated such that this region is inside the sector defined by constraints (6.97) and (6.98) and so meets the closed loop poles location requirements. The circular region as indicated in Fig. 6.10, tangent to the sector boundary defined by $\alpha>0$ and $\theta=\sin ^{-1}\left(\xi_{\min }\right)$, is given in (6.99) with

$$
r=\frac{\alpha \cos (\theta)}{1-\cos (\theta)}
$$

The rationale behind the choice of the circular region $\mathcal{R}$ instead of the sector will be clear in the sequel. The main point is that the circular region imposes convex constraints on the optimal $\mathrm{RH}_{2}$ control design problem. The next lemma characterizes by means of a modified Lyapunov inequality the matrices with all eigenvalues inside a given circular region $\mathcal{R}$.

Lemma 6.3 Let the circular $\mathcal{R}$ be given. Matrix $A$ with dimension $n \times n$ has all its eigenvalues inside $\mathcal{R}$ if and only if for any matrix $Q=Q^{\prime}>0$, there exists a matrix $P=P^{\prime}>0$ such that

$$
\begin{equation*}
0 \geq A_{\alpha} P+P A_{\alpha}^{\prime}+r^{-1} A_{\alpha} P A_{\alpha}^{\prime}+Q \tag{6.100}
\end{equation*}
$$

where $A_{\alpha}:=A+\alpha I$.

Proof Let us first prove the necessity. Assume all eigenvalues of matrix $A$ are in the circular region $\mathcal{R}$. It is a simple matter to verify that all eigenvalues $\tilde{\lambda}$ of the matrix $\tilde{A}:=r^{-1}[A+(\alpha+r) I]$ are such that $|\tilde{\lambda}|<1$ which implies that matrix

$$
\bar{A}:=(\tilde{A}-I)^{-1}(\tilde{A}+I)
$$

is stable. From the Extended Lyapunov lemma we can say that for any $Q=Q^{\prime}>0$ chosen, there exits $P=P^{\prime}>0$ satisfying the linear inequality

$$
0 \geq \bar{A} P+P \bar{A}^{\prime}+2 r^{-1}(\tilde{A}-I)^{-1} Q\left(\tilde{A}^{\prime}-I\right)^{-1}
$$

or equivalently

$$
0 \geq \tilde{A} P \tilde{A}^{\prime}-P+r^{-1} Q
$$

Using the definition of matrix $A_{\alpha}$ we then conclude that inequality (6.100) holds. For the sufficiency, take $x$ an eigenvector of matrix $A^{\prime}$ associated to an arbitrary eigenvalue $\lambda$. Multiplying inequality (6.100) to left by $x^{\sim}$ and to the right by $x$ it follows that

$$
0 \geq\left[2 R e(\lambda+\alpha)+r^{-1}|\lambda+\alpha|^{2}\right] x^{\sim} P x+x^{\sim} Q x
$$

which together with the positive definiteness of both involved matrices provides

$$
\left[2 \operatorname{Re}(\lambda+\alpha)+r^{-1}|\lambda+\alpha|^{2}\right]<0
$$

Finally, using this fact we get

$$
\begin{aligned}
r^{-2}|\lambda+\alpha+r|^{2} & =r^{-2}\left[|\lambda+\alpha|^{2}+2 \operatorname{Re}(\lambda+\alpha) r+r^{2}\right] \\
& =r^{-1}\left[2 \operatorname{Re}(\lambda+\alpha)+r^{-1}|\lambda+\alpha|^{2}\right]+1 \\
& <1
\end{aligned}
$$

which proves that $\lambda \in \mathcal{R}$. The sufficiency is proved because $A$ and $A^{\prime}$ have the same eigenvalues.

For a given circular region $\mathcal{R}$ and a given matrix $A$, the inequality (6.100) defines a convex constraint with respect to $P$. This is clearly true because its right hand side is an affine function of $P$. More surprisingly is that convexity still holds when $A$ is not constant but depends upon a state feedback gain matrix.

### 6.4.1 State feedback design

Let a circular region $\mathcal{R}$ be given. The dynamic system under consideration has the following state space representation

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{6.101}\\
z & =C_{1} x+D_{12} u  \tag{6.102}\\
y & =x \tag{6.103}
\end{align*}
$$

where the only assumption we a priori need is i) $D_{12}^{\prime} D_{12}=I$. The goal is to determine a state feedback control law of the form

$$
\begin{equation*}
u=F x \tag{6.104}
\end{equation*}
$$

such that the gain matrix $F$ solves the optimal control problem

$$
\min \left\{\|T(z, w ; s)\|_{2}^{2}: F \in \mathcal{K}_{R}\right\}
$$

where $\mathcal{K}_{R}$ denotes the set of matrices $F$ such that all eigenvalues of matrix $A_{c c}=$ $A+B_{2} F$ are in the circular region $\mathcal{R}$. It is clear that since $\mathcal{R}$ is a circular region on the left part of the complex plane then $F \in \mathcal{K}_{R}$ always provides $T(z, w ; s)$ in $R H_{2}$. Of course, we need the set $\mathcal{K}_{R}$ be nonempty which is the case if and only if the unreachable part of $\left(A, B_{2}\right)$ has all eigenvalues belonging to $\mathcal{R}$. This problem is not exactly solved. Instead we propose here an overbounding objective function which has the main advantage to preserve convexity. Throughout this section we redefine the matrix

$$
M_{c}:=\left[\begin{array}{cc}
A+\alpha I & B_{2} \\
0 & 0
\end{array}\right]
$$

which has the same structure of matrix $M_{c}$ defined in (6.6). All other matrices remain unchanged.

Theorem 6.9 Consider $B_{1} B_{1}^{\prime}>0$, let the circular region $\mathcal{R}$ be given and define the set

$$
\begin{equation*}
\mathcal{C}_{R}:=\left\{W: W \geq 0, \quad v^{\prime} \Theta_{R}(W) v \leq 0, \forall v \in \mathcal{N}_{c}\right\} \tag{6.105}
\end{equation*}
$$

where $\Theta_{R}(W):=\Theta_{c}(W)+r^{-1} M_{c} W M_{c}^{\prime}$. The following hold
a) $\mathcal{C}_{R}$ is a convex set.
b) Each $W \in \mathcal{C}_{R}$ is such that $W_{1}>0$.
c) $\mathcal{K}_{R}=\left\{W_{2}^{\prime} W_{1}^{-1}: W \in \mathcal{C}_{R}\right\}$.
d) The optimal solution $\bar{W}$ of the convex programming problem

$$
\begin{equation*}
J_{s u b}:=\min \left\{\operatorname{trace}\left[R_{c} W\right]: W \in \mathcal{C}_{R}\right\} \tag{6.106}
\end{equation*}
$$

provides $F=\bar{W}_{2}^{\prime} \bar{W}_{1}^{-1} \in \mathcal{K}_{R}$ such that the upper bound of the $R H_{2}$ norm $\|T(z, w ; s)\|_{2} \leq \operatorname{trace}\left[R_{c} W\right]$, valid for all $W \in \mathcal{C}_{R}$, is minimized.

Proof The proof of Points a) and b) follow immediately from the fact that $\mathcal{C}_{R} \subset \mathcal{C}_{c}$ and the matrix function $\Theta_{R}(W)$ is affine together with Theorem 6.1.

Point c) Assume $\mathcal{K}_{R} \neq \emptyset$, for an arbitrary $F \in \mathcal{K}_{R}$, Lemma 6.3 applies for $Q=B_{1} B_{1}^{\prime}$ providing thus a symmetric and positive definite matrix $P$ satisfying the inequality

$$
\begin{aligned}
0 \geq & \left(A_{c c}+\alpha I\right) P+P\left(A_{c c}+\alpha I\right)^{\prime}+ \\
& +r^{-1}\left(A_{c c}+\alpha I\right) P\left(A_{c c}+\alpha I\right)^{\prime}+B_{1} B_{1}^{\prime}
\end{aligned}
$$

Choosing

$$
W=\left[\begin{array}{cc}
P & P F^{\prime} \\
F P & F P F^{\prime}
\end{array}\right]
$$

it is seen that $W \geq 0$ and all $v \in \mathcal{N}_{c}$ yields

$$
\begin{aligned}
v^{\prime} \Theta_{R}(W) v^{\prime}= & x^{\prime}\left[\left(A_{c c}+\alpha I\right) P+P\left(A_{c c}+\alpha I\right)^{\prime}+\right. \\
& \left.+r^{-1}\left(A_{c c}+\alpha I\right) P\left(A_{c c}+\alpha I\right)^{\prime}+B_{1} B_{1}^{\prime}\right] x \\
\leq & 0
\end{aligned}
$$

implying that $W \in \mathcal{C}_{R}$. Since, in addition $W_{2}^{\prime} W_{1}^{-1}=F P P^{-1}=F$, the necessity is proved. Conversely, with $\mathcal{C}_{R} \neq \emptyset$, for any $W \in \mathcal{C}_{R}$ and any $v \in \mathcal{N}_{c}$ we get

$$
\begin{aligned}
0 \geq & v^{\prime} \Theta_{R}(W) v \\
\geq & x^{\prime}\left[\left(A_{\alpha}+B_{2} W_{2}^{\prime} W_{1}^{-1}\right) W_{1}+W_{1}\left(A_{\alpha}+B_{2} W_{2}^{\prime} W_{1}^{-1}\right)^{\prime}+\right. \\
& +r^{-1}\left(A_{\alpha}+B_{2} W_{2}^{\prime} W_{1}^{-1}\right) W_{1}\left(A_{\alpha}+B_{2} W_{2}^{\prime} W_{1}^{-1}\right)^{\prime}+B_{1} B_{1}^{\prime}+ \\
& \left.+r^{-1} B_{2}\left(W_{3}-W_{2}^{\prime} W_{1}^{-1} W_{2}\right) B_{2}^{\prime}\right] x \\
\geq & x^{\prime}\left[\left(A_{\alpha}+B_{2} W_{2}^{\prime} W_{1}^{-1}\right) W_{1}+W_{1}\left(A_{\alpha}+B_{2} W_{2}^{\prime} W_{1}^{-1}\right)^{\prime}+\right. \\
& \left.+r^{-1}\left(A_{\alpha}+B_{2} W_{2}^{\prime} W_{1}^{-1}\right) W_{1}\left(A_{\alpha}+B_{2} W_{2}^{\prime} W_{1}^{-1}\right)^{\prime}+B_{1} B_{1}^{\prime}\right] x
\end{aligned}
$$

which shows, from Lemma 6.3 that $F=W_{2}^{\prime} W_{1}^{-1} \in \mathcal{K}_{R}$. The result follows from the fact that $\mathcal{K}_{R}=\emptyset$ implies $\mathcal{C}_{R}=\emptyset$ and vice versa.

Point d) This point follows from the fact that $\vec{W} \in \mathcal{C}_{R}$ generates $F=\bar{W}_{2}^{\prime} \bar{W}_{1}^{-1} \in$ $\mathcal{K}_{R}$. Furthermore, any $W \in \mathcal{C}_{R}$ satisfies

$$
0 \geq A_{c c} W_{1}+W_{1} A_{c c}^{\prime}+B_{1} B_{1}^{\prime}
$$

which imposes that $W_{1} \geq P_{r}$, where

$$
0=A_{c c} P_{r}+P_{r} A_{c c}^{\prime}+B_{1} B_{1}^{\prime}
$$

Simple calculation of the $R \mathrm{H}_{2}$ norm then shows that

$$
\begin{aligned}
\operatorname{trace}\left[R_{c} W\right]= & \operatorname{trace}\left[\left(C_{1}+D_{12} W_{2}^{\prime} W_{1}^{-1}\right) W_{1}\left(C_{1}+D_{12} W_{2} W_{1}^{-1}\right)^{\prime}\right]+ \\
& +\operatorname{trace}\left[D_{12}\left(W_{3}-W_{2}^{\prime} W_{1}^{-1} W_{2}\right) D_{12}^{\prime}\right] \\
\geq & \operatorname{trace}\left[\left(C_{1}+D_{12} W_{2} W_{1}^{-1}\right) P_{r}\left(C_{1}+D_{12} W_{2} W_{1}^{-1}\right)^{\prime}\right] \\
\geq & \|T(z, w ; s)\|_{2}^{2}
\end{aligned}
$$

holds for all $W \in \mathcal{C}_{R}$ and the proof is complete.
If the convex programming problem (6.106) admits a solution then the optimal feedback gain places the closed-loop poles in a desired circular region and minimizes an upper bound of $\|T(z, w ; s)\|_{2}^{2}$. Clearly, we can not say that this norm has been minimized.

Remark 6.31 Since $\mathcal{C}_{R} \subset \mathcal{C}_{c}$, the same conclusion related to the boundedness of the feasible set of problem (6.106) applies. Its numeric solution by means of convex programming methods depends on the assumption that the pair $\left(-A_{c}, C_{1 c}\right)$ is detectable.

Remark 6.32 As perhaps already occurred to the reader, there are many other regions in the complex plane that can be recasted in the same framework of the circular region.

For $\alpha>0$ fixed and $r$ arbitrarily large, the circular region degenerates to

$$
\mathcal{R}=\{s: \operatorname{Re}(s)<-\alpha\}
$$

and Lemma 6.3 and Theorem 6.9 still hold.
Any region in the complex plane that can be written as a convex set of matrices $W$ such that $W_{1} \geq P_{r}$, is also handled with no additional theoretical difficulty. One of such regions is the vertical strip defined as

$$
\mathcal{R}_{S}:=\{s:-\beta<\operatorname{Re}(s)<-\alpha\}
$$

with $\beta>\alpha>0$. Notice however, that we have to work with the intersection of two regions generating thus a feasible set which is the intersection of two convex constraints in the same variable $W$. Consequently, only the sufficient part of Theorem 6.9 still holds.


Figure 6.11: Closed-loop poles

Example 6.4 (Root locus) Let us consider the system (6.101) - (6.103) with

$$
\begin{gathered}
A=\left[\begin{array}{rr}
0 & 1 \\
-1 & 1
\end{array}\right], \quad B_{1}=\left[\begin{array}{l}
1 \\
0
\end{array}\right], \quad B_{2}=\left[\begin{array}{l}
0 \\
1
\end{array}\right] \\
C_{1}=\left[\begin{array}{ll}
\rho & 0 \\
0 & 0
\end{array}\right], \quad D_{12}=\left[\begin{array}{l}
0 \\
1
\end{array}\right]
\end{gathered}
$$

where $0 \leq \rho \leq 20$.
First, the State feedback problem in $R H_{2}$ has been solved for all values of $\rho$ in the given interval. The optimal feedback gain $F_{2}$ as a function of this parameter has been determinated. The root locus of the closed-loop system is then plotted in fig. 6.11. By inspection, it is possible to verify that the root locus never enters in the circular region defined by the parameters $\alpha=1$ and $r=1 /(\sqrt{2}-1)$. For $\rho=0$ the optimal gain is $F_{2}=\left[\begin{array}{ll}0 & -2\end{array}\right]$ which imposes to the system a performance such that $\|T(z, w ; s)\|_{2}^{2}=2$.

Again, for $\rho=0$ and imposing the above circular region for pole placement, the convex programming problem (6.106) has been solved, providing $F_{s u b}=[-3.0678-4.9841]$ and the associated upper bound on the minimum cost $J_{s u b}=51.9434$. The closed-loop poles are also shown in fig. 6.11. This example illustrates a very important fact. Indeed, there is no possibility to choose a penalty term in the $(1,1)$ entry of matrix $C_{1}$ such that the poles of the closed-loop system are all inside the circular region which corresponds to a minimum damping factor $\xi_{\text {min }}=1 / \sqrt{2}$ or in other terms $\theta=\pi / 4$ radians. Of course the design procedure introduced in Theorem 6.9 circumvents this drawback of the pure State feedback design in $\mathrm{RH}_{2}$.

### 6.4.2 Output feedback design

The system to be dealt with is of the form

$$
\begin{align*}
& \dot{x}=A x+B_{1} w+B_{2} u  \tag{6.107}\\
& z=C_{1} x+D_{12} u  \tag{6.108}\\
& y=C_{2} x+D_{21} w \tag{6.109}
\end{align*}
$$

and the controller with transfer function $K(s)$ has to be determined in order to place the closed-loop poles in a desired circular region while the quantity $\|T(z, w ; s)\|_{2}$ is minimized. Let us elaborate more on this point. First all feasible controllers are parametrized as

$$
K(s):=\left[\begin{array}{c|c}
A+B_{2} F+L C_{2} & -L  \tag{6.110}\\
\hline F & 0
\end{array}\right]
$$

for all pairs of matrices $(F, L)$. The rationale behind this choice has been put on evidence before. The closed-loop system has the state space representation (recall fig. 6.5)

$$
\Sigma_{F}:=\left[\begin{array}{c|c}
\tilde{A} & \tilde{B} \\
\hline \tilde{C} & 0
\end{array}\right]
$$

where the indicated matrices are

$$
\tilde{A}=\left[\begin{array}{cc}
A+B_{2} F & B_{2} F \\
0 & A+L C_{2}
\end{array}\right]
$$

and

$$
\tilde{B}=\left[\begin{array}{c}
B_{1} \\
-B_{1}-L D_{21}
\end{array}\right], \quad \tilde{C}=\left[\begin{array}{cc}
C_{1}+D_{12} F & D_{12} F
\end{array}\right]
$$

From the very particular structure of matrix $\tilde{A}$ it is apparent that the closed loop poles are those of matrix $A+B_{2} F$ and $A+L C_{2}$. As far as pole placement is under consideration, in principle it is possible to determine ( $F, L$ ) such that all poles of the closed loop system lie in some region of the complex plane. However, since the state reconstruction from the output depends only on matrix $L$, in our present design procedure it is imposed as the optimal solution of the Output estimation problem, that is $L=L_{2}$. Thus, for a given circular region $\mathcal{R}$, the Output feedback design problem is formulated as

$$
\min \left\{\|T(z, w ; s)\|_{2}^{2}: F \in \mathcal{K}_{R}\right\}
$$

where it is only necessary to make explicit the dependence of $\|T(z, w ; s)\|_{2}^{2}$ with $F$. The optimal solution of this problem is not possible to be determined exactly. So, we proceed by overbounding its objective function. The assumptions i) The pair $\left(A, C_{2}\right)$ is detectable and no eigenvalue of the unreachable part of the pair $\left[\left(A-B_{1} D_{21}^{\prime} C_{2}\right), B_{1}\left(I-D_{21}^{\prime} D_{21}\right)\right]$ lies on the imaginary axis and ii) $D_{12}^{\prime} D_{12}=I$ and $D_{21} D_{21}^{\prime}=I$ are made. Under these assumptions, there exists $\Pi_{2}$ a positive semidefinite and stabilizing solution of the Riccati equation (in the unknown $\Pi$ )

$$
\begin{equation*}
0=\Pi A_{f}^{\prime}+A_{f} \Pi-\Pi C_{2}^{\prime} C_{2} \Pi+B_{1 f} B_{1 f}^{\prime} \tag{6.111}
\end{equation*}
$$

which provides $L_{2}=-\Pi_{2} C_{2}^{\prime}-B_{1} D_{21}^{\prime}$. The controller defined in this way has the following important design property which are obtained as a limit case of Lemma 6.2.

Lemma 6.4 For all $F \in \mathcal{K}_{c}$, the symmetric and positive semidefinite solution to the Lyapunov equation

$$
\begin{equation*}
0=\left(A+B_{2} F\right) X+X\left(A+B_{2} F\right)^{\prime}+L_{2} L_{2}^{\prime} \tag{6.112}
\end{equation*}
$$

is such that the controller with transfer function $K(s)$ given in (6.110) imposes to the closed-loop system the performance

$$
\begin{equation*}
\|T(z, w ; s)\|_{2}^{2}=\operatorname{trace}\left[C_{1} \Pi_{2} C_{1}^{\prime}\right]+\operatorname{trace}\left[\left(C_{1}+D_{12} F\right) X\left(C_{1}+D_{12} F\right)^{\prime}\right] \tag{6.113}
\end{equation*}
$$

Proof The state space representation of the closed-loop system being given by $\Sigma_{F}$, enables us to get immediately (recall Remark 6.14)

$$
\|T(z, w ; s)\|_{2}^{2}=\operatorname{trace}\left[\tilde{C} \tilde{P}_{r} \tilde{C}^{\prime}\right]
$$

where $\tilde{P}_{r}$, solution of the linear equation (6.36) is given by

$$
\tilde{P}_{r}=\left[\begin{array}{cc}
\Pi_{2}+X & -\Pi_{2} \\
-\Pi_{2} & \Pi_{2}
\end{array}\right]
$$

Simple substitution shows that equality (6.113) holds, proving thus the lemma proposed.

This result opens the possibility to reduce the Output feedback design problem to the previous State feedback design problem. In fact, once the equality $L=L_{2}$ holds, the transfer function $K(s)$ of the output feedback controller is completely parametrized by matrix $F$ only. The impact of this decision in the global cost is the first term in (6.113) which does not need to be considered further since it remains constant for all possible choices of $F$. Hence, if the state feedback design problem is applied to the auxiliary plant

$$
\begin{align*}
\dot{x} & =A x+L_{2} D_{21} w+B_{2} u  \tag{6.114}\\
z & =C_{1} x+D_{12} u  \tag{6.115}\\
y & =x \tag{6.116}
\end{align*}
$$

that is, if problem (6.106) is solved with $B_{1}$ replaced by $L_{2} D_{21}$ then the optimal gain $F=W_{2}^{\prime} W_{1}^{-1}$ is such that the controller

$$
K(s):=\left[\begin{array}{c|c}
A+B_{2} W_{2}^{\prime} W_{1}^{-1}+L_{2} C_{2} & -L_{2}  \tag{6.117}\\
\hline W_{2}^{\prime} W_{1}^{-1} & 0
\end{array}\right]
$$

imposes all the eigenvalues of matrix $A+B_{2} W_{2}^{\prime} W_{1}^{-1}$ in the circular region $\mathcal{R}$ while the upper bound

$$
\|T(z, w ; s)\|_{2} \leq \sqrt{\operatorname{trace}\left[C_{1} \Pi_{2} C_{1}^{\prime}\right]+J_{s u b}}
$$

is globally minimized. Unfortunately, as illustrated before, the global minimization of this upper bound does not necessarily means that the global optimum of the true design problem is attained.

### 6.5 Time-domain specifications

One of the most important time-domain specifications of control systems design is the limitation, to some prespecified level, of the time-response overshoot. This section is completely devoted to generalize the previous results to this particular situation. Once again, the important feature is that convexity is preserved and similar manipulations for both the state and output feedback cases are allowed. Some few preliminary calculations are needed before we define and solve the associated optimal control problem.

Consider the system, specified by the following state space minimal realization

$$
\begin{align*}
& \dot{x}=A x+B w, \quad x(0)=0  \tag{6.118}\\
& z=C x \tag{6.119}
\end{align*}
$$

where matrix $A$ is assumed to be stable, the pair $(A, B)$ is reachable and $w \in$ $R L_{2}[0 \infty)$. The quantity we want to compute depends on the transfer function $T(z, w ; s)$ of the above system and is defined as

$$
\begin{equation*}
\mathcal{G}(T(z, w ; s)):=\sup _{\|w\|_{2} \leq 1}\|z\|_{\infty} \tag{6.120}
\end{equation*}
$$

where, for the norm of the output variable $\|z\|_{\infty}$, two different cases are considered, namely

$$
\begin{equation*}
\|z\|_{\infty}:=\sup _{t \geq 0} \sqrt{z^{\prime}(t) z(t)} \tag{6.121}
\end{equation*}
$$

and

$$
\begin{equation*}
\|z\|_{\infty}:=\sup _{t \geq 0} \max _{i}\left|z_{i}(t)\right| \tag{6.122}
\end{equation*}
$$

where $z_{i}(t)$ denotes the $i$-th scalar component of $z(t)$. We are now in position to interpret the function $\mathcal{G}(\cdot)$ defined above. Suppose, for a certain positive scalar $\alpha$,

$$
\mathcal{G}(T(z, w ; s)) \leq \alpha
$$

then for each time $t \geq 0$ the worst case overshoot of $z(t)$ is limited by $\alpha$. Hence, the possibility to take into account this time constraint in control system design is of great practical interest.

Lemma 6.5 Let $P$ be the symmetric and positive definite solution of the Lyapunov equation

$$
\begin{equation*}
0=A P+P A^{\prime}+B B^{\prime} \tag{6.123}
\end{equation*}
$$

The following are true
a) For the norm (6.121) then $\mathcal{G}(T(z, w ; s))=\lambda_{\text {max }}^{1 / 2}\left(C P C^{\prime}\right)$ where $\lambda_{\max }(\cdot)$ denotes the maximum eigenvalue of $(\cdot)$.
b) For the norm (6.122) then $\mathcal{G}(T(z, w ; s))=d_{\max }^{1 / 2}\left(C P C^{\prime}\right)$ where $d_{\max }(\cdot)$ denotes the maximum diagonal element of (.).

Proof Define $v(x):=x^{\prime} P^{-1} x$ and consider the system (6.118) - (6.119) with an arbitrary input $w$ satisfying $\|w\|_{2} \leq 1$. The time derivative of $v(\cdot)$ along a trajectory of that system yields

$$
\begin{aligned}
\dot{v}(x) & =x^{\prime}\left(A^{\prime} P^{-1}+P^{-1} A\right) x+2 w^{\prime} B^{\prime} P^{-1} x \\
& =-x^{\prime} P^{-1} B B^{\prime} P^{-1} x+2 w^{\prime} B^{\prime} P^{-1} x \\
& =w^{\prime} w-\left(w-B^{\prime} P^{-1} x\right)^{\prime}\left(w-B^{\prime} P^{-1} x\right) \\
& \leq\|w\|^{2}
\end{aligned}
$$

which, after integration of both sides from 0 to $t \geq 0$ provides

$$
v(x(t))=x^{\prime}(t) P^{-1} x(t) \leq\|w\|_{2}^{2} \leq 1
$$

This inequality means that, in the state space, the trajectories $x(t)$, for all $t \geq 0$ are confined in the set $x^{\prime} P^{-1} x \leq 1$ whenever $w$ remains bounded by $\|w\|_{2} \leq 1$.

Point a) With $z=C x$ and $\tilde{x}:=P^{-1 / 2} x$, we have

$$
\begin{aligned}
\mathcal{G}(T(z, w ; s))^{2} & =\sup _{\|w\|_{2} \leq 1}\|z\|_{\infty}^{2} \\
& \leq \max \left\{x^{\prime} C^{\prime} C x ; x^{\prime} P^{-1} x \leq 1\right\} \\
& \leq \max \left\{\tilde{x}^{\prime} P^{1 / 2} C^{\prime} C P^{1 / 2} \tilde{x} ; \tilde{x}^{\prime} \tilde{x} \leq 1\right\} \\
& \leq \lambda_{\max }\left(P^{1 / 2} C^{\prime} C P^{1 / 2}\right) \\
& \leq \lambda_{\max }\left(C P C^{\prime}\right)
\end{aligned}
$$

and it remains to show that there exists a feasible trajectory $w(t)$ such that $\|z\|_{\infty}^{2}$ is arbitrarily close to $\lambda_{\max }\left(C P C^{\prime}\right)$. To this end, consider $T>0$ fixed but arbitrary

$$
0<S(T):=\int_{0}^{T} e^{A t} B B^{\prime} e^{A^{\prime} t} d t \leq P
$$

and the input signal such that $w(t)=0$ for all $t>T$ and

$$
w(t)=B^{\prime} e^{A^{\prime}(T-t)} S(T)^{-1 / 2} \psi, \quad 0 \leq t \leq T
$$

where $\psi$ is a vector to be determined. Simple calculations show that

$$
\|w\|_{2}^{2}=\int_{0}^{\infty} w^{\prime}(t) w(t) d t=\int_{0}^{T} w^{\prime}(t) w(t) d t=\psi^{\prime} \psi
$$

and

$$
z(T)=C \int_{0}^{T} e^{A(T-\tau)} B w(\tau) d \tau=C S(T)^{1 / 2} \psi
$$

Consequently, choosing $\psi$ as being the unitary norm eigenvector associated to the maximum eigenvalue of matrix $S(T)^{1 / 2} C^{\prime} C S(T)^{1 / 2}$, the feasibility of the input signal $w$ is guaranteed and

$$
\begin{aligned}
\|z\|_{\infty}^{2} & =\sup _{t \geq 0} z^{\prime}(t) z(t) \\
& \geq \psi^{\prime} S(T)^{1 / 2} C^{\prime} C S(T)^{1 / 2} \psi \\
& \geq \lambda_{\max }\left(C S(T) C^{\prime}\right)
\end{aligned}
$$

the proof is then concluded because $S(T)$ becomes arbitrarily close to $P$ as $T$ increases.
Point b) With $z_{i}=C_{i} x$ with $C_{i}$ being the $i$-th row of matrix $C$ and $\tilde{x}:=P^{-1 / 2} x$, we have

$$
\begin{aligned}
\mathcal{G}(T(z, w ; s))^{2} & =\sup _{\|w\|_{2} \leq 1}\|z\|_{\infty}^{2} \\
& \leq \max _{x, i}\left\{x^{\prime} C_{i}^{\prime} C_{i} x ; x^{\prime} P^{-1} x \leq 1\right\} \\
& \leq \max _{i} \max _{\tilde{x}}\left\{\tilde{x}^{\prime} P^{1 / 2} C_{i}^{\prime} C_{i} P^{1 / 2} \tilde{x} ; \tilde{x}^{\prime} \tilde{x} \leq 1\right\} \\
& \leq \max _{i} \lambda_{\max }\left(P^{1 / 2} C_{i}^{\prime} C_{i} P^{1 / 2}\right) \\
& \leq \max _{i} C_{i} P C_{i}^{\prime} \\
& \leq d_{\max }\left(C P C^{\prime}\right)
\end{aligned}
$$

As before, it remains to determine a feasible input such that the equality holds. This is accomplished by the same function $w(t)$ already defined and a convenient choice of vector $\psi$. Indeed, take $\psi_{i}$ as being the unitary norm eigenvector associated to the maximum eigenvalue of matrix $S(T)^{1 / 2} C_{i}^{\prime} C_{i} S(T)^{1 / 2}$, and choose $\psi=\psi_{l}$ where $C_{i} S(T) C_{i}^{\prime} \leq C_{l} S(T) C_{l}^{\prime}$ for all index $i=1,2, \cdots$, then $w$ is feasible and

$$
\begin{aligned}
\|z\|_{\infty}^{2} & =\sup _{t \geq 0} \max _{i} z_{i}^{\prime}(t) z_{i}(t) \\
& \geq \max _{i} \psi^{\prime} S(T)^{1 / 2} C_{i}^{\prime} C_{i} S(T)^{1 / 2} \psi \\
& \geq C_{l} S(T) C_{l}^{\prime} \\
& \geq d_{\max }\left(C S(T) C^{\prime}\right)
\end{aligned}
$$

the proof is then concluded since as said before, $S(T)$ becomes arbitrarily close to $P$ as $T$ increases.

Remark 6.33 In the proof of Lemma 6.5, it is assumed that the solution of the linear matrix equation (6.123) is positive definite. This occurs whenever the pair $(A, B)$ is reachable. If this assumption is not verified, the result still holds true. In this case, using Kalman's canonical controllability form, it is immediate to see that the output $z(t)$ in (6.119) depends only on the reachable part of the system.

Remark 6.34 The relationship between both $\mathcal{G}(\cdot)$ for the norms (6.121) and (6.122) are

$$
d_{\max }\left(C P C^{\prime}\right) \leq \lambda_{\max }\left(C P C^{\prime}\right) \leq \operatorname{trace}\left(C P C^{\prime}\right)
$$

which also implies, in both cases, that $\mathcal{G}(T(z, w ; s)) \leq\|T(z, w ; s)\|_{2}$.
Remark 6.35 (Convexity) The real valued function $g(X): \mathcal{P} \longrightarrow R$ defined as $g(X):=$ $\lambda_{\max }(X)$ is convex (recall Remark 6.3). The same is true for the function $g(X):=d_{\max }(X)$. To prove this, take $X_{0} \in \mathcal{P}$ and $e_{0}$ the column of the identity matrix such that $e_{0}^{\prime} X_{0} e_{0}=$ $g\left(X_{0}\right)$. For all $X \in \mathcal{P}$ we get

$$
\begin{aligned}
g(X) & =d_{\max }(X) \\
& \geq e_{0}^{\prime} X e_{0} \\
& \geq g\left(X_{0}\right)+<e_{0} e_{0}^{\prime}, X-X_{0}>
\end{aligned}
$$

and inequality (6.1) is verified for $\Lambda_{0}=e_{0} e_{0}^{\prime}$, then convexity follows. Notice further that both functions are not differentiable in $\mathcal{P}$ and are non decreasing functions in the sense that for any $X_{1}, X_{2} \in \mathcal{P}$ such that $X_{1} \leq X_{2}$ then $g\left(X_{1}\right) \leq g\left(X_{2}\right)$.

Example 6.5 Consider the system (6.118) - (6.119) with

$$
A=\left[\begin{array}{rrr}
0 & 1 & 0 \\
0 & 0 & 1 \\
-10 & -9 & -4
\end{array}\right], \quad B=\left[\begin{array}{r}
2 \\
-2 \\
4
\end{array}\right], \quad C=\left[\begin{array}{lll}
0 & 1 & 0 \\
0 & 0 & 1
\end{array}\right]
$$

Our purpose is to illustrate the result of Lemma 6.5. The definite positive solution $P$ of the Lyapunov equation (6.123), provides

$$
\mathcal{G}(T(z, w ; s))=\left\{\begin{array}{l}
\lambda_{\max }^{1 / 2}\left(C P C^{\prime}\right)=2.16 \\
d_{\max }^{1 / 2}\left(C P C^{\prime}\right)=1.68
\end{array}\right.
$$

On the other hand, taking $T=5$ it can be verified that $S(T) \approx P$ and so, corresponding to the input

$$
w(t)=\left\{\begin{array}{lr}
B^{\prime} e^{A^{\prime}(T-t)} P^{-1 / 2} \psi & , \quad 0 \leq t \leq T \\
0 & , \quad t>T
\end{array}\right.
$$



Figure 6.12: Norms of the output $z(t)$
the system produces the output $z(t)$ which enables us to calculate the time varying function

$$
f(t):=\left\{\begin{array}{c}
\sqrt{z_{1}(t)^{2}+z_{2}(t)^{2}} \\
\max \left\{\left|z_{1}(t)\right|,\left|z_{2}(t)\right|\right\}
\end{array}\right.
$$

for each norm used to define $\mathcal{G}(\cdot)$. These functions are shown in fig 6.12 where the labels $\lambda_{\max }$ and $d_{\text {max }}$ identify the norms introduced in (6.121) and (6.122) respectively. In fig. 6.13 the corresponding inputs $w(t)$ are also shown. It is interesting to verify that in both cases

$$
\max _{t \geq 0} f(t) \approx \mathcal{G}(T(z, w ; s))
$$

and $\|w\|_{2} \approx 1$ as required in the proof of Lemma 6.5.
Throughout the remaining of this section, we define the convex function

$$
\begin{equation*}
g\left(C P C^{\prime}\right):=\mathcal{G}(T(z, w ; s))^{2} \tag{6.124}
\end{equation*}
$$

to indicate both cases treated before. Since this function is convex in the domain $\mathcal{P}$, the constraint $g\left(C P C^{\prime}\right) \leq \alpha$, for $\alpha>0$ fixed, can be handled with no additional difficulty because convexity is preserved. The same obviously occurs if $g\left(C P C^{\prime}\right)$ is used as a objective function to be minimized. This case, generalizes the control design problem in $R \mathrm{H}_{2}$ in the sense that it is obtained from the above formulation for the choice $g\left(C P C^{\prime}\right)=\operatorname{trace}\left(C P C^{\prime}\right)$.

### 6.5.1 State feedback design

The dynamic system is described by the equations

$$
\begin{align*}
& \dot{x}=A x+B_{1} w+B_{2} u  \tag{6.125}\\
& z=C_{1} x+D_{12} u  \tag{6.126}\\
& y=x \tag{6.127}
\end{align*}
$$



Figure 6.13: The inputs $w(t)$
where our purpose is to determine the optimal matrix $F$ such that with $u=F x$ the following control problem is solved

$$
\begin{equation*}
\min \left\{\mathcal{G}(T(z, w ; s)): F \in \mathcal{K}_{c}\right\} \tag{6.128}
\end{equation*}
$$

To this end we only consider the assumption i) $D_{12}^{\prime} D_{12}=I$. For each $F \in \mathcal{K}_{c}$ then $T(z, w ; s)$ is in $R H_{2}$ which means that if $\mathcal{K}_{c} \neq \emptyset$ then (6.128) is well-posed. The main feature of the result to be presented is that the optimal global solution of (6.128) is actually attained.

Theorem 6.10 Consider $B_{1} B_{1}^{\prime}>0$ and define the matrix

$$
\bar{R}_{c}:=\left[\begin{array}{ll}
C_{1} & D_{12} \tag{6.129}
\end{array}\right]
$$

The convex programming problem

$$
\begin{equation*}
J_{o p t}:=\min \left\{g\left(\bar{R}_{c} W \bar{R}_{c}^{\prime}\right): W \in \mathcal{C}_{c}\right\} \tag{6.130}
\end{equation*}
$$

is equivalent to problem (6.128) in the sense that both present the same global solution.
Proof Recall Theorem 6.1 where it is proved that all elements of the set $\mathcal{K}_{c}$ are generated from those of $\mathcal{C}_{c}$. Hence, as far as feasibility is concerned, both problems are equivalent. We proceed by assuming that problem (6.128) has an optimal solution. In this case, from Lemma 6.5 there exist matrices $\bar{P}>0$ and $\bar{F} \in \mathcal{K}_{c}$ such that

$$
g\left[\left(C_{1}+D_{12} \bar{F}\right) \bar{P}\left(C_{1}+D_{12} \bar{F}\right)^{\prime}\right]=\min \mathcal{G}(T(z, w ; s))
$$

and

$$
0=\left(A+B_{2} \bar{F}\right) \bar{P}+\bar{P}\left(A+B_{2} \bar{F}\right)^{\prime}+B_{1} B_{1}^{\prime}
$$

Defining the matrix

$$
\bar{W}:=\left[\begin{array}{cc}
\bar{P} & \bar{P} \bar{F}^{\prime} \\
\bar{F} \bar{P} & \bar{F} \bar{P} \bar{F}^{\prime}
\end{array}\right]
$$

simple algebraic manipulations yield the conclusion that $\bar{W} \in \mathcal{C}_{c}$ provides the minimum cost. It remains to prove that $\bar{W}$ is the global optimum of problem (6.130). To this end, taking any feasible $W$ we have $F=W_{2} W_{1}^{-1} \in \mathcal{K}_{c}$ and

$$
\mathcal{G}(T(z, w ; s))=g\left[\left(C_{1}+D_{12} W_{2}^{\prime} W_{1}^{-1}\right) P_{r}\left(C_{1}+D_{12} W_{2}^{\prime} W_{1}^{-1}\right)\right]
$$

where $W_{1} \geq P_{r}$. Using again the fact that $W \geq 0$ together with the Schur complement formula we get

$$
\begin{aligned}
g\left(\bar{R}_{c} \bar{W} \bar{R}_{c}^{\prime}\right) & =\min \mathcal{G}(T(z, w ; s)) \\
& \leq \mathcal{G}(T(z, w ; s)) \\
& \leq g\left[\left(C_{1}+D_{12} W_{2}^{\prime} W_{1}^{-1}\right) W_{1}\left(C_{1}+D_{12} W_{2}^{\prime} W_{1}^{-1}\right)\right] \\
& \leq g\left[\bar{R}_{c} W \bar{R}_{c}^{\prime}-D_{12}\left(W_{3}-W_{2}^{\prime} W_{1}^{-1} W_{2}\right) D_{12}^{\prime}\right] \\
& \leq g\left[\bar{R}_{c} W \bar{R}_{c}^{\prime}\right]
\end{aligned}
$$

which being true for all $W \in \mathcal{C}_{c}$ completes the proof.
The proof of this theorem is almost the same of that of Theorem 6.3 where the special case $g(\cdot)=$ trace $(\cdot)$ has been considered. Notice further that the above proof depends basically on the convexity and on the non decreasing property of the function $g(\cdot)$ introduced in Remark 6.35. Consequently, the same result also applies to any other function with these properties.

### 6.5.2 Output feedback design

Once again consider the system

$$
\begin{align*}
& \dot{x}=A x+B_{1} w+B_{2} u  \tag{6.131}\\
& z=C_{1} x+D_{12} u  \tag{6.132}\\
& y=C_{2} x+D_{21} w \tag{6.133}
\end{align*}
$$

where the transfer function $K(s)$ of the output controller is to be determined. It is adopted the same reasoning as before, that is all feasible controllers are completely parametrized by only one matrix $F$ which is used to meet the design requirements. Its transfer function is given by

$$
K(s):=\left[\begin{array}{c|c}
A+B_{2} F+L_{2} C_{2} & -L_{2}  \tag{6.134}\\
\hline F & 0
\end{array}\right]
$$

being thus apparent that the choice $L=L_{2}$ is the best we can do as far as the reconstruction of the state from the output is concerned. Let us keep in mind that $L_{2}=-\Pi_{2} C_{2}^{\prime}-B_{1} D_{21}^{\prime}$, where $\Pi_{2}$ is the positive semidefinite and stabilizing solution of the Riccati equation (in the unknown $\Pi$ )

$$
\begin{equation*}
0=\Pi A_{f}^{\prime}+A_{f} \Pi-\Pi C_{2}^{\prime} C_{2} \Pi+B_{1 f} B_{1 f}^{\prime} \tag{6.135}
\end{equation*}
$$

The feedback connection drawn in fig. 6.5 puts in evidence that the internal stability of the closed-loop system is assured if $F \in \mathcal{K}_{c}$ and so our goal is to solve the associated optimal control problem

$$
\min \left\{\mathcal{G}(T(z, w ; s)): F \in \mathcal{K}_{c}\right\}
$$

The next lemma provides the generalization of the result introduced in Lemma 6.4 to deal with its objective function.

Lemma 6.6 For all $F \in \mathcal{K}_{c}$ the symmetric and positive semidefinite solution to the Lyapunov equation

$$
\begin{equation*}
0=\left(A+B_{2} F\right) X+X\left(A+B_{2} F\right)^{\prime}+L_{2} L_{2}^{\prime} \tag{6.136}
\end{equation*}
$$

is such that the controller with transfer function $K(s)$ given in (6.134) imposes to the closed-loop system the performance

$$
\begin{equation*}
\mathcal{G}(T(z, w ; s))=g\left[C_{1} \Pi_{2} C_{1}^{\prime}+\left(C_{1}+D_{12} F\right) X\left(C_{1}+D_{12} F\right)^{\prime}\right] \tag{6.137}
\end{equation*}
$$

Proof From Lemma 6.5 together with the state space realization of the closed-loop system, yields

$$
\mathcal{G}(T(z, w ; s))=g\left[\tilde{C} \tilde{P}_{r} \tilde{C}^{\prime}\right]
$$

where $\tilde{P}_{r}$, solution of the linear equation (6.36) is given by

$$
\tilde{P}_{r}=\left[\begin{array}{cc}
\Pi_{2}+X & -\Pi_{2} \\
-\Pi_{2} & \Pi_{2}
\end{array}\right]
$$

Furthermore, simple calculations enables us to write

$$
\begin{equation*}
\tilde{C} \tilde{P}_{r} \tilde{C}^{\prime}=C_{1} \Pi_{2} C_{1}^{\prime}+\left(C_{1}+D_{12} F\right) X\left(C_{1}+D_{12} F\right)^{\prime} \tag{6.138}
\end{equation*}
$$

which proves the lemma proposed.
The equality (6.138) is of particular importance. With it, the specific properties of the function $g(\cdot)$ is not used in the proof of the above lemma. Consequently, the same result also holds for any other function $g(\cdot)$. For those functions under consideration in this section, the Output feedback problem is reduced to the State feedback problem applied to the auriliary plant

$$
\begin{align*}
& \dot{x}=A x+L_{2} D_{21} w+B_{2} u  \tag{6.139}\\
& z=C_{1} x+D_{12} u  \tag{6.140}\\
& y=x \tag{6.141}
\end{align*}
$$

Doing this, it is important to keep in mind that both, matrix $B_{1}$ should be replaced by $L_{2} D_{21}$ and the objective function to be minimized over $\mathcal{C}_{c}$ should be replaced accordingly, leading to

$$
\begin{equation*}
J_{o p t}=\min \left\{g\left[C_{1} \Pi_{2} C_{1}^{\prime}+\bar{R}_{c} W \bar{R}_{c}^{\prime}\right]: W \in \mathcal{C}_{c}\right\} \tag{6.142}
\end{equation*}
$$

which again is a convex programming problem.
Remark 6.36 Due to the fact that

$$
g\left[C_{1} \Pi_{2} C_{1}^{\prime}+\bar{R}_{c} W \bar{R}_{c}^{\prime}\right] \leq g\left[C_{1} \Pi_{2} C_{1}^{\prime}\right]+g\left[\bar{R}_{c} W \bar{R}_{c}^{\prime}\right]
$$

it is possible to simplify the objective function of Problem (6.142) by retaining in the optimization problem, only the second term in the above expression (since the first one is constant). Unfortunately, this may produces an important degree of sub-optimality because the overbound is, in general, very conservative.

Example 6.6 Consider the system (6.139) - (6.141) with the matrices as indicated below

$$
\begin{gathered}
A=\left[\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right], \quad B_{1}=\left[\begin{array}{ll}
1 & 0 \\
1 & 0
\end{array}\right], \quad B_{2}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad D_{12}=\left[\begin{array}{l}
0 \\
1
\end{array}\right] \\
C_{1}=\left[\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right], \quad C_{2}=\left[\begin{array}{ll}
1 & 0
\end{array}\right], \quad D_{21}=\left[\begin{array}{ll}
0 & 1
\end{array}\right]
\end{gathered}
$$

Solving the Riccati equation (6.135) we get $L_{2}=-[1.73211]^{\prime}$ and

$$
C_{1} \Pi_{2} C_{1}^{\prime}=\left[\begin{array}{cc}
1.7321 & 0 \\
0 & 0
\end{array}\right]
$$

Then we solved problem (6.142) considering $g(\cdot)=$ trace $(\cdot), g(\cdot)=\lambda_{\max }(\cdot)$ and $g(\cdot)=$ $d_{\text {max }}(\cdot)$. The optimal controllers and the minimum cost associated, for each case are respectively

$$
\begin{aligned}
& K_{o p t}(s)=\frac{-3.14 s-0.99}{s^{2}+3.14 s+4.44}, \quad J_{o p t}=10.85 \\
& K_{o p t}(s)=\frac{-4.59 s-1.32}{s^{2}+4.03 s+6.31}, \quad J_{o p t}=8.59 \\
& K_{o p t}(s)=\frac{-4.26 s-1.47}{s^{2}+3.44 s+5.44}, \quad J_{o p t}=5.73
\end{aligned}
$$

Each controller and the associated cost are quite different which indicates that the optimum of each problem solved are distinct, even though they obey the inequality given in Remark 6.34. The optimal solution of the upper bound of all $g(\cdot)$ namely trace $(\cdot)$ may furnish a poor suboptimal solution to the other cases under consideration.

Following the discussion in Remark 6.36, we also calculated the following suboptimal cost

$$
J_{s u b}=g\left(C_{1} \Pi_{2} C_{1}^{\prime}\right)+\min \left\{g\left[\bar{R}_{c} W \bar{R}_{c}^{\prime}\right]: W \in \mathcal{C}_{c}\right\}
$$

for each function $g(\cdot)$ as before. The result is

$$
\begin{aligned}
& K_{s u b}(s)=\frac{-3.14 s-0.99}{s^{2}+3.14 s+4.44}, \quad J_{s u b}=10.85 \\
& K_{s u b}(s)=\frac{-4.16 s-1.16}{s^{2}+3.88 s+5.88}, \quad J_{s u b}=9.27 \\
& K_{s u b}(s)=\frac{-3.60 s-1.19}{s^{2}+3.27 s+4.86}, \quad J_{s u b}=6.35
\end{aligned}
$$

Comparing these controllers with the optimal ones, it can be verified an important loss on the performance index. This occurs in all cases but the first one on which the trace function is used.

### 6.6 Controllers with structural constraints

In practice one is frequently faced to control design problems where the controller must exhibit some desired structure. For instance, if the system to be controlled is composed by many coupled subsystems, robustness considerations requires the controller should use only local informations for feedback, that is, in the state feedback case, the matrix gain $F$ must present a block-diagonal structure. This particular structure defines the important class of Decentralized control design problems. As well as, if only the system output is available for feedback we already know how to design a dynamic controller to meet certain performance criteria. However, it is also of great interest to determine (if any) a static feedback gain with the same paradigm. This is on the origin of the so called Static output control design problems. This section is entirely devoted to analyze these problems in the framework of convex analysis.


Figure 6.14: The interconnected System

### 6.6.1 Decentralized control design

The system to be dealt with is represented in Fig. 6.14. It is composed by a number $N$ of subsystems each of them with transfer function $P_{i}(s)$, coupled together by means of an interconnection structure. The dynamic model of the $i$-th subsystem with transfer function $P_{i}(s)$ is given by

$$
\begin{align*}
\dot{x}_{i} & =A^{i i} x_{i}+B_{1}^{i} w+B_{2}^{i} u_{i}+\sum_{j \neq i}^{N} A^{i j} x_{j}  \tag{6.143}\\
z_{i} & =C_{1}^{i} x_{i}+D_{12}^{i} u_{i}  \tag{6.144}\\
y_{i} & =x_{i} \tag{6.145}
\end{align*}
$$

In order to impose to the overall system a desirable performance, it is asked to design $N$ local controllers, each of then using only the information available in the local state variable $x_{i}, i=1,2, \cdots, N$. That is,

$$
\begin{equation*}
u_{i}=F_{i} x_{i}, \quad i=1,2, \cdots, N \tag{6.146}
\end{equation*}
$$

Clearly, if the overall system is rewritten in the standard form

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{6.147}\\
z & =C_{1} x+D_{12} u  \tag{6.148}\\
y & =x \tag{6.149}
\end{align*}
$$

then the control is given by

$$
\begin{equation*}
u=F x, \quad F=\operatorname{blockdiag}\left[F_{1}, F_{2}, \cdots, F_{N}\right] \tag{6.150}
\end{equation*}
$$

where each block of $F$ is the local matrix gain with appropriate dimension. To ease the presentation, let us introduce the following notation. The subscript " $D$ " in a matrix, for instance $F_{D}$ means that this matrix is constraint to have a block-diagonal structure. In other words, $F_{D}$ is obtained from any $F$ by simply zeroing all off blockdiagonal elements. The problem to be faced in this section is how to incorporate the above structural constraint in several design procedures of interest. It is important to keep in mind that the results of this section are related to the general linear system (6.147)-(6.149) and do not depend upon any particular system structure like that of system (6.143)-(6.145).

The first important feature to be analyzed is the internal stability of the closedloop system. Clearly, internal stability is guaranteed whenever there exist $F$ such that

$$
F=F_{D} \in \mathcal{K}_{c}
$$

Furthermore, all decentralized matrices $F_{D}$ with this property constitute the set of all stabilizing decentralized gains for the system under consideration. This set is not convex and in many cases may be constituted by disjoint subsets. So, we work here with a particular subset of it defined as follows

Definition 6.3 (Structural D - Stabilizability) The pair $\left(A, B_{2}\right)$ is said to be structurally $D$-stabilizable if there exist matrices $P_{D}$ symmetric and positive definite and $F_{D}$ such that

$$
\begin{equation*}
0 \geq\left(A+B_{2} F_{D}\right) P_{D}+P_{D}\left(A+B_{2} F_{D}\right)^{\prime}+Q \tag{6.151}
\end{equation*}
$$

for some matrix $Q=Q^{\prime}>0$. The set of all such matrices $F_{D}$ is denoted $\mathcal{K}_{D}$.
We want to stress that $F_{D} \in \mathcal{K}_{D}$ implies that $F_{D} \in \mathcal{K}_{c}$ but the converse is not necessarily true. Moreover, it is possible to have $\mathcal{K}_{D}=\emptyset$ while there exists $F_{D} \in \mathcal{K}_{c}$. The additional constraint in Definition 6.3 is that the Lyapunov inequality (6.151) must present a block-diagonal solution $P=P_{D}$. At a first glance, it may appears that the existence of a pair $\left(F_{D}, P_{D}\right)$ depends on a particular choice of matrix $Q$. Fortunately this is not true as can be simply demonstrated as follows. For a given matrix $F_{D}$, suppose $\bar{P}_{D}$ satisfies the Lyapunov inequality (6.151) with $Q=\bar{Q}>0$. For any other matrix $Q>0$, choosing the scalar $\beta>0$ such that $\beta \bar{Q} \geq Q$, it is simple to verify that (6.151) is also satisfied for $P_{D}=\beta \bar{P}_{D}>0$. Hence, matrix $Q>0$ in Definition 6.3 can be chosen arbitrarily. Before presenting the next result we introduce the notation used for the partitioned matrix $W$. The subscript " D " is used as follows

$$
W_{D}:=\left[\begin{array}{cc}
W_{1 D} & W_{2 D} \\
W_{2 D}^{\prime} & W_{3}
\end{array}\right]
$$

which indicates that only the sub-blocks $W_{1}$ and $W_{2}$ have to present the decentralized structure, namely

$$
\begin{aligned}
& W_{1 D}=\operatorname{blockdiag}\left[W_{11}, W_{12}, \cdots, W_{1 N}\right] \\
& W_{2 D}=\operatorname{blockdiag}\left[W_{21}, W_{22}, \cdots, W_{2 N}\right]
\end{aligned}
$$

with $W_{1 i}$ and $W_{2 i}$ being $n_{i} \times n_{i}$ and $n_{i} \times m_{i}$ matrices where $n_{i}$ is the local state vector dimension and $m_{i}$ is the local control vector dimension respectively for all $i=1,2, \cdots, N$.

Theorem 6.11 Consider $B_{1} B_{1}^{\prime}>0$ and define the set

$$
\begin{equation*}
\mathcal{C}_{D}:=\left\{W: W=W_{D}\right\} \cap \mathcal{C}_{c} \tag{6.152}
\end{equation*}
$$

The following are true
a) $\mathcal{C}_{D}$ is a convex set.
b) Each $W \in \mathcal{C}_{D}$ is such that $W_{1}>0$.
c) $\mathcal{K}_{D}=\left\{W_{2}^{\prime} W_{1}^{-1}: W \in \mathcal{C}_{D}\right\}$.

Proof Since the constraint $W=W_{D}$ is linear and the set $\mathcal{C}_{c}$ is convex then point a) is proved. Moreover, the proof of point b) is a consequence of $B_{1} B_{1}^{\prime}>0$ together with $\mathcal{C}_{D} \subset \mathcal{C}_{c}$.

Point c) For the necessity, assume $F_{D} \in \mathcal{K}_{D} \neq \emptyset$. Setting $Q=B_{1} B_{1}^{\prime}$ in (6.151), it is simple to verify that

$$
W=W_{D}=\left[\begin{array}{cc}
P_{D} & P_{D} F_{D}^{\prime} \\
F_{D} P_{D} & F_{D} P_{D} F_{D}^{\prime}
\end{array}\right] \in \mathcal{C}_{c}
$$

which is the same to say that $W \in \mathcal{C}_{D}$. Furthermore, $W_{2}^{\prime} W_{1}^{-1}=F_{D} P_{D} P_{D}^{-1}=F_{D}$. Conversely, for any $W \in \mathcal{C}_{D} \neq \emptyset$ and any $v \in \mathcal{N}_{c}$ we get

$$
\begin{aligned}
0 & \geq v^{\prime} \Theta_{c}(W) v \\
& \geq x^{\prime}\left[\left(A+B_{2} W_{2}^{\prime} W_{1}^{-1}\right) W_{1}+W_{1}\left(A+B_{2} W_{2}^{\prime} W_{1}^{-1}\right)^{\prime}+B_{1} B_{1}^{\prime}\right] x
\end{aligned}
$$

which shows that inequality (6.151) holds for $P_{D}=W_{1}=W_{1 D}$ and $F_{D}=W_{2}^{\prime} W_{1}^{-1}=$ $W_{2 D}^{\prime} W_{1 D}^{-1}$, that is both matrices exhibit the desired block-diagonal structure. From the above, the equality in point c) also holds when $\mathcal{K}_{D}=\emptyset$ or $\mathcal{C}_{D}=\emptyset$.

The linear constraint $W=W_{D}$ is essential to get this result. It provides all stabilizing decentralized matrices $F=F_{D}$ and the quadratic function

$$
v(x)=x^{\prime} W_{1 D}^{-1} x=\sum_{i=1}^{N} x_{i}^{\prime} W_{1 i}^{-1} x_{i}
$$

is a Lyapunov function associated to the closed-loop system with $w=0$ since, its time-derivative along an arbitrary trajectory of the system is

$$
\dot{v}(x)=-\left\|B_{1}^{\prime} W_{1 D}^{-1} x\right\|^{2}<0, \quad \forall x \neq 0
$$

The interpretation of Theorem 6.11 is now clear. It generates all stabilizing matrices $F_{D}$ such that the closed-loop system stability is tested by an additively separable Lyapunov function.

In this framework, it is possible to solve an approximate version of the decentralized state feedback design problem in $R H_{2}$, written in the form

$$
\begin{equation*}
\min \left\{\|T(z, w ; s)\|_{2}^{2}: F \in \mathcal{K}_{D}\right\} \tag{6.153}
\end{equation*}
$$

which makes once again explicit that the feasible set is restricted to those gains satisfying Definition 6.3.

Theorem 6.12 Assume $B_{1} B_{1}^{\prime}>0$ and let $\bar{W}$ be the optimal solution of the convex programming problem

$$
\begin{equation*}
J_{s u b}:=\min \left\{\operatorname{trace}\left[R_{c} W\right]: W \in \mathcal{C}_{D}\right\} \tag{6.154}
\end{equation*}
$$

Then, $F=F_{D}=\bar{W}_{2}^{\prime} \bar{W}_{1}^{-1} \in \mathcal{K}_{D}$ minimizes an upper bound of the objective function of Problem ( 6.153 ) in the sense that $\|T(z, w ; s)\|_{2}^{2} \leq \operatorname{trace}\left[R_{c} W\right]$, for all $W \in \mathcal{C}_{D}$.

Proof The infeasibility of one problem implies the same is true to the other and vice versa. Hence, assuming they are feasible, that the optimal solution of Problem
(6.154) provides $\bar{W}_{2}^{\prime} \bar{W}_{1}^{-1} \in \mathcal{K}_{D}$ is a consequence of Theorem 6.11. Additionally, for any $W \in \mathcal{C}_{D}$ we have $W_{1 D} \geq P_{r}$ where

$$
0=\left(A+B_{2} W_{2 D}^{\prime} W_{1 D}^{-1}\right) P_{r}+P_{r}\left(A+B_{2} W_{2 D}^{\prime} W_{1 D}^{-1}\right)^{\prime}+B_{1} B_{1}^{\prime}
$$

yielding

$$
\begin{aligned}
\operatorname{trace}\left[R_{c} W\right]= & \operatorname{trace}\left[\left(C_{1}+D_{12} W_{2 D}^{\prime} W_{1 D}^{-1}\right) W_{1 D}\left(C_{1}+D_{12} W_{2 D}^{\prime} W_{1 D}^{-1}\right)^{\prime}\right]+ \\
& +\operatorname{trace}\left[D_{12}\left(W_{3}-W_{2 D}^{\prime} W_{1 D}^{-1} W_{2 D}\right) D_{12}^{\prime}\right] \\
\geq & \operatorname{trace}\left[\left(C_{1}+D_{12} W_{2 D}^{\prime} W_{1 D}^{-1}\right) P_{r}\left(C_{1}+D_{12} W_{2 D}^{\prime} W_{1 D}^{-1}\right)^{\prime}\right] \\
\geq & \|T(z, w ; s)\|_{2}^{2}
\end{aligned}
$$

completing the proof.
At this point it is important to keep in mind why Problem (6.154) corresponds to minimize only an upper bound to the objective function of Problem (6.153). The reason is that even though $W \in \mathcal{C}_{D}$ generates all feasible gains $F=F_{D}$, the matrix $P_{r}$ used to determine the corresponding value of $\|T(z, w ; s)\|_{2}^{2}$ does not necessarily satisfies the decentralized constraint $P_{r}=\left(P_{r}\right)_{D}$. Generally the inequality $W_{1 D} \geq P_{r}$ is only strictly satisfied for all $W \in \mathcal{C}_{D}$.

Remark 6.37 The exact Decentralized state feedback design in $R H_{2}$ is the optimal control problem

$$
\min \left\{\|T(z, w ; s)\|_{2}^{2}: F=F_{D} \in \mathcal{K}_{c}\right\}
$$

Using Theorem 6.1, this is equivalent to

$$
\min \left\{\operatorname{trace}\left[R_{c} W\right]: W \in \mathcal{C}_{c D}\right\}
$$

where

$$
\mathcal{C}_{c D}:=\left\{W: W_{2}^{\prime} W_{1}^{-1}=\left(W_{2}^{\prime} W_{1}^{-1}\right)_{D}\right\} \cap \mathcal{C}_{c}
$$

Unfortunately, the nonlinear equality constraint present in the set $\mathcal{C}_{c D}$ makes it nonconvex. The way to circumvent this difficulty is to impose the decentralized structure on matrix $W_{2}$ and $W_{1}$ simultaneously, as required by Definition 6.3.

Remark 6.38 The convex constraint $W=W_{D}$ can be added to any other design problem in order to search a decentralized stabilizing control. For instance, any $W$ belonging to the convex set,

$$
\mathcal{C}_{\gamma D}:=\left\{W: W=W_{D}\right\} \cap \mathcal{C}_{\gamma c}
$$

provides $F=F_{D} \in \mathcal{K}_{\gamma c}$.
Example 6.7 Consider the interconnected system (6.143) - (6.145) composed by $N=2$ local subsystems with matrices $(i=1,2)$

$$
\begin{gathered}
A^{i i}=\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right], \quad B_{1}^{i}=\left[\begin{array}{l}
1 \\
1
\end{array}\right], \quad B_{2}^{i}=\left[\begin{array}{l}
0 \\
1
\end{array}\right] \\
C_{1}^{i}=\left[\begin{array}{ll}
0.5 & 0
\end{array}\right], \quad D_{12}^{i}=1
\end{gathered}
$$

coupled together by the interconnection matrices

$$
A^{12}=A^{21}=\left[\begin{array}{rr}
-1 & 0 \\
2 & 0
\end{array}\right]
$$



Figure 6.15: The static output feedback control system

Solving the optimal control problem in $\mathrm{RH}_{2}$, the optimal feedback gain and the associated minimum cost are

$$
F_{2}=\left[\begin{array}{llll}
-2.84 & -2.15 & -0.34 & -0.15 \\
-0.34 & -0.15 & -2.84 & -2.15
\end{array}\right], \quad J_{o p t}=\|T(z, w ; s)\|_{2}^{2}=21.79
$$

while the optimal solution of Problem (6.154) provides

$$
F_{D}=\left[\begin{array}{rrrr}
-3.62 & -2.50 & 0 & 0 \\
0 & 0 & -3.62 & -2.50
\end{array}\right], \quad J_{s u b}=30.27
$$

However, for $F=F_{D}$ the exact value of the objective function can be calculated as being $\|T(z, w ; s)\|_{2}^{2}=23.04$ which confirms numerically the fact that only an upper bound of the true cost has been minimized.

### 6.6.2 Static output control design

The system to be analyzed is given in fig. 6.15. The transfer function $P(s)$ has the state space realization

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{6.155}\\
z & =C_{1} x+D_{12} u  \tag{6.156}\\
y & =C_{2} x \tag{6.157}
\end{align*}
$$

where in opposition to the Output feedback control problems solved before, the measured output variable $y(\cdot) \in R^{r}$ is not corrupted by the external disturbance $w(t)$. Moreover it is assumed that $C_{2}$ is full row rank. The goal is to design a static output feedback law of the form

$$
\begin{equation*}
u=G y \tag{6.158}
\end{equation*}
$$

where the internal stability as well as some previously defined performance behavior are assured.

The control law (6.158) can be rewritten as $u=F_{O} x$, where the subscript " O " in matrix $F$ means that there exists $G$ such that $F_{O}=G C_{2}$. In this case the internal stability of the closed-loop system is preserved by means of the static output control law (6.158) provided $F=F_{O} \in \mathcal{K}_{c}$. Before proceed, let us made the assumption that matrix $C_{2}$ presents the following structure

$$
C_{2}=\left[\begin{array}{ll}
I & 0
\end{array}\right]
$$

If this is not the case, it is always possible to put the system in this form by a suitable choice of a similarity transformation. Under this assumption, matrix $F_{O}$ must present the particular structure

$$
F_{O}=\left[\begin{array}{ll}
F_{1} & 0
\end{array}\right]
$$

which implies that, $G=F_{1}$ as required before. For a symmetric matrix $P$ the subscript "O" has a slightly different meaning, that is

$$
P=P_{O}=\left[\begin{array}{cc}
P_{11} & 0 \\
0 & P_{22}
\end{array}\right]
$$

where $P_{11}$ and $P_{22}$ are $r \times r$ and $n-r \times n-r$ symmetric sub-matrices respectively. Unfortunately, the set of all stabilizing output feedback gains, being

$$
F=F_{O} \in \mathcal{K}_{c}
$$

is nonconvex and in many instances may be constituted of disjoint subsets. To circumvent this difficulty, we work here with a subset of it which as will be proven can be converted to a convex set. This subset is characterized by the following definition

Definition 6.4 (Structural O-Stabilizability) The triple $\left(A, B_{2}, C_{2}\right)$ is said to be structurally $O$ - stabilizable if there exist matrices $P_{O}$ symmetric and positive definite and $F_{O}$ such that

$$
\begin{equation*}
0 \geq\left(A+B_{2} F_{O}\right) P_{O}+P_{O}\left(A+B_{2} F_{O}\right)^{\prime}+Q \tag{6.159}
\end{equation*}
$$

for some matrix $Q=Q^{\prime}>0$. The set of all such matrices $F_{O}$ is denoted $\mathcal{K}_{O}$.
For any matrix such that $F=F_{O}$, it can be factorized as $F=G C_{2}$ for some $G$, consequently structural O - stability implies that the eigenvalues of the closed-loop matrix $A+B_{2} F=A+B_{2} G C_{2}$ are all in the open left hand side part of the complex plane. However, as in the case of decentralized control, Definition 6.4 requires the Lyapunov inequality solution, used to test stability presents a particular structure. The consequence is that $F_{O} \in \mathcal{K}_{O}$ implies $F_{O} \in \mathcal{K}_{c}$ but the inverse is not true in general. The price to be paid to handle convex sets only is to retain a subset of the entire set of static output feedback stabilizing gains. With no loss of generality, matrix $Q>0$ in definition 6.4 can be a priori fixed. This fact is proved with no major difficulty since for any $\bar{P}=\bar{P}_{O}>0$ and any scalar $\beta>0$ then $P=\beta \bar{P}_{O}$ satisfies $P=P_{O}$. Considering a matrix $W$ partitioned in four blocks, the subscript "O" stands for

$$
W_{O}:=\left[\begin{array}{ll}
W_{1 O} & W_{2 O} \\
W_{2 O}^{\prime} & W_{3}
\end{array}\right]
$$

which indicates that only the sub-blocks $W_{1}$ and $W_{2}$ have to present the output structure, namely

$$
W_{1 O}=\left[\begin{array}{cc}
W_{11} & 0 \\
0 & W_{22}
\end{array}\right], \quad W_{2 O}=\left[\begin{array}{c}
W_{21} \\
0
\end{array}\right]
$$

with $W_{11}, W_{22}$ and $W_{21}$ being $r \times r, n-r \times n-r$ and $r \times m$ matrices where $n$ is the state vector dimension, $m$ is the control vector dimension and $r$ is the output vector dimension respectively.

Theorem 6.13 Consider $B_{1} B_{1}^{\prime}>0$ and define the set

$$
\begin{equation*}
\mathcal{C}_{O}:=\left\{W: W=W_{O}\right\} \cap \mathcal{C}_{c} \tag{6.160}
\end{equation*}
$$

The following are true
a) $\mathcal{C}_{O}$ is a convex set.
b) Each $W \in \mathcal{C}_{O}$ is such that $W_{1}>0$.
c) $\mathcal{K}_{O}=\left\{W_{2}^{\prime} W_{1}^{-1}: W \in \mathcal{C}_{O}\right\}$.

Proof As in the proof of Theorem 6.11, the constraint $W=W_{O}$ is linear and hence convex. This fact together with $B_{1} B_{1}^{\prime}>0$ and $\mathcal{C}_{O} \subset \mathcal{C}_{c}$ prove both points a) and b).

Point c) For the necessity, assume $F_{O} \in \mathcal{K}_{O} \neq \emptyset$. Setting $Q=B_{1} B_{1}^{\prime}$ in (6.159), it is simple to verify that

$$
W=W_{O}=\left[\begin{array}{cc}
P_{O} & P_{O} F_{O}^{\prime} \\
F_{O} P_{O} & F_{O} P_{O} F_{O}^{\prime}
\end{array}\right] \in \mathcal{C}_{O}
$$

and $W_{2}^{\prime} W_{1}^{-1}=F_{O} P_{O} P_{O}^{-1}=F_{O}$. Conversely, for any $W \in \mathcal{C}_{O} \neq \emptyset$ and any $v \in \mathcal{N}_{c}$ we get

$$
\begin{aligned}
0 & \geq v^{\prime} \Theta_{c}(W) v \\
& \geq x^{\prime}\left[\left(A+B_{2} W_{2}^{\prime} W_{1}^{-1}\right) W_{1}+W_{1}\left(A+B_{2} W_{2}^{\prime} W_{1}^{-1}\right)^{\prime}+B_{1} B_{1}^{\prime}\right] x
\end{aligned}
$$

which shows that inequality (6.159) holds for $P_{O}=W_{1}=W_{1 O}$ and $F_{O}=W_{2}^{\prime} W_{1}^{-1}=$ $W_{2 O}^{\prime} W_{1 O}^{-1}$, that is both matrices exhibit the desired output feedback structure. In case one of the sets $\mathcal{K}_{O}$ or $\mathcal{C}_{O}$ is empty then the equality stated in point c) follows trivially.

From this theorem, we can see that any feasible matrix $W \in \mathcal{C}_{O}$ generates an output feedback stabilizing gain which is very simple to be determined. Actually, the imposed structure constraint

$$
W_{1 O}=\left[\begin{array}{cc}
W_{11} & 0 \\
0 & W_{22}
\end{array}\right], \quad W_{2 O}=\left[\begin{array}{c}
W_{21} \\
0
\end{array}\right]
$$

provides $F_{O} \in \mathcal{K}_{O}$ which can be factorized as $F_{O}=G C_{2}$ where

$$
\begin{equation*}
G=F_{1}=W_{21}^{\prime} W_{11}^{-1} \tag{6.161}
\end{equation*}
$$

in addition, the complete parametrization of the set $\mathcal{K}_{O}$ by means of point $c$ ) opens the possibility, with no major difficulty, to involve it in an optimization procedure. For instance consider the Static output feedback control design problem

$$
\begin{equation*}
\min \left\{\|T(z, w ; s)\|_{2}^{2}: F \in \mathcal{K}_{O}\right\} \tag{6.162}
\end{equation*}
$$

which makes once again explicit that the feasible set is restrict to those gains satisfying Definition 6.4.

Theorem 6.14 Assume $B_{1} B_{1}^{\prime}>0$ and let $\bar{W}$ be the optimal solution of the convex programming problem

$$
\begin{equation*}
J_{s u b}:=\min \left\{\operatorname{trace}\left[R_{c} W\right]: W \in \mathcal{C}_{O}\right\} \tag{6.163}
\end{equation*}
$$

Then, $F=F_{O}=\bar{W}_{2}^{\prime} \bar{W}_{1}^{-1} \in \mathcal{K}_{O}$ minimizes an upper bound of the objective function of Problem (6.162) in the sense that $\|T(z, w ; s)\|_{2}^{2} \leq \operatorname{trace}\left[R_{c} W\right]$, for all $W \in \mathcal{C}_{O}$.
Proof From Theorem 6.13, it suffices to consider that both problems are feasible. In this case, the optimal solution of problem (6.163) provides $\bar{W}_{2}^{\prime} \bar{W}_{1}^{-1} \in \mathcal{K}_{O}$. Additionally, for any $W \in \mathcal{C}_{O}$ we have $W_{1 O} \geq P_{r}$ where

$$
0=\left(A+B_{2} W_{2 O}^{\prime} W_{1 O}^{-1}\right) P_{r}+P_{r}\left(A+B_{2} W_{2 O}^{\prime} W_{1 O}^{-1}\right)^{\prime}+B_{1} B_{1}^{\prime}
$$

yielding immediately trace $\left[R_{c} W\right] \geq\|T(z, w ; s)\|_{2}^{2}$. The proof is complete.
In general, the optimal solution of problem (6.162) is not generated by means of the global optimal solution of the convex problem (6.163). Even though all gains in $\mathcal{K}_{O}$ are generated by the proposed convex parametrization, matrix $P_{r}$ needed to calculate $\|T(z, w ; s)\|_{2}^{2}$ may not satisfy the constraint $P_{r}=\left(P_{r}\right)_{O}$, in this case there is no matrix $W \in \mathcal{C}_{O}$ for which the equality $W_{1 O}=P_{r}$ holds.

Remark 6.39 Following the lines of Theorem 6.13 we notice that the set of all $F=F_{O} \in$ $\mathcal{K}_{c}$ can be generated by $F=W_{2}^{\prime} W_{1}^{-1}$ where

$$
W \in \mathcal{C}_{c O}:=\left\{W: W_{2}^{\prime} W_{1}^{-1}=\left(W_{2}^{\prime} W_{1}^{-1}\right) o\right\} \cap \mathcal{C}_{c}
$$

Unfortunately this set is not convex. The same theorem provides a subset of $\mathcal{C}_{c O}$, namely $\mathcal{C}_{O}$ which has the important property to be convex.

Remark 6.40 Consider the system (6.155) - (6.157) and the matrix

$$
S:=\left[\begin{array}{l}
C_{2} \\
E_{2}
\end{array}\right]
$$

where $E_{2}$ is any matrix of appropriate dimension such that $S^{-1}$ exists. Defining the new state variable $\tilde{x}:=S x$, the new state space realization turns out to be such that $y=\tilde{C}_{2} \tilde{x}$ where

$$
\tilde{C}_{2}=C_{2} S^{-1}=\left[\begin{array}{ll}
I & 0
\end{array}\right]
$$

The concept of Structural $O$ - stabilizability depends on the similarity matrix $S$ by means of the arbitrary submatrix $E_{2}$. The next lemma puts in evidence this important point.

Lemma 6.7 The following statements are equivalent:
a) The triple $\left(A, B_{2}, C_{2}\right)$ is stabilizable by output feedback, that is there exists $G$ such that the matrix $A+B_{2} G C_{2}$ is stable.
b) There exists a matrix $E_{2}$ such that the triple $\left(\tilde{A}, \tilde{B}_{2}, \tilde{C}_{2}\right)$ in the new state space representation is structurally $O$ - stabilizable.

Proof To prove that a) implies b), let us suppose that there exists an output feedback gain $G$ such that $A+B_{2} G C_{2}$ is stable. From the Extended Lyapunov lemma, there exists $P=P^{\prime}>0$ solution to the linear equation

$$
0=\left(A+B_{2} G C_{2}\right) P+P\left(A+B_{2} G C_{2}\right)^{\prime}+Q
$$

for $Q=Q^{\prime}>0$ given. Chosen $E_{2}=U_{2}^{\prime} P^{-1}$ where $U_{2}$ defines an orthogonal basis to the null space of $C_{2}$, that is $C_{2} U_{2}=0$ and $U_{2}^{\prime} U_{2}=I$ and multiplying the above equation to the right by $S^{\prime}$ and to the left by $S$ we get

$$
0=\left(\tilde{A}+\tilde{B}_{2} G \tilde{C}_{2}\right) \tilde{P}+\tilde{P}\left(\tilde{A}+\tilde{B}_{2} G \tilde{C}_{2}\right)^{\prime}+\tilde{Q}
$$

where $\tilde{P}:=S P S^{\prime}$ and $\tilde{Q}:=S Q S^{\prime}$. Simple calculations show that

$$
\tilde{F}:=G \tilde{C}_{2}=\left[\begin{array}{ll}
G & 0
\end{array}\right]=\tilde{F}_{O}
$$

and

$$
\tilde{P}=S P S^{\prime}=\left[\begin{array}{cc}
C_{2} P C_{2}^{\prime} & 0 \\
0 & U_{2}^{\prime} P^{-1} U_{2}
\end{array}\right]=\tilde{P}_{O}
$$

meaning that point $b$ ) holds. The converse is immediate.

Remark 6.41 The numerical solution of problem (6.162) by means of convex programming methods is very efficient. However the main drawback of this approach is that only a subset of the true feasible set of stabilizing output feedback gains is considered. Let us verify the difficulties we have to face for the numerical solution of the true Output feedback design problem in $\mathrm{RH}_{2}$

$$
\min \left\{\|T(z, w ; s)\|_{2}^{2}: G C_{2} \in \mathcal{K}_{c}\right\}
$$

which reduces to the determination of matrices $G$ and $P=P^{\prime} \geq 0$ such that

$$
\min \left\{\operatorname{trace}\left[B_{1}^{\prime} P B_{1}\right]: 0=A_{c g}^{\prime} P+P A_{c g}+C_{c g}^{\prime} C_{c g}\right\}
$$

where $A_{c g}:=A+B_{2} G C_{2}$ and $C_{c g}:=C_{1}+D_{12} G C_{2}$. Defining $\Lambda=\Lambda^{\prime}$ the matrix of Lagrange multiplier associated to the equality constraint, the necessary optimality conditions are (recall Remark 6.7)

$$
\begin{aligned}
& 0=A_{c g}^{\prime} P+P A_{c g}+C_{c g}^{\prime} C_{c g} \\
& 0=A_{c g} \Lambda+\Lambda A_{c g}^{\prime}+B_{1} B_{1}^{\prime} \\
& 0=\left(G C_{2}+D_{12}^{\prime} C_{1}+B_{2}^{\prime} P\right) \Lambda C_{2}^{\prime}
\end{aligned}
$$

Assuming as before that $B_{1} B_{1}^{\prime}>0$ then for any $G C_{2} \in \mathcal{K}_{c}$ matrix $\Lambda$ is positive definite, in which case the last equation yields

$$
G=-\left(B_{2}^{\prime} P+D_{12}^{\prime} C_{1}\right) \Lambda C_{2}^{\prime}\left(C_{2} \Lambda C_{2}^{\prime}\right)^{-1}
$$

This formula for the optimal gain couples the first two equations in a very nonlinear manner. The optimal gain can not be expressed in terms of a Riccati equation unless matrix $C_{2}$ is square and nonsingular. The following algorithm is useful for numerical purposes

1) Choose $G_{0}$ such that $F_{0}:=G_{0} C_{2} \in \mathcal{K}_{c}$ and iterate until convergence with $k=0,1, \cdots$.
2) Set $F_{k}:=G_{k} C_{2}$ and let $P_{k} \geq 0$ be the solution of the Lyapunov equation (in the unknown $P$ )

$$
0=\left(A+B_{2} F_{k}\right)^{\prime} P+P\left(A+B_{2} F_{k}\right)+\left(C_{1}+D_{12} F_{k}\right)^{\prime}\left(C_{1}+D_{12} F_{k}\right)^{\prime}
$$

Determine $\tilde{F}_{k}:=-B_{2}^{\prime} P_{k}-D_{12}^{\prime} C_{1}$.
3) Let $\Lambda_{k+1}>0$ and $G_{k+1}$ be the solution of the nonlinear equations (in the unknown $\Lambda$ and $G$ )

$$
\begin{gathered}
0=\left(A+B_{2} G C_{2}\right) \Lambda+\Lambda\left(A+B_{2} G C_{2}\right)^{\prime}+B_{1} B_{1}^{\prime} \\
G=\tilde{F}_{k} \Lambda C_{2}^{\prime}\left(C_{2} \Lambda C_{2}^{\prime}\right)^{-1}
\end{gathered}
$$

Lemma 6.8 With $B_{1} B_{1}^{\prime}>0$, the above algorithm has the following properties
a) For all $k=0,1, \cdots$ the gain $F_{k+1}=G_{k+1} C_{2} \in \mathcal{K}_{c}$.
b) For all $k=0,1, \cdots$ the criterion $J_{k+1}:=\operatorname{trace}\left[B_{1}^{\prime} P_{k+1} B_{1}\right] \leq \operatorname{trace}\left[B_{1}^{\prime} P_{k} B_{1}\right]:=J_{k}$

Proof Point a) is a consequence of step 3). Actually, it implies that $\Lambda_{k+1}>0$ satisfies

$$
0=\left(A+B_{2} F_{k+1}\right) \Lambda_{k+1}+\Lambda_{k+1}\left(A+B_{2} F_{k+1}\right)^{\prime}+B_{1} B_{1}^{\prime}
$$

thus, $F_{k+1}=G_{k+1} C_{2}$ is a stabilizing gain.
Point b) Simple but tedious algebraic manipulations show that for two subsequent iterations

$$
\begin{aligned}
0= & \left(A+B_{2} F_{k+1}\right)^{\prime}\left(P_{k}-P_{k+1}\right)+\left(P_{k}-P_{k+1}\right)\left(A+B_{2} F_{k+1}\right)+ \\
& +\left(F_{k}-\tilde{F}_{k}\right)^{\prime}\left(F_{k}-\tilde{F}_{k}\right)-\left(F_{k+1}-\tilde{F}_{k}\right)^{\prime}\left(F_{k+1}-\tilde{F}_{k}\right)
\end{aligned}
$$

which allows us to determine the current value of the criterion as being

$$
\begin{aligned}
J_{k}-J_{k+1}= & \operatorname{trace}\left[\int_{0}^{\infty} B_{1} e^{\left(A+B_{2} F_{k+1}\right) t}\left(F_{k}-\tilde{F}_{k}\right)^{\prime}\left(F_{k}-\tilde{F}_{k}\right) e^{\left(A+B_{2} F_{k+1}\right)^{\prime} t} B_{1}^{\prime} d t\right]- \\
& -\operatorname{trace}\left[\int_{0}^{\infty} B_{1} e^{\left(A+B_{2} F_{k+1}\right) t}\left(F_{k+1}-\tilde{F}_{k}\right)^{\prime}\left(F_{k+1}-\tilde{F}_{k}\right) e^{\left(A+B_{2} F_{k+1}\right)^{\prime} t} B_{1}^{\prime} d t\right] \\
= & \operatorname{trace}\left[\left(F_{k}-\tilde{F}_{k}\right) \Lambda_{k+1}\left(F_{k}-\tilde{F}_{k}\right)^{\prime}-\left(F_{k+1}-\tilde{F}_{k}\right) \Lambda_{k+1}\left(F_{k+1}-\tilde{F}_{k}\right)^{\prime}\right] \\
= & \operatorname{trace}\left[\left(G_{k} C_{2}-\tilde{F}_{k}\right) \Lambda_{k+1}\left(G_{k} C_{2}-\tilde{F}_{k}\right)^{\prime}\right]- \\
& -\operatorname{trace}\left[\left(G_{k+1} C_{2}-\tilde{F}_{k}\right) \Lambda_{k+1}\left(G_{k+1} C_{2}-\tilde{F}_{k}\right)^{\prime}\right]
\end{aligned}
$$

using now the value of $G_{k+1}$ provided in step 3)

$$
G_{k+1}=\tilde{F}_{k} \Lambda_{k+1} C_{2}^{\prime}\left(C_{2} \Lambda_{k+1} C_{2}^{\prime}\right)^{-1}
$$

we get

$$
\begin{aligned}
J_{k}-J_{k+1}= & \operatorname{trace}\left[\left(G_{k} C_{2}-\tilde{F}_{k}\right) \Lambda_{k+1}\left(G_{k} C_{2}-\tilde{F}_{k}\right)^{\prime}\right]- \\
& -\operatorname{trace}\left[\tilde{F}_{k}\left(\Lambda_{k+1}-\Lambda_{k+1} C_{2}^{\prime}\left(C_{2} \Lambda_{k+1} C_{2}^{\prime}\right)^{-1} C_{2} \Lambda_{k+1}\right) \tilde{F}_{k}^{\prime}\right] \\
= & \operatorname{trace}\left[G_{k} C_{2} \Lambda_{k+1} C_{2}^{\prime} G_{k}-G_{k+1} C_{2} \Lambda_{k+1} C_{2}^{\prime} G_{k}^{\prime}-\right. \\
& \left.-G_{k} C_{2} \Lambda_{k+1} C_{2}^{\prime} G_{k+1}^{\prime}+G_{k+1} C_{2} \Lambda_{k+1} C_{2}^{\prime} G_{k+1}^{\prime}\right] \\
= & \operatorname{trace}\left[\left(G_{k}-G_{k+1}\right) C_{2} \Lambda_{k+1} C_{2}^{\prime}\left(G_{k}-G_{k+1}\right)^{\prime}\right]
\end{aligned}
$$

This equality finally enables us to conclude that

$$
\begin{aligned}
J_{k+1} & =J_{k}-\operatorname{trace}\left[\left(G_{k}-G_{k+1}\right) C_{2} \Lambda_{k+1} C_{2}^{\prime}\left(G_{k}-G_{k+1}\right)^{\prime}\right] \\
& \leq J_{k}
\end{aligned}
$$

which proves point b) of the lemma proposed.
It is important to recognize that the second part of Lemma 6.8 is also a consequence of the joint determination of matrices $\Lambda_{k+1}$ and $G_{k+1}$ in step 3) from the solution of two simultaneous equations, one of then nonlinear. It would be very attractive to change the previous algorithm in order to solve only linear equations. However, doing this it is no more possible to be sure that

$$
F_{k}=G_{k} C_{2} \in \mathcal{K}_{c} \Longrightarrow F_{k+1}=G_{k+1} C_{2} \in \mathcal{K}_{c}
$$

in which occurrence the properties introduced in Lemma 6.8 are both lost. Hence, we can say that the price to be paid to keep these convergence features is to solve in step 3) a nonlinear matrix equation.

Example 6.8 Consider the system (6.155) - (6.157) with the following data

$$
\begin{gathered}
A=\left[\begin{array}{rr}
0 & 1 \\
-1 & 0
\end{array}\right], \quad B_{1}=\left[\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right], \quad B_{2}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \\
C_{1}=\left[\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right], \quad D_{12}=\left[\begin{array}{l}
0 \\
1
\end{array}\right], \quad C_{2}=\left[\begin{array}{ll}
0 & 1
\end{array}\right]
\end{gathered}
$$

Solving the State feedback control problem in $\mathrm{RH}_{2}$ we get

$$
F_{2}=\left[\begin{array}{ll}
-0.41 & -0.91
\end{array}\right], \quad J_{o p t}=\operatorname{trace}\left[P_{2}\right]=2.19
$$

Since matrix $C_{2}$ is not in the standard form, from the discussion in Remark 6.39 we have first defined a similarity transformation with

$$
S=\left[\begin{array}{ll}
0 & 1 \\
1 & 1
\end{array}\right]
$$

together with the optimal solution of problem (6.163) provides

$$
G_{s u b}=-2.0, \quad J_{s u b}=3.50
$$

In this case, the optimal solution of the exact design problem (recall Remark 6.41) is known to be $G_{o p t}=-0.81$ with the associated cost $J_{o p t}=2.46$. This example shows that the quality of the solution provided by the solution of the convex problem (6.163) may be poor. Of course the quality, measured in terms of the deviation from the optimal solution depends on the similarity transformation matrix $S$ which unfortunately is not easy to choose. Again, the main attractive feature of problem (6.163) is its convexity.

### 6.7 Notes and references

The various aspects of Convex Analysis, from the basic facts to most important and deep results are included in the seminal book by Rockafellar [53]. The book [10] exhaustively treats the most important topics related to systems and control theory in the framework of Linear Matrix Inequalities so that the problems to be handled are convex. As well as, this book is also an important source of references to those interested in going deeper into optimal control design using convex programming techniques. Section 6.2 is mainly based on papers [6], [22], [47], [48] and on the results introduced in the previous chapters of this book. Section 6.3 follows the same lines of [32] and connections to the results of [5] are put in evidence. The approach proposed in [40] is summarized in Remark 6.28 where also a numeric example is included for comparison purpose. In Section 6.4 the pole placement problem in a circular region is analyzed in a convex programming view point. Other regions for pole placement as well as a different approach to solve the same problem are considered in [28]. The time-domain specifications treated in Section 6.5 is based on the result of Lemma 6.5 due to Wilson [62]. The proof of this result is different and new. The related optimal control design problems appeared first in [54]. Finally, the Decentralized control design in Section 6.6 parallels the results of [23]. In the same section the Static output control design problem considered, first appeared in [39]. From this reference comes both the Example 6.8 solved and the algorithm discussed in Remark 6.41. The proof of convergence is new.

## Chapter 7

## Uncertain Systems Control Design

### 7.1 Introduction

This chapter is devoted to uncertain systems control design. The main goal is to provide a simple and easy to follow introduction to this important subject of control theory. The necessity to analyze dynamic systems subject to uncertainties stems from the fact that the model to be used for design purpose is generally only an approximation of reality. In other words, for control design purpose we need to handle simple models. However, the controller obtained must work when connected to the real system. The way to take this feature into account is to consider a simplified nominal model which is corrupted by uncertainties belonging to a prespecified domain. To make clear this point it is interesting to remember the mixed $R H_{2} / R H_{\infty}$ control design problem. The controller calculated from its solution imposes optimality to the nominal closed loop system while stability is preserved against the considered uncertainties. In this chapter, we go beyond this point. Indeed, not only robust stability is considered but also robust performance is taken into account. That is, the main purpose is to design controllers which preserve stability and minimize the performance loss due to existence of uncertainties. In this way, two important concepts are introduced, namely quadratic stabilizability and guaranteed cost which are both on the basis of the results that follow. Once again, the optimal control problems to be solved are all convex and so the same machinery provided in Chapter 6 is intensively used. It is important to make explicit that only linear control design is considered.

### 7.2 Robust stability and performance

In this section we generalize the concepts of stability and optimality to deal with uncertainties belonging to a precise although general domain. Control problems design involving norms of transfer functions in $R H_{2}$ and $R H_{\infty}$ spaces are considered in the special case that the whole state is available for feedback.

The basic system structure is depicted in fig. 7.1 making explicit the linear state feedback control law to be used. The transfer function $P_{\mathcal{D}}(s)$ depends on the subscript $\mathcal{D}$ which will be made clear in the sequel. The linear system under consideration thus


Figure 7.1: The control system structure
have the following state space representation

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{7.1}\\
z & =C_{1} x+D_{12} u  \tag{7.2}\\
y & =x \tag{7.3}
\end{align*}
$$

and the controller is completely defined by means of the matrix $F$ which has to be determined (if any) such that with

$$
\begin{equation*}
u=F x \tag{7.4}
\end{equation*}
$$

some desired performance is assured. Equations (7.1) - (7.3) represent an uncertain linear system. That is, matrices $A$ and $B_{2}$ are not exactly known. On the contrary we only know that they belong to some prespecified matrix set $\mathcal{D}$. In this sense, equations (7.1) - (7.3) represent in fact a family of linear systems any member of which has to be controlled by means of the same control law. Of course it is important, if possible, to determine $F$ such that all members of the family are internally stable and, at the same time, some desired performance level is guaranteed.

To put this discussion in more precise terms, let us recall that for any pair $\left(A, B_{2}\right) \in$ $\mathcal{D}$ the internal stability of the corresponding system is assured whenever $F \in \mathcal{K}_{c}$. Consequently, the matrix gain $F$ stabilizes all members of the family if and only if

$$
\begin{equation*}
F \in \mathcal{K}_{\mathcal{D}}:=\bigcap_{\left(A, B_{2}\right) \in \mathcal{D}} \mathcal{K}_{c} \tag{7.5}
\end{equation*}
$$

This fact has a very clear interpretation. The set $\mathcal{K}_{\mathcal{D}}$ is composed by all matrices $F$ which stabilize each feasible pair $\left(A, B_{2}\right) \in \mathcal{D}$. The set $\mathcal{K}_{\mathcal{D}}$ does not present any important property as for instance convexity. A crucial point is that although each element of the intersection above can be converted to a convex set (as has successfully been done in Chapter 6), the same is not true for the intersection itself, that is for the set $\mathcal{K}_{\mathcal{D}}$. A way to circumvent this difficulty is to work with a subset of it, characterized by means of the following definition.

Definition 7.1 (Quadratic stabilizability) The pair $\left(A, B_{2}\right)$ is said to be quadratic stabilizable if there exist matrices $P$, symmetric and positive definite, and $F$ such that

$$
\begin{equation*}
0 \geq\left(A+B_{2} F\right) P+P\left(A+B_{2} F\right)^{\prime}+Q, \quad \forall\left(A, B_{2}\right) \in \mathcal{D} \tag{7.6}
\end{equation*}
$$

for some matrix $Q=Q^{\prime}>0$. The set of all such matrices $F$ is denoted $\mathcal{K}_{Q}$.

Once again it is important to observe that in this definition, the choice of matrix $Q>0$ is immaterial as far as the existence of a matrix $P>0$ satisfying (7.6) is concerned. If this inequality is satisfied for $P=\bar{P}$ when $Q=\bar{Q}$ then for any other $Q>0$ the same inequality also holds true for $P=\beta \bar{P}$ provided $\beta>0$ is such that $\beta \bar{Q} \geq Q$. Definition 7.1 has a very interesting interpretation. Pick from $\mathcal{D}$ a pair ( $A, B_{2}$ ), consider the system (7.1) with $w(t)=0$ and define the Lyapunov function

$$
v(x)=x^{\prime} P^{-1} x
$$

If inequality (7.6) holds then the time derivative of $v(\cdot)$ along any trajectory of the closed loop system satisfies

$$
\begin{aligned}
\dot{v}(x) & =x^{\prime}\left[\left(A+B_{2} F\right)^{\prime} P^{-1}+P^{-1}\left(A+B_{2} F\right)\right] x \\
& =x^{\prime} P^{-1}\left[\left(A+B_{2} F\right) P+P\left(A+B_{2} F\right)^{\prime}\right] P^{-1} x \\
& \leq-x^{\prime} P^{-1} Q P^{-1} x \\
& <0, \forall x \neq 0
\end{aligned}
$$

showing that the closed-loop system is internally stable indeed. The point is that the above calculation holds for any pair $\left(A, B_{2}\right) \in \mathcal{D}$ and so quadratic stability means that only one Lyapunov function can be used to test stability of all systems generated by all pairs $\left(A, B_{2}\right)$ in $\mathcal{D}$. Thus, it is clear that $\mathcal{K}_{Q}$ is only a subset of $\mathcal{K}_{\mathcal{D}}$ defined in (7.5) since for the latter many different Lyapunov functions can be used to test stability.

Remark 7.1 An uncertain linear system is said to be robustly stabilizable if $\mathcal{K}_{\mathcal{D}} \neq \emptyset$. This means that there exists a feedback gain matrix $F$ such that for all $\left(A, B_{2}\right) \in \mathcal{D}$, the eigenvalues of $A+B_{2} F$ lie in the open left complex plane.

At this point we have to move our attention to the uncertainty domain $\mathcal{D}$ since Definition 7.1 depends essentially on its mathematical description. There are several possibilities for that. Let us consider one of such description whose generality will be discussed in the sequel. Notice first that we are assuming only matrices $A$ and $B_{2}$ of tioo open-loop system (7.1) - (7.3) to be uncertain. This is the same to say that matıix (recall Chapter 6)

$$
M_{c}:=\left[\begin{array}{cc}
A & B_{2}  \tag{7.7}\\
0 & 0
\end{array}\right]
$$

is not exactly known but belongs to an uncertainty domain $\mathcal{D}_{c}$ such that

$$
\left(A, B_{2}\right) \in \mathcal{D} \Longleftrightarrow M_{c} \in \mathcal{D}_{c}
$$

Then, we proceed by defining the uncertainty domain $\mathcal{D}_{c}$. It is characterized as a polyhedral convex bounded domain given by

$$
\begin{equation*}
\mathcal{D}_{c}:=\operatorname{co}\left\{M_{c i}, i=1,2, \cdots, N\right\} \tag{7.8}
\end{equation*}
$$

where $\operatorname{co}\{\cdot\}$ denotes convex hull (recall Appendix H) generated by matrices $M_{c i}, i=$ $1,2, \cdots, N$. From this, any matrix $M_{c} \in \mathcal{D}_{c}$ can be determined by a convex combination of the extreme matrices $M_{c i}, i=1,2, \cdots, N$. More precisely,

$$
\begin{equation*}
M_{c}=\sum_{i=1}^{N} \xi_{i} M_{c i}, \quad \xi_{i} \geq 0, \quad \sum_{i=1}^{N} \xi_{i}=1 \tag{7.9}
\end{equation*}
$$



Figure 7.2: The uncertainty domain

This uncertainty description is very general. Any polyhedral and convex bounded domain can be exactly represented by a proper choice of the extreme matrices $M_{c i}$. On the other hand, a convex domain (not necessarily polyhedral) can be represented only approximately but the degree of approximation involved can be made arbitrarily small if the number of extreme matrices becomes sufficiently large. This fact is illustrated in fig. 7.2.

Remark 7.2 Another important uncertainty description is the so called norm bounded domain defined as

$$
\mathcal{D}_{n}:=\left\{M_{c}=M_{c n}+B_{n} \Omega C_{n}:\|\Omega\| \leq 1\right\}
$$

where matrix $M_{c n}$ defines the open-loop nominal system and matrices $B_{n}$ and $C_{n}$ of appropriate dimension define the uncertainty structure. First of all, notice that this set is convex. Indeed, the generic element $M_{c}$ is affine with respect to $\Omega$ which varies in a convex set defined by the norm constraint $\|\Omega\| \leq 1$. Furthermore, this set is not polyhedral in general and so, as it has been commented before, only an approximation of it can be generated by means of the convex hull (7.8). However, in some particular although important cases, it degenerates to a polyhedral convex set. For instance, this occurs when the uncertain matrix $\Omega$ is constrained to be diagonal.

Example 7.1 Consider the following linear uncertain system

$$
\dot{x}=A x+B_{2} u
$$

where matrices $A$ and $B_{2}$ are given by

$$
A=\left[\begin{array}{cc}
0 & \alpha-1 \\
\beta & 0
\end{array}\right], \quad B_{2}=\left[\begin{array}{c}
\alpha \\
1-\beta
\end{array}\right]
$$

and the parameters $\alpha$ and $\beta$, representing the uncertainties are such that

$$
|\alpha-0.5| \leq 0.3, \quad|\beta-0.5| \leq 0.3
$$

Letting the nominal values for these parameters be $\alpha_{0}=\beta_{0}=0.5$, the norm bounded domain $\mathcal{D}_{n}$ is completely defined by matrices

$$
M_{c n}=\left[\begin{array}{rrr}
0 & -0.5 & 0.5 \\
0.5 & 0 & 0.5 \\
0 & 0 & 0
\end{array}\right], \quad B_{n}=\left[\begin{array}{rr}
0.3 & 0 \\
0 & 0.3 \\
0 & 0
\end{array}\right], \quad C_{n}^{\prime}=\left[\begin{array}{rr}
0 & 1 \\
1 & 0 \\
1 & -1
\end{array}\right]
$$

since the uncertain parameters may be retained in matrix $\Omega$, which for any $\|\Omega\| \leq 1$ exhibits the diagonal structure

$$
\Omega=\operatorname{diag}\left\{\frac{\alpha-0.5}{0.3}, \frac{\beta-0.5}{0.3}\right\}
$$

The same system can be modeled by means of a convex bounded domain $\mathcal{D}_{c}$. This is done from the observation that $N=4$ and the corresponding extreme matrices are generated from the four vertex of the rectangle $[0.2,0.8] \times[0.2,0.8]$, giving thus

$$
\begin{array}{ll}
M_{c 1}=\left[\begin{array}{rrr}
0 & -0.8 & 0.2 \\
0.2 & 0 & 0.8 \\
0 & 0 & 0
\end{array}\right], & M_{c 2}=\left[\begin{array}{rrr}
0 & -0.8 & 0.2 \\
0.8 & 0 & 0.2 \\
0 & 0 & 0
\end{array}\right] \\
M_{c 3}=\left[\begin{array}{rrr}
0 & -0.2 & 0.8 \\
0.2 & 0 & 0.8 \\
0 & 0 & 0
\end{array}\right], & M_{c 4}=\left[\begin{array}{rrr}
0 & -0.2 & 0.8 \\
0.8 & 0 & 0.2 \\
0 & 0 & 0
\end{array}\right]
\end{array}
$$

From this example we conclude that any uncertain system represented by matrix $M_{c}$ belongs to a polyhedral uncertain domain whenever each entry of $M_{c}$ is an affine function of the uncertainty. The number of independent uncertain parameters is immaterial since the domain $\mathcal{D}_{c}$ depends only on the extreme matrices $M_{c i}, i=1,2, \cdots, N$ which are a priori determined from the data.

We move now our attention to the characterization of uncertain systems optimality. Consider the closed-loop system of fig. 7.1 with the state space representation (7.1) - (7.3) and assume that the performance index can be expressed in terms of the transfer function from the input $w$ to the output $z$, that is

$$
\mathcal{J}\left(M_{c}, F\right):=\mathcal{G}(T(z, w ; s))
$$

where it is made explicit the dependence of the performance index with respect to both the control gain $F$ and the open-loop uncertain matrix $M_{c}$. The former has to be determined such that $\mathcal{J}(\cdot)$ is optimized in a sense to be precisely defined. Thus, restricting our attention to the quadratic stabilizing gains $F \in \mathcal{K}_{Q}$, the following definition is of particular importance.

Definition 7.2 (Guaranteed cost) The scalar $\rho$ is said to be a guaranteed cost associated with the feedback gain matrix $F \in \mathcal{K}_{Q}$ if

$$
\begin{equation*}
\mathcal{J}\left(M_{c}, F\right) \leq \rho, \quad \forall M_{c} \in \mathcal{D}_{c} \tag{7.10}
\end{equation*}
$$

The minimum value of $\rho$ satisfying (7.10), denoted $\rho_{Q}$, is called the minimum guaranteed cost.

From this definition, two points have to be kept in mind. The first one is that the guaranteed cost is a function of $F \in \mathcal{K}_{Q}$. Indeed, any $\rho \geq \rho(F)$ is also a guaranteed cost provided

$$
\rho(F):=\max _{M_{c} \in \mathcal{D}_{c}} \mathcal{J}\left(M_{c}, F\right)
$$

and second, the best choice for the control gain matrix $F \in \mathcal{K}_{Q}$ is the one such that the minimum of $\rho(F)$ is achieved, that is

$$
\rho_{Q}:=\min _{F \in \mathcal{K}_{Q}} \rho(F)
$$

Clearly, the scalar $\rho_{Q}$ is the minimum guaranteed cost associated to the uncertain system under consideration, it does not depend on a particular value of $F \in \mathcal{K}_{Q}$ and
it is the optimal solution of a min/max optimization problem. Actually, to show this, notice that

$$
\begin{align*}
\rho_{Q} & =\min _{F \in \mathcal{K}_{Q}} \rho(F) \\
& =\min _{F \in \mathcal{K}_{Q}} \max _{M_{c} \in \mathcal{D}_{c}} \mathcal{J}\left(M_{c}, F\right) \\
& =\min _{\rho, F \in \mathcal{K}_{Q}}\left\{\rho: \mathcal{J}\left(M_{c}, F\right) \leq \rho, \quad \forall M_{c} \in \mathcal{D}_{c}\right\} \tag{7.11}
\end{align*}
$$

which allows a very interesting interpretation of the minimum guaranteed cost, in terms of a game between $M a n$ and Nature. The Nature plays by fixing an openloop model $M_{c} \in \mathcal{D}_{c}$ and the Man defines the best feedback gain matrix by choosing $F \in \mathcal{K}_{Q}$ such that $\mathcal{J}\left(M_{c}, F\right)$ is minimized. Once this information is available, the Nature plays again but trying to destroy the Man's action by choosing a new openloop model $M_{c} \in \mathcal{D}_{c}$ such that $\mathcal{J}\left(M_{c}, F\right)$ is maximized. The guaranteed cost $\rho_{Q}$ is the equilibrium of this game. The equilibrium value $\rho_{Q}$ is generally very difficult to be calculated and so it must be thought as a paradigm for uncertain systems control design. A much more simple strategy, which is used in the sequel stems from the determination of a function $\bar{\rho}(F)$ such that

$$
\mathcal{J}\left(M_{c}, F\right) \leq \bar{\rho}(F), \quad \forall M_{c} \in \mathcal{D}_{c}
$$

which allows the determination of a guaranteed cost $\bar{\rho}_{Q}$ and the associated feedback gain by solving

$$
\bar{\rho}_{Q}:=\min _{F \in \mathcal{K}_{Q}} \bar{\rho}(F)
$$

clearly implying that

$$
\begin{aligned}
\bar{\rho}_{Q} & =\min _{F \in \mathcal{K}_{Q}} \bar{\rho}(F) \\
& \geq \min _{F \in \mathcal{K}_{Q}} \max _{M_{c} \in \mathcal{D}_{c}} \mathcal{J}\left(M_{c}, F\right) \\
& \geq \rho_{Q}
\end{aligned}
$$

In the next sections, our goal will be to solve this problem for the class of linear uncertain systems with polyhedral convex bounded domains. Once again, it will be possible to convert the optimization problem to be dealt with into a convex programming problem.

### 7.2.1 Quadratic stabilizing controllers

The main purpose here is to analyze the geometry of the set of all quadratic stabilizing feedback gains $\mathcal{K}_{Q}$ defined before (recall Definition 7.1). Consider the uncertain linear system

$$
\begin{equation*}
\dot{x}=A x+B_{2} u \tag{7.12}
\end{equation*}
$$

with $n$ states, $m$ inputs and where matrices $A$ and $B_{2}$ of appropriate and known dimension are such that $\left(A, B_{2}\right) \in \mathcal{D}$ or equivalently $M_{c} \in \mathcal{D}_{c}$. The set $\mathcal{D}_{c}$ is a polyhedral convex bounded domain. Assuming the whole state variable is available for feedback, we want to characterize all state feedback gain matrices $F$ such that with $u=F x$, the uncertain closed-loop system is quadratic stable, that is

$$
\begin{equation*}
F \in \mathcal{K}_{Q} \tag{7.13}
\end{equation*}
$$

To this end, we need to consider again the $p \times p$ matrices $W$ and $Q_{c}$ with $p:=n+m$, partitioned as (recall Section 6.2.1)

$$
W:=\left[\begin{array}{ll}
W_{1} & W_{2}  \tag{7.14}\\
W_{2}^{\prime} & W_{3}
\end{array}\right], \quad Q_{c}:=\left[\begin{array}{cl}
Q_{1 c} & 0 \\
0 & 0
\end{array}\right]
$$

where in both matrices, the $(1,1)$ block has dimension $n \times n$. We also consider, for $i=1,2, \cdots, N$ the convex sets

$$
\begin{equation*}
\mathcal{C}_{c i}:=\left\{W: W \geq 0, \quad v^{\prime} \Theta_{c i}(W) v \leq 0, \forall v \in \mathcal{N}_{c}\right\} \tag{7.15}
\end{equation*}
$$

where $\Theta_{c i}(W):=M_{c i} W+W M_{c i}^{\prime}+Q_{c}$. In other terms, the convex set defined in (7.15) is exactly the same as $\mathcal{C}_{c}$ but with matrix $M_{c}$ replaced by the extreme matrix $M_{c i}$. Of course, Theorem 6.1 remains valid for each extreme vertex of the polyhedral convex bounded domain $\mathcal{D}_{c}$. More interestingly, these convex sets can generate all quadratic stabilizing gains as is proved in the next theorem.

Theorem 7.1 Assume $Q_{1 c}$ is a positive definite matrix and consider the set

$$
\begin{equation*}
\mathcal{C}_{Q}:=\bigcap_{i=1}^{N} \mathcal{C}_{c i} \tag{7.16}
\end{equation*}
$$

The following hold
a) $\mathcal{C}_{Q}$ is a convex set.
b) Each $W \in \mathcal{C}_{Q}$ is such that $W_{1}>0$.
c) $\mathcal{K}_{Q}=\left\{W_{2}^{\prime} W_{1}^{-1}: W \in \mathcal{C}_{Q}\right\}$.

Proof The first two points follow immediately from Theorem 6.1 where it is proven that for each $i=1,2, \cdots, N$ the set $\mathcal{C}_{c i}$ is convex and each $W \in \mathcal{C}_{c i}$ is such that $W_{1}>0$. Consequently, point a) and point b) are both true from the definition of $\mathcal{C}_{Q}$ which is the intersection of the former convex sets. The proof of the last part of the theorem is much more involved.

Point c) It is done by construction. First take $F \in \mathcal{K}_{Q} \neq \emptyset$ and remember that $Q_{1 c}$ is positive definite. From Definition 7.1, there exists a symmetric positive definite matrix $P$ satisfying the matrix inequality

$$
\left(A+B_{2} F\right) P+P\left(A+B_{2} F\right)^{\prime}+Q_{1 c} \leq 0
$$

for all pairs $\left(A, B_{2}\right)$ such that $M_{c} \in \mathcal{D}_{c}$. Consequently, the same inequality holds true for the $N$ pairs $\left(A_{i}, B_{2 i}\right)$ corresponding to the extreme matrices $M_{c i}, i=1,2, \cdots, N$. Using this fact, it is readily seen that the matrix

$$
W=\left[\begin{array}{cc}
P & P F^{\prime} \\
F P & F P F^{\prime}
\end{array}\right] \in \mathcal{C}_{c i}
$$

for all $i=1,2, \cdots, N$ that is, $W \in \mathcal{C}_{Q}$ and $W_{2}^{\prime} W_{1}^{-1}=F P P^{-1}=F$. Conversely, pick $W \in \mathcal{C}_{Q} \neq \emptyset$ and notice that from (7.16) $W \in \mathcal{C}_{c i}$ for all $i=1,2, \cdots, N$. On the other hand, we know that a generic matrix $M_{c} \in \mathcal{D}_{c}$ can be written as

$$
M_{c}=\sum_{i=1}^{N} \xi_{i} M_{c i}, \quad \xi_{i} \geq 0, \quad \sum_{i=1}^{N} \xi_{i}=1
$$

which yields

$$
\begin{align*}
\Theta_{c}(W) & :=M_{c} W+W M_{c}^{\prime}+Q_{1 c} \\
& =\sum_{i=1}^{N} \xi_{i}\left(M_{c i} W+W M_{c i}^{\prime}+Q_{1 c}\right) \\
& =\sum_{i=1}^{N} \xi_{i} \Theta_{c i}(W) \tag{7.17}
\end{align*}
$$

The consequence is that for all $v \in \mathcal{N}_{c}$ and $\forall\left(A, B_{2}\right) \in \mathcal{D}$

$$
\begin{align*}
0 & \geq \sum_{i=1}^{N} \xi_{i} v^{\prime} \Theta_{c i}(W) v \\
& \geq v^{\prime} \Theta_{c}(W) v \\
& \geq x^{\prime}\left[\left(A+B_{2} W_{2}^{\prime} W_{1}^{-1}\right) W_{1}+W_{1}\left(A+B_{2} W_{2}^{\prime} W_{1}^{-1}\right)^{\prime}+Q_{1 c}\right] x \tag{7.18}
\end{align*}
$$

This inequality, together with the fact that both matrices $W_{1}$ and $Q_{1 c}$ are positive definite implies, by Definition 7.1, that $F=W_{2}^{\prime} W_{1}^{-1} \in \mathcal{K}_{Q}$ proving thus the theorem proposed since the case in which one of the sets is empty follows immediately.

This theorem is the uncertain systems stability counterpart of Theorem 6.1 proved in Chapter 6. The main point to be retained in mind is that it provides a convex description of the set of all quadratic stabilizing gains $\mathcal{K}_{Q}$ in terms of the same nonlinear mapping namely $W_{2}^{\prime} W_{1}^{-1}$ but with $W$ varying now in the convex set $\mathcal{C}_{Q}$. As discussed before, this theorem shows how to generate the set $\mathcal{K}_{Q}$ which is only a subset of $\mathcal{K}_{\mathcal{D}}$. Thus, it is clear that using it we can not capture all robust stabilizing gains as defined in Remark 7.1 but only those inside $\mathcal{K}_{Q}$ which are easily obtained by a simple convex feasibility problem.

Remark 7.3 From part c) of Theorem 7.1 it is evident that the pair $\left(A, B_{2}\right) \in \mathcal{D}$ is quadratic stabilizable if and only if the convex set $\mathcal{C}_{Q}$ is not empty. Indeed, if there are no matrices $P>0$ and $F$ satisfying Definition 7.1 then $\mathcal{C}_{Q}=\emptyset$.

Remark 7.4 The matrix function $\Theta_{c}(W)$ defined in (7.17) is of great importance. For $W$ fixed, it can be viewed as an affine function of $M_{c} \in \mathcal{D}_{c}$. Since the uncertain domain $\mathcal{D}_{c}$ is a polyhedral convex and bounded set, each matrix $M_{c}$ can be written as a convex combination (although unknown) of the extreme matrices $M_{c i}, i=1,2, \cdots, N$. Consequently, it follows that

$$
\Theta_{c}(W)=\sum_{i=1}^{N} \xi_{i} \Theta_{c i}(W)
$$

which means that $\Theta_{c}(W) \leq 0$ if and only if $\Theta_{c i}(W) \leq 0, i=1,2, \cdots, N$. In other words, we can say that the uncertain system defined by the pair $\left(A, B_{2}\right)$ is quadratic stabilizable if and only if the collection of $N$ systems represented by the pairs $\left(A_{i}, B_{2 i}\right)$ is quadratic stabilizable. The quadratic stabilizability of any other pair $\left(A, B_{2}\right) \in \mathcal{D}_{c}$ is a mere consequence of the quadratic stabilizability of the extreme matrices which define the domain $\mathcal{D}_{c}$. Even though the set $\mathcal{D}_{c}$ is composed by an infinity number of matrices $M_{c}$, only $N$ of them have to be used to check quadratic stabilizability.

To put in evidence the importance of Definition 7.1, let us proceed in the above discussion under the assumption that the collection of extreme matrices, namely ( $A_{i}, B_{2 i}$ ), $i=$ $1,2, \cdots, N$ is robustly stabilizable (recall Remark 7.1). That is, there exists a matrix $F$
such that the eigenvalues of $A_{i}+B_{2 i} F$ are all in the open left complex plane. From the Extended Lyapunov lemma, this is equivalent to the existence of matrices $P_{i}=P_{i}^{\prime}>0$ and $Q_{i}=Q_{i}^{\prime}>0$ such that

$$
0 \geq\left(A_{i}+B_{2 i} F\right) P_{i}+P_{i}\left(A_{i}+B_{2 i} F\right)^{\prime}+Q_{i}
$$

for all $i=1,2, \cdots, N$. Since the right hand side of the inequality

$$
0 \geq\left(A+B_{2} F\right) P+P\left(A+B_{2} F\right)^{\prime}+Q
$$

is not a convex function of the unknown matrices $A, B_{2}, P$ and $Q$ then the robust stabilizability of the collection of the extreme matrices does not provide enough information to guarantee the same for all matrices in the convex domain $\mathcal{D}_{c}$.

Remark 7.5 (Comparison between $\mathcal{D}_{n}$ and $\mathcal{D}_{c}$ ) In many instances (recall Example 7.1) the uncertainty of a linear dynamic system can be modeled as $M_{c} \in \mathcal{D}_{n}$ only at expense to impose to matrix $\Omega$ a prespecified structure. In Example 7.1 this occurred since $\Omega$ was taken to be a diagonal matrix. Let us analyze, in this particular case, the degree of conservativeness introduced. From the definition of $\mathcal{D}_{n}$ in Remark 7.1, we notice that the nominal matrices must be partitioned as

$$
M_{\mathrm{c} n}=\left[\begin{array}{cc}
A_{0} & B_{20} \\
0 & 0
\end{array}\right], \quad B_{n}=\left[\begin{array}{c}
B_{1} \\
0
\end{array}\right], \quad C_{n}=\left[\begin{array}{ll}
C_{1} & D_{12}
\end{array}\right]
$$

The consequence is that there exists a stabilizing gain $F$ such that with $u=F x$ the closed loop system is stable for all $M_{c} \in \mathcal{D}_{n}$ if $A_{0}+B_{20} F$ is stable and

$$
\left\|\left(C_{1}+D_{12} F\right)\left[s I-\left(A_{0}+B_{20} F\right)\right]^{-1} B_{1}\right\|_{\infty}<1
$$

Of course, if this is true then the closed-loop system is stable for all $\Omega$ such that $\|\Omega\| \leq 1$. In other words, this condition does not take into account the known fact that $\Omega$ is a diagonal matrix.

On the contrary, for this particular case (in which $\Omega$ is diagonal), the same uncertain domain can also be exactly described by $\mathcal{D}_{c}$ and consequently, applying Theorem 7.1 the whole set $\mathcal{K}_{Q}$ is generated which contributes to decrease the degree of conservativeness involved. This feature is numerically addressed in the example that follows.

Example 7.2 Consider again the linear uncertain system (recall Example 7.1)

$$
\dot{x}=A x+B_{2} u
$$

where matrices $A$ and $B_{2}$ are given by

$$
A=\left[\begin{array}{cc}
0 & \alpha-1 \\
\beta & 0
\end{array}\right], \quad B_{2}=\left[\begin{array}{c}
\alpha \\
1-\beta
\end{array}\right]
$$

and the parameters $\alpha$ and $\beta$, representing the uncertainties are such that

$$
|\alpha-0.5| \leq \gamma, \quad|\beta-0.5| \leq \gamma
$$

Letting the nominal value for these parameters be $\alpha_{0}=\beta_{0}=0.5$, the norm bounded domain $\mathcal{D}_{n}$ is completely defined by matrices (recall Remark 7.5)

$$
\begin{gathered}
A_{0}=\left[\begin{array}{rr}
0 & -0.5 \\
0.5 & 0
\end{array}\right], B_{20}=\left[\begin{array}{l}
0.5 \\
0.5
\end{array}\right] \\
B_{1}=\sqrt{2}\left[\begin{array}{cc}
\gamma & 0 \\
0 & \gamma
\end{array}\right], \quad C_{1}=\frac{1}{\sqrt{2}}\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right], \quad D_{12}=\frac{1}{\sqrt{2}}\left[\begin{array}{r}
1 \\
-1
\end{array}\right]
\end{gathered}
$$



Figure 7.3: Uncertain system closed-loop poles
where matrix $D_{12}$ is such that $D_{12}^{\prime} D_{12}=1$. In addition, matrix $\Omega$ such that $\|\Omega\| \leq 1$ exhibits the diagonal structure

$$
\Omega=\operatorname{diag}\left\{\frac{\alpha-0.5}{\gamma}, \frac{\beta-0.5}{\gamma}\right\}
$$

The central controller which satisfies the $R H_{\infty}$ constraint just discussed in the previous remark is provided by the solution (if any) of the Full information problem, that is $F=$ $F_{\infty}=-B_{20}^{\prime} P_{\infty}-D_{12}^{\prime} C_{1}$ where $P_{\infty}$ is the positive semidefinite and stabilizing solution of the Riccati equation (in the unknown $P$ )

$$
0=P A_{c 0}+A_{c 0}^{\prime} P-P\left(B_{20} B_{20}^{\prime}-B_{1} B_{1}^{\prime}\right) P+C_{1 c}^{\prime} C_{1 c}
$$

with $A_{c 0}:=A_{0}-B_{20} D_{12}^{\prime} C_{1}$ and $C_{1 c}:=\left(I-D_{12} D_{12}^{\prime}\right) C_{1}$. Numerically, it is possible to verify that the above Riccati equation admits a positive semidefinite and stabilizing solution for $\gamma \in[0,0.27]$. Thus, setting $\gamma=0.27$, we can conclude that the associated state feedback gain $F=F_{\infty}=[3.80-21.49]$ is stabilizing for all $M_{c} \in \mathcal{D}_{n}$.

On the other hand, the same system can also be modeled by means of a convex bounded domain $\mathcal{D}_{c}$. This is done with $N=4$ and the corresponding extreme matrices, whose entries, depending on the parameter $\gamma$, are easily determined. From Theorem 7.1, we observe that the convex set $\mathcal{C}_{Q}$ depends on $\gamma$. Numerically it is verified that with $Q_{1 c}=I$, the set $\mathcal{C}_{Q}$ is not empty for all $\gamma \in[0,0.36]$. The associated quadratic stabilizing gains, valid for all $M_{c} \in \mathcal{D}_{c}$ have been computed in two situations, namely

$$
\begin{aligned}
& \gamma=0.27 \Longrightarrow F=\left[\begin{array}{ll}
-0.13 & -1.91
\end{array}\right] \\
& \gamma=0.36 \Longrightarrow F=\left[\begin{array}{ll}
-0.26 & -4.94
\end{array}\right]
\end{aligned}
$$

As we have said before, in this case, Theorem 7.1 provides a better result when compared to the former one. Figure 7.3 shows the root locus of the uncertain closed-loop system with the quadratic stabilizing gain above, corresponding to $\gamma=0.36$ which is obtained with the parameters $(\alpha, \beta)$ varying in the rectangle $[0.14,0.86] \times[0.14,0.86]$. It is numerically confirmed that the closed-loop system is in the limit of stability since there exists a closedloop pole with real part equal to -0.0058 . Hence a small increase in $\gamma$ will move some closed-loop poles to the right part of the complex plane, causing instability.

### 7.2.2 $\quad \mathrm{RH}_{2}$ guaranteed cost control

Let the uncertain linear system under consideration be

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{7.19}\\
z & =C_{1} x+D_{12} u  \tag{7.20}\\
y & =x \tag{7.21}
\end{align*}
$$

and represented in fig. 7.1. It is assumed that i) the pair $\left(A, B_{2}\right)$ is quadratic stabilizable and ii) $D_{12}^{\prime} D_{12}=I$. For this system, our interest is to determine a $R H_{2}$ guaranteed cost control which is defined as follows. Choosing

$$
\begin{equation*}
\mathcal{G}(T(z, w ; s)):=\|T(z, w ; s)\|_{2}^{2} \tag{7.22}
\end{equation*}
$$

we have to determine an upper bound, depending only on $F \in \mathcal{K}_{Q}$ such that

$$
\begin{equation*}
\|T(z, w ; s)\|_{2}^{2}=\mathcal{J}\left(M_{c}, F\right) \leq \bar{\rho}(F), \quad \forall M_{c} \in \mathcal{D}_{c} \tag{7.23}
\end{equation*}
$$

and solve the optimal state feedback control problem

$$
\begin{equation*}
\bar{\rho}_{Q}:=\min _{F \in \mathcal{K}_{Q}} \bar{\rho}(F) \tag{7.24}
\end{equation*}
$$

Before we proceed, it is important to keep in mind that our main objective is to convert the design problem (7.24) into an equivalent convex programming problem. This fact of course has to be kept in mind during the determination of the upper bound satisfying inequality (7.23). To this end we need to introduce the symmetric matrices

$$
R_{c}:=\left[\begin{array}{c}
C_{1}^{\prime}  \tag{7.25}\\
D_{12}^{\prime}
\end{array}\right]\left[\begin{array}{ll}
C_{1} & D_{12}
\end{array}\right], \quad Q_{c}:=\left[\begin{array}{cc}
B_{1} B_{1}^{\prime} & 0 \\
0 & 0
\end{array}\right]
$$

which together with the next lemma define the function $\bar{\rho}(\cdot)$ which exhibits the properties just discussed.

Lemma 7.1 Assume $B_{1} B_{1}^{\prime}$ is a positive definite matrix, consider $F \in \mathcal{K}_{Q}$ arbitrary and define the convex set

$$
\overline{\mathcal{C}}_{Q}(F):=\mathcal{C}_{Q} \bigcap\left\{W: W_{2}^{\prime}=F W_{1}\right\}
$$

then

$$
\begin{equation*}
\bar{\rho}(F):=\min \left\{\operatorname{trace}\left\{R_{c} W\right]: W \in \overline{\mathcal{C}}_{Q}(F)\right\} \tag{7.26}
\end{equation*}
$$

is a valid upper bound to $\|T(z, w ; s)\|_{2}^{2}$ for all feasible $M_{c} \in \mathcal{D}_{c}$.
Proof First notice that for any $F \in \mathcal{K}_{Q}$ given, the set $\overline{\mathcal{C}}_{Q}(F)$ is convex and nonempty. The convexity follows from the fact that it is the intersection of a convex set (recall Theorem 7.1) with another one defined by a linear and hence convex constraint. Again from Theorem 7.1, $W \in \mathcal{C}_{Q}$ generates all quadratic stabilizing state feedback gains and so for any $F \in \mathcal{K}_{Q}$ fixed the linear constraint $W_{2}^{\prime}=F W_{1}$ must be satisfied for some $W \in \mathcal{C}_{Q}$. Let us now consider an arbitrary $M_{c} \in \mathcal{D}_{c}$ which can be written as a convex combination of the extreme matrices, that is

$$
M_{c}=\sum_{i=1}^{N} \xi_{i} M_{c i}, \quad \xi_{i} \geq 0, \quad \sum_{i=1}^{N} \xi_{i}=1
$$

and pick $W \in \overline{\mathcal{C}}_{Q}(F)$. For all $v \in \mathcal{N}_{c}$ we have

$$
\begin{aligned}
0 & \geq \sum_{i=1}^{N} \xi_{i} v^{\prime} \Theta_{c i}(W) v \\
& \geq x^{\prime}\left[\left(A+B_{2} F\right) W_{1}+W_{1}\left(A+B_{2} F\right)^{\prime}+B_{1} B_{1}^{\prime}\right] x
\end{aligned}
$$

which enables us to conclude that $W_{1} \geq P_{r}$ where $P_{r}>0$ solves the linear matrix equation

$$
\left(A+B_{2} F\right) P_{r}+P_{r}\left(A+B_{2} F\right)+B_{1} B_{1}^{\prime}=0
$$

On the other hand, applying the Schur complement formula to $W \geq 0$ it is simple to see that $W_{3} \geq F W_{1} F^{\prime}$. Then, for the same matrix $M_{c}$ and using the fact that $W_{2}^{\prime}=F W_{1}$, we also get

$$
\begin{aligned}
\operatorname{trace}\left[R_{c} W\right]= & \operatorname{trace}\left[\left[\begin{array}{ll}
C_{1} & D_{12}
\end{array}\right] W\left[\begin{array}{c}
C_{1}^{\prime} \\
D_{12}^{\prime}
\end{array}\right]\right] \\
= & \operatorname{trace}\left[\left(C_{1}+D_{12} F\right) W_{1}\left(C_{1}+D_{12} F\right)^{\prime}\right]+ \\
& +\operatorname{trace}\left[D_{12}\left(W_{3}-F W_{1} F^{\prime}\right) D_{12}^{\prime}\right] \\
\geq \geq & \operatorname{trace}\left[\left(C_{1}+D_{12} F\right) P_{r}\left(C_{1}+D_{12} F\right)^{\prime}\right] \\
\geq & \|T(z, w ; s)\|_{2}^{2}
\end{aligned}
$$

Since this inequality holds for an arbitrary $M_{c} \in \mathcal{D}_{c}$ and for all feasible $W \in \overline{\mathcal{C}}_{Q}(F)$, the final conclusion is that $\bar{\rho}(F)$ defined in (7.26) satisfies

$$
\|T(z, w ; s)\|_{2}^{2} \leq \bar{\rho}(F), \quad \forall M_{c} \in \mathcal{D}_{c}
$$

which proves the lemma proposed.
Remark 7.6 There are two reasons to define $\bar{\rho}(F)$ as indicated in (7.26). First, any feasible $W$ produces an upper bound to the $R H_{2}$ norm of the transfer function under consideration but the optimal solution of the convex problem indicated in the previous lemma provides the smaller upper bound as far as the choice of matrix $W$ is concerned. Second, when $N=1$, that is when $\mathcal{D}_{c}$ is composed by only on matrix $M_{c}$ (there is no uncertainty at all) then from the proof above it is easily verified that the optimal solution of problem (7.26) is

$$
W=\left[\begin{array}{cc}
P_{r} & P_{r} F^{\prime} \\
F P_{r} & F P_{r} F^{\prime}
\end{array}\right] \in \overline{\mathcal{C}}_{Q}(F)
$$

providing thus

$$
\bar{\rho}(F)=\|T(z, w ; s)\|_{2}^{2}
$$

In this special case, the function $\bar{\rho}(F)$ reduces to the $R H_{2}$ norm of the closed-loop transfer function. Thus, in the general case of uncertain systems with polyhedral convex bounded domain $\mathcal{D}_{c}$, this function can be interpreted as a generalized $\mathrm{RH}_{2}$ norm which is finite for all quadratic stabilizing state feedback gains.

We are now in position to solve the basic design problem (7.24) by means of an equivalent convex programming problem. This result is summarized in the next theorem.

Theorem 7.2 Assume $B_{1} B_{1}^{\prime}$ is a positive definite matrix. Then,

$$
\begin{equation*}
\min _{F \in \mathcal{K}_{Q}} \bar{\rho}(F)=\min \left\{\operatorname{trace}\left[R_{c} W\right]: W \in \mathcal{C}_{Q}\right\} \tag{7.27}
\end{equation*}
$$

Furthermore, being $\bar{F}$ and $\bar{W}$ the optimal solution of each problem they are related one to each other by $\bar{F}=\bar{W}_{2}^{\prime} \bar{W}_{1}^{-1}$.

Proof The proof is a mere consequence of Theorem 7.1 together with Lemma 7.1. Actually, the minimum of $\bar{\rho}(F)$ over $F \in \mathcal{K}_{Q}$ is obtained by the joint minimization of (7.26) with respect to both $F$ and $W$. Due to the part c) of Theorem 7.1, this joint minimization can be calculated in two independent steps. In the first one, the optimal solution $\bar{W}$ of the convex programming problem in the right hand side of (7.27) is determined. This solution allows the determination of the optimal gain namely, $\bar{F}=\bar{W}_{2}^{\prime} \bar{W}_{1}^{-1} \in \mathcal{K}_{Q}$ in the second and last step.

Remark 7.7 Assuming $B_{1} B_{1}^{\prime}>0$, the function $\bar{\rho}(F)$ introduced in Lemma 7.1 can also be calculated as follows

$$
\bar{\rho}(F)=\min _{P>0}\left\{\operatorname{trace}\left[C_{c c} P C_{c c}^{\prime}\right]: 0 \geq A_{c i} P+P A_{c i}^{\prime}+B_{1} B_{1}^{\prime}\right\}
$$

where $A_{c i}:=A_{i}+B_{2 i} F$ for $i=1,2, \cdots, N$ and $C_{c c}:=C_{1}+D_{12} F$. Since (recall Definition 7.1) the constraint set in the above problem is feasible if and only if $F \in \mathcal{K}_{Q}$ then it is simple to see that

$$
\bar{\rho}_{Q}=\min _{F, P>0}\left\{\operatorname{trace}\left[C_{c c} P C_{c c}^{\prime}\right]: 0 \geq A_{c i} P+P A_{c i}^{\prime}+B_{1} B_{1}^{\prime}\right\}
$$

Defining $\Lambda_{i}=\Lambda_{i}^{\prime} \geq 0$ as being the Kuhn-Tucker multipliers associated to the $i-t h$ inequality constraint, the Lagrangian function turns out to be

$$
\mathcal{L}=\operatorname{trace}\left[C_{c c} P C_{c c}^{\prime}\right]+\sum_{i=1}^{N} \operatorname{trace}\left[\Lambda_{i}\left(A_{c i} P+P A_{c i}^{\prime}+B_{1} B_{1}^{\prime}\right)\right]
$$

from which the Kuhn-Tucker necessary conditions for optimality can be written in terms of the unknown matrices $F, P>0, Q_{i} \geq 0$ and $\Lambda_{i} \geq 0$ satisfying

$$
\begin{aligned}
& 0 \geq A_{c i} P+P A_{c i}^{\prime}+B_{1} B_{1}^{\prime}:=-Q_{i} \\
& 0=\sum_{i=1}^{N}\left(A_{c i}^{\prime} \Lambda_{i}+\Lambda_{i} A_{c i}\right)+C_{c c}^{\prime} C_{c c} \\
& 0=\left(F+D_{12}^{\prime} C_{1}+\sum_{i=1}^{N} B_{2 i}^{\prime} \Lambda_{i}\right) P \\
& 0=\operatorname{trace}\left[\Lambda_{i} Q_{i}\right]
\end{aligned}
$$

These conditions put in evidence several nonlinear relationship among the unknown matrices implying that they can be solved only in some special cases. Indeed, the third equation can be used to determine the optimal gain matrix $F$ explicitly, however the second condition alone is not of great help in finding matrices $\Lambda_{i}, i=1,2, \cdots, N$ since it represents only one equation with $N$ unknown. The consequence is that putting aside the third equation, the remaining ones have to be solved simultaneously. This is accomplished by solving the right hand side of (7.27) directly. It is important to stress that this discussion is no longer true in the particular case $N=1$. The domain $\mathcal{D}_{c}$ degenerates to a single matrix and the above conditions reduce to those of Remark 6.7 with matrices $P$ and $\Lambda$ replaced one by the other since the problem now under consideration is exactly the dual of the one treated there.

Remark 7.8 The use of any convex programming method to get the global solution of the optimal $\mathrm{RH}_{2}$ guaranteed cost control problem

$$
\bar{\rho}_{Q}=\min \left\{\operatorname{trace}\left[R_{c} W\right]: W \in \mathcal{C}_{Q}\right\}
$$

depends on the boundedness of the following convex set (recall Remark 6.8)

$$
\mathcal{C}_{Q \mu}:=\mathcal{C}_{Q} \bigcap\left\{W: \operatorname{trace}\left[R_{c} W\right] \leq \mu\right\}
$$

Since the set $\mathcal{C}_{Q}$ is itself an intersection of $N$ convex sets, then the requirement that there exists an index $1 \leq i \leq N$ for which the pair ( $-A_{c i}, C_{1 c}$ ) is detectable is a sufficient condition for the boundedness of $\mathcal{C}_{Q \mu}$ for all $\mu$ such that $\bar{\rho}_{Q}<\mu<\infty$.

Remark 7.9 In this section we have defined the generalized norm $\bar{\rho}(F)$ satisfying

$$
\|T(z, w ; s)\|_{2}^{2} \leq \bar{\rho}(F), \quad \forall M_{c} \in \mathcal{D}_{c}
$$

Let us now verify the difficulties we have to face if we desire to work, instead of $\bar{\rho}(F)$, with the true generalized $\mathrm{RH}_{2}$ norm

$$
\rho(F)=\max _{M_{c} \in \mathcal{D}_{c}}\|T(z, w ; s)\|_{2}^{2}
$$

Since $F \in \mathcal{K}_{Q}$ assures that all closed-loop transfer functions $T(z, w ; s)$ are stable then with $A_{c c}:=A+B_{2} F$ and $C_{c c}:=C_{1}+D_{12} F$ we have

$$
\rho(F)=\max _{M_{c} \in \mathcal{D}_{c}} \operatorname{trace}\left[\int_{0}^{\infty} C_{c c} e^{A_{c c} t} B_{1} B_{1}^{\prime} e^{A_{c c}^{\prime} t} C_{c c}^{\prime} d t\right]
$$

On the other hand, from Theorem 7.1 we already know that all $F \in \mathcal{K}_{Q}$ are generated by $F=W_{2}^{\prime} W_{1}^{-1}$ with $W \in \mathcal{C}_{Q}$. Then we get

$$
\begin{aligned}
\rho_{Q} & =\min _{F \in \mathcal{K}_{Q}} \rho(F) \\
& =\min _{W \in \mathcal{C}_{Q}} \rho\left(W_{2}^{\prime} W_{1}^{-1}\right)
\end{aligned}
$$

In the above problem the constraint $W \in \mathcal{C}_{Q}$ is convex. However, the difficulty is that it is not easy to prove if the objective function is convex or not. In the affirmative case, no major difficulty exists to approach the problem by means of convex programming methods. The convexity of $\rho(\cdot)$ remains until now an open question.

Remark 7.10 Let us now consider the $\mathrm{RH}_{2}$ guaranteed cost control problem but with norm bounded uncertainty, that is $M_{c} \in \mathcal{D}_{n}$. The uncertain system (7.19) - (7.21) is now written with a slightly different notation

$$
\begin{aligned}
\dot{x} & =A x+B_{0} w+B_{2} u \\
z_{0} & =C_{0} x+D_{0} u \\
y & =x
\end{aligned}
$$

where $B_{0} B_{0}^{\prime}>0$ and $D_{0}^{\prime} D_{0}=I$. Following Remarks 7.2 and 7.5 , using $u=F x$ the closedloop system matrix can be expressed as

$$
A+B_{2} F=A_{c c}+B_{1} \Omega C_{c c}, \quad\|\Omega\| \leq 1
$$

where $A_{c c}:=A_{0}+B_{20} F$ and $C_{c c}:=C_{1}+D_{12} F$. Notice that in this setting, the external perturbation, and the parametric uncertainty enter in the system through two different matrices, namely $B_{0}$ and $B_{1}$ respectively. First of all, it is claimed that for the above uncertain system, the upper bound $\bar{\rho}(F)$ is given by

$$
\begin{aligned}
\bar{\rho}(F) & =\min _{\beta \in \mathcal{B}} \operatorname{trace}\left[B_{0}^{\prime} P(\beta) B_{0}\right] \\
& \geq\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}, \quad \forall M_{c} \in \mathcal{D}_{n}
\end{aligned}
$$

where $\mathcal{B}$ is the set of all $\beta>0$ such that there exists a symmetric, stabilizing and positive semidefinite solution $P(\beta)$ to the Riccati equation (in the unknown $P$ )

$$
0=P A_{c c}+A_{c c}^{\prime} P+\beta P B_{1} B_{1}^{\prime} P+\beta^{-1} C_{c c}^{\prime} C_{c c}+C_{0 c}^{\prime} C_{0 c}
$$

where $C_{0 c}:=C_{0}+D_{0} F$. To show this, notice first that any $\beta>0$ yields

$$
\begin{aligned}
C_{c c}^{\prime} \Omega^{\prime} B_{1}^{\prime} P+P B_{1} \Omega C_{c c}= & \beta^{-1} C_{c c}^{\prime} \Omega^{\prime} \Omega C_{c c}+\beta P B_{1}^{\prime} B_{1} P- \\
& -\left(\Omega C_{c c} \beta^{-1 / 2}-B_{1}^{\prime} P \beta^{1 / 2}\right)^{\prime}\left(\Omega C_{c c} \beta^{-1 / 2}-B_{1}^{\prime} P \beta^{1 / 2}\right) \\
\leq & \beta^{-1} C_{c c}^{\prime} C_{c c}+\beta P B_{1}^{\prime} B_{1} P
\end{aligned}
$$

where to get the last inequality we have used the fact that $\Omega^{\prime} \Omega \leq I$. Using this inequality together with the previously defined Riccati equation, for all $\beta \in \mathcal{B}$ we can draw two important conclusions. The first one is that (recall Remark 7.5)

$$
\left\|C_{c c}\left[s I-\left(A_{0}+B_{20} F\right)\right]^{-1} B_{1}\right\|_{\infty}<1
$$

meaning that the closed-loop system matrix $A+B_{2} F$ is stable for all $M_{c} \in \mathcal{D}_{n}$. Second, from the above inequality we also have

$$
\left(A+B_{2} F\right)^{\prime} P(\beta)+P(\beta)\left(A+B_{2} F\right)+C_{0 c}^{\prime} C_{0 c} \leq 0
$$

for all $M_{c} \in \mathcal{D}_{n}$ implying that

$$
\begin{aligned}
\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2} & =\operatorname{trace}\left[\int_{0}^{\infty} B_{0}^{\prime} e^{\left(A+B_{2} F\right)^{\prime} t} C_{0 c}^{\prime} C_{0 c} e^{\left(A+B_{2} F\right) t} B_{0} d t\right] \\
& \leq \operatorname{trace}\left[B_{0}^{\prime} P(\beta) B_{0}\right] \\
& \leq \bar{\rho}(F)
\end{aligned}
$$

Based on this result, let us now determine the best guaranteed cost denoted $\bar{\rho}_{N}$, that is

$$
\bar{\rho}_{N}=\min _{F, \beta \in \mathcal{B}} \operatorname{trace}\left[B_{0}^{\prime} P(\beta) B_{0}\right]
$$

The necessary conditions for optimality with respect to $F$ is readily obtained (recall Remark 6.19) from the associated Lagrangian function. After tedious algebraic manipulations, a solution to these conditions can be written as

$$
F(\beta):=\tilde{F}-(1+\beta)^{-1} B_{20}^{\prime} W(\beta)
$$

where

$$
\tilde{F}:=-(1+\beta)^{-1}\left(D_{12}^{\prime} C_{1}+\beta D_{0}^{\prime} C_{0}\right)
$$

and $W(\beta):=\beta P(\beta)$ solves the Riccati equation (in the unknown W )

$$
0=W \tilde{A}+\tilde{A}^{\prime} W+W \tilde{N} W+\tilde{M}
$$

with

$$
\begin{aligned}
& \tilde{A}:=A_{0}+B_{20} \tilde{F} \\
& \tilde{N}:=B_{1} B_{1}^{\prime}-(1+\beta)^{-1} B_{20} B_{20}^{\prime} \\
& \tilde{M}:=\left(C_{1}+D_{12} \tilde{F}\right)^{\prime}\left(C_{1}+D_{12} \tilde{F}\right)+\beta\left(C_{0}+D_{0} \tilde{F}\right)^{\prime}\left(C_{0}+D_{0} \tilde{F}\right)
\end{aligned}
$$

Finally, using this result we get

$$
\bar{\rho}_{N}=\min _{\beta \in \mathcal{B}} \beta^{-1} \operatorname{trace}\left[B_{0}^{\prime} W(\beta) B_{0}\right]
$$

which is numerically easy to solve since it involves the search of only one positive parameter. The main conclusion to be kept in mind is that the guaranteed cost associated to the uncertain domain $\mathcal{D}_{n}$ is numerically tractable by means of a unidimensional search procedure. It must be clear that in the present case, it may occur that $\mathcal{B}=\emptyset$, meaning that $\bar{\rho}_{N}$ is unbounded and so no guaranteed cost control is obtained. This aspect is illustrated in the next example.


Figure 7.4: Actual value of the $\mathrm{RH}_{2}$ norm

Example 7.3 Let the uncertain system be again the one of Example 7.2. The state space representation is

$$
\begin{aligned}
\dot{x} & =A x+B_{0} w+B_{2} u \\
z_{0} & =C_{0} x+D_{0} u \\
y & =x
\end{aligned}
$$

where the uncertainty is modeled as $M_{c} \in \mathcal{D}_{c}$ as well as $M_{c} \in \mathcal{D}_{n}$. The main goal is to determine the associated guaranteed $\mathrm{RH}_{2}$ cost control. The data are those of Example 7.2 but with only one uncertain parameter obtained by imposing $\alpha=\beta$ such that

$$
|\alpha-0.5| \leq \gamma
$$

where $\gamma=0.20$. The controlled output is defined by matrices

$$
B_{0}=B_{1}=\sqrt{2}\left[\begin{array}{ll}
\gamma & 0 \\
0 & \gamma
\end{array}\right], C_{0}=\left[\begin{array}{ll}
1 & 1 \\
0 & 0
\end{array}\right], \quad D_{0}=\left[\begin{array}{l}
0 \\
1
\end{array}\right]
$$

Three different situations have been considered :

1) First, assuming no uncertainty acts on the open-loop system, the Full information problem in $\mathrm{RH}_{2}$ has been solved using the nominal data. The optimal solution found is

$$
\min \left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}=0.32 \Longrightarrow F_{2}=\left[\begin{array}{ll}
-1.00 & -1.00
\end{array}\right]
$$

Using this state feedback gain, it is simple to verify that the closed-loop system remains stable for all $\alpha$ varying in the interval defined as above with $\gamma=(\sqrt{5}-2) / 2=0.1180$. This fact is illustrated in fig. 7.4 where it is clear that the corresponding value of $\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}$, denoted by the symbol " $F_{2}$ ", becomes arbitrarily large as $\alpha$ goes to 0.3820 . For all other values of $\alpha$, the closed-loop system is unstable.
2) We considered $M_{c} \in \mathcal{D}_{c}$. In this case we have two extreme matrices only. The optimal solution of the convex programming problem in Theorem 7.2 provides

$$
\bar{\rho}_{Q}=1.18 \Longrightarrow F_{Q}=\bar{F}=\left[\begin{array}{ll}
-0.02 & -2.39
\end{array}\right]
$$

As indicated in fig.7.4, the closed-loop system is now stable in the whole interval $0.3 \leq \alpha \leq 0.7$ implying that the associated cost is always finite. This is an important
improvement of the guaranteed cost control design proposed when compared with the previous situation. It is interesting to notice that the closed-loop system transfer function satisfies

$$
\max _{0.3 \leq \alpha \leq 0.7}\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2}=0.77<\bar{\rho}_{Q}
$$

which illustrates the fact discussed before (recall Lemma 7.1) that function $\bar{\rho}(F)$ is only an upper bound to the true value of the minimized cost.
3) In this situation we modeled the uncertainty as $M_{c} \in \mathcal{D}_{n}$. Following the results of Remark 7.10, it has been verified that $\mathcal{B}=(0,0.7]$ in order to preserve $W(\beta)>0$. Using this, an unidimensional search has been implemented to get the optimal $\bar{\beta}=0.32$ which provides

$$
\bar{\rho}_{N}=2.45 \Longrightarrow F_{N}=F(\beta=0.32)=\left[\begin{array}{ll}
-0.13 & -2.52
\end{array}\right]
$$

Figure 7.4 shows that the closed-loop system is also stable for all values of $\alpha$ in the real interval $0.3 \leq \alpha \leq 0.7$. There is no big difference between the cost variation in the last two situations. However, comparing them, it is apparent that $\bar{\rho}_{Q}<\bar{\rho}_{N}$ which indicates that in this example the design based on $\mathcal{D}_{c}$ is preferable.
Figure 7.4 illustrates also the advantages of the guaranteed cost control design when compared with the nominal Full information control design. For the nominal system it is obvious that the latter is better. However, the performance imposed by the optimal Full information controller may become worse when a small change in the system parameters occurs. Even the closed-loop stability can be lost. On the contrary, the guaranteed cost controller assures closed-loop stability for the parameters varying in the prespecified range and impose a reduced level of performance deterioration.

### 7.2.3 $\quad R H_{\infty}$ guaranteed cost control

We now move our attention to the $R H_{\infty}$ counterpart of the results provided in the last section. The uncertain dynamic system under consideration is again the same depicted in fig. 7.1 where as indicated the full state vector is available for feedback. To ease the presentation the state space equations are given once again,

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{7.28}\\
z & =C_{1} x+D_{12} u  \tag{7.29}\\
y & =x \tag{7.30}
\end{align*}
$$

where $A$ and $B_{2}$ are uncertain matrices such that is $M_{c} \in \mathcal{D}_{c}$ and as before we assume that i) the pair ( $A, B_{2}$ ) is quadratic stabilizable and ii) $D_{12}^{\prime} D_{12}=I$. The above equations define a family of linear dynamic systems for which our goal is to determine an associated $R H_{\infty}$ guaranteed cost control. Following the results of Chapter 6, let $\gamma$ denotes an arbitrary positive scalar and choose

$$
\begin{equation*}
\mathcal{G}(T(z, w ; s)):=\|T(z, w ; s)\|_{\infty}^{2} \tag{7.31}
\end{equation*}
$$

Then the minimum guaranteed $R H_{\infty}$ cost can be calculated from (7.11), yielding

$$
\begin{equation*}
\rho_{Q}=\min _{\rho, F \in \mathcal{K}_{Q}}\left\{\rho: \mathcal{J}\left(M_{c}, F\right) \leq \rho, \quad \forall M_{c} \in \mathcal{D}_{c}\right\} \tag{7.32}
\end{equation*}
$$

Parallel to the case of $\mathrm{RH}_{2}$ guaranteed cost control problem, the solution of (7.32) is extremely difficult to get, then the strategy to be used is to define a suboptimal easy to determine solution provided by

$$
\begin{equation*}
\bar{\rho}_{Q}:=\min \left\{\gamma^{2}:(F, \gamma) \in \mathcal{K}_{\gamma Q}\right\} \geq \rho_{Q} \tag{7.33}
\end{equation*}
$$

where $\mathcal{K}_{\gamma Q}$ is a subset of the set of all pairs $(F, \gamma)$ such that $F \in \mathcal{K}_{Q}$ and

$$
\begin{equation*}
\|T(z, w ; s)\|_{\infty}^{2}=\mathcal{J}\left(M_{c}, F\right) \leq \gamma^{2}, \forall M_{c} \in \mathcal{D}_{c} \tag{7.34}
\end{equation*}
$$

It is clear that the choice of the set $\mathcal{K}_{\gamma Q}$ is crucial to get good suboptimal solutions in terms of inequality (7.33). Besides it must be possible to convert it into an equivalent convex set. To this end, considering matrix $W$ partitioned as indicated in (7.14) and matrices $R_{c}$ and $Q_{c}$ given in (7.25) let us define the following convex set (recall Theorem 6.5) associated with each extreme matrix of the polyhedral convex bounded uncertain domain $\mathcal{D}_{c}$,

$$
\begin{equation*}
\mathcal{C}_{\gamma i}:=\left\{(W, \rho): W \geq 0, \quad \rho>0, v^{\prime} \Theta_{\gamma i}(W, \rho) v \leq 0, \forall v \in \mathcal{N}_{c}\right\} \tag{7.35}
\end{equation*}
$$

where $\Theta_{\gamma i}(W, \rho):=\Theta_{c i}(W)+\rho^{-1} W R_{c} W$ for each $i=1,2, \cdots, N$. Thus, it is apparent that the above set has been obtained from $\mathcal{C}_{\gamma c}$ by simply replacing matrix $M_{c}$ by the extreme matrix $M_{c i}$. The next theorem provides the solution of the problem indicated in (7.33).
Theorem 7.3 Assume $B_{1} B_{1}^{\prime}$ is a positive definite matrix and consider the set

$$
\begin{equation*}
\mathcal{C}_{\gamma Q}:=\bigcap_{i=1}^{N} \mathcal{C}_{\gamma i} \tag{7.36}
\end{equation*}
$$

The following hold
a) $\mathcal{C}_{\gamma Q}$ is a convex set.
b) Each $(W, \rho) \in \mathcal{C}_{\gamma Q}$ is such that $W_{1}>0$.
c) The subset $\mathcal{K}_{\gamma Q}$ defined above can be chosen as

$$
\begin{equation*}
\mathcal{K}_{\gamma Q}=\left\{\left(W_{2}^{\prime} W_{1}^{-1}, \sqrt{\rho}\right):(W, \rho) \in \mathcal{C}_{\gamma Q}\right\} \tag{7.37}
\end{equation*}
$$

Proof The first two points are straightforward consequences of Theorem 6.5. Since, both are true for each set $\mathcal{C}_{\gamma i}, i=1,2, \cdots, N$ then the same is also true for their intersection.

Point c) We have to prove that any $(W, \rho) \in \mathcal{C}_{\gamma Q}$ provides $F=W_{2}^{\prime} W_{1}^{-1} \in \mathcal{K}_{Q}$ and $\gamma=\sqrt{\rho}$ such that (7.34) holds. This is proved by selecting an arbitrary pair $(W, \rho) \in \mathcal{C}_{\gamma Q}$ and noticing that $\Theta_{c i}(W) \leq \Theta_{\gamma i}(W, \rho)$ for all $i=1,2, \cdots, N$ implies

$$
(W, \rho) \in \mathcal{C}_{\gamma Q} \Longrightarrow W \in \mathcal{C}_{Q}
$$

which from part c) of Theorem 7.1 allows us to conclude that $F=W_{2}^{\prime} W_{1}^{-1} \in \mathcal{K}_{Q}$. Let us now pick an arbitrary matrix $M_{c} \in \mathcal{D}_{c}$. Recalling Theorem 6.5, for this matrix the set $\mathcal{C}_{\gamma c}$ is well defined. So, expressing $M_{c}$ as a convex combination of the extreme matrices

$$
M_{c}=\sum_{i=1}^{N} \xi_{i} M_{c i}, \quad \xi_{i} \geq 0, \quad \sum_{i=1}^{N} \xi_{i}=1
$$

yields (recall Theorem 6.5) for all $v \in \mathcal{N}_{c}$

$$
\begin{align*}
v^{\prime} \Theta_{\gamma c}(W, \rho) v & :=v^{\prime}\left[\Theta_{c}(W)+\rho^{-1} W R_{c} W\right] v \\
& =v^{\prime}\left[\sum_{i=1}^{N} \xi_{i}\left(\Theta_{c i}(W)+\rho^{-1} W R_{c} W\right)\right] v \\
& =\sum_{i=1}^{N} \xi_{i} v^{\prime} \Theta_{\gamma i}(W, \rho) v \leq 0 \tag{7.38}
\end{align*}
$$

which means that $(W, \rho) \in \mathcal{C}_{\gamma c}$. Since the inequality (7.38) is always true for all $M_{c} \in \mathcal{D}_{c}$ we get from part c) of Theorem 6.5 that

$$
(F, \gamma)=\left(W_{2}^{\prime} W_{1}^{-1}, \sqrt{\rho}\right) \in \bigcap_{M_{c} \in \mathcal{D}_{c}} \mathcal{K}_{\gamma c}
$$

which means that (7.34) actually holds.
The importance of this theorem is twofold. First, Problem (7.33) which defines the guaranteed $R H_{\infty}$ cost reduces to

$$
\begin{equation*}
\bar{\rho}_{Q}=\min \left\{\rho:(W, \rho) \in \mathcal{C}_{\gamma Q}\right\} \tag{7.39}
\end{equation*}
$$

and this is a convex programming problem which can be solved efficiently. Second, the convex set $\mathcal{C}_{\gamma Q}$ allows the determination of a set of feedback matrix gains such that for a given $\gamma>0$, the closed loop system is quadratically stable and $\|T(z, w ; s)\|_{\infty}$ is bounded above by $\gamma$ for all $M_{c} \in \mathcal{D}_{c}$.

Remark 7.11 The boundedness of the set $\mathcal{C}_{\gamma Q}$ can be analyzed following the same lines of Remark 6.22. The key property to assure the adequate use of convex programming methods in solving the approximate version of Problem (7.39), namely

$$
\bar{\rho}_{Q} \approx \min \left\{\rho:(W, \rho) \in \mathcal{C}_{\gamma Q}, \quad \operatorname{trace}\left[R_{C} W\right] \leq \beta\right\}
$$

where the scalar $\beta>0$ is finite and appropriately large, is to impose that there exists an index $1 \leq i \leq N$ for which the corresponding pair $\left(-A_{c i}, C_{1 c}\right)$ is detectable.

Remark 7.12 The convex set $\mathcal{C}_{\gamma Q}$ presents two additional properties. First, in the case of certain dynamic systems, the domain $\mathcal{D}_{c}$ is defined by only one matrix, then $\mathcal{C}_{\gamma Q}$ equals the set $\mathcal{C}_{\gamma c}$ introduced in Theorem 6.5. As a consequence the optimal value of $\bar{\rho}_{Q}$ reduces to the minimum possible $R H_{\infty}$ norm level. Second, in the case of uncertain systems, with $\gamma$ arbitrarily large $(W, \gamma) \in \mathcal{C}_{\gamma Q}$ implies $W \in \mathcal{C}_{Q}$. Once again, the set of all quadratic stabilizing feedback gains is obtained.

### 7.2.4 Miscellaneous design problems

This section is devoted to treat in the framework of uncertain systems control design the other problems already solved in Chapter 6 , namely Mixed $R H_{2} / R H_{\infty}$ control, $R H_{2}$ control with regional pole placement, time domain specifications and controllers with structural constraints. The basic idea is to generalize the previous results on guaranteed cost control to cope with the new performance constraints. All design problems are solved for the special case of static linear state feedback controllers.

Let us first consider the Mixed $R H_{2} / R H_{\infty}$ control problem associated to the linear dynamic system of fig. 7.1 and with the following state space representation

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{7.40}\\
z_{0} & =C_{0} x+D_{0} u  \tag{7.41}\\
z_{1} & =C_{1} x+D_{12} u  \tag{7.42}\\
y & =x \tag{7.43}
\end{align*}
$$

where it is assumed that the uncertainty is described by $M_{c} \in \mathcal{D}_{c}$. Furthermore, it is also assumed that i) the pair ( $A, B_{2}$ ) is quadratic stabilizable and ii) $D_{12}^{\prime} D_{12}=I$. The controller is given by $u=F x$ where matrix $F$ is to be determined. In the present
context, the guaranteed mixed $\mathrm{RH}_{2} / \mathrm{RH}_{\infty}$ cost control problem can be formulated as being to determine (if one exists) a state feedback gain matrix $F$ such that for a given scalar $\gamma>0$, the closed-loop system is quadratically stable and

$$
\begin{equation*}
\rho_{Q}:=\min \left\{\rho:\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2} \leq \rho,\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty} \leq \gamma, \forall M_{c} \in \mathcal{D}_{c}\right\} \tag{7.44}
\end{equation*}
$$

Based on the previous results, a suboptimal solution of this problem is given in the next theorem.

Theorem 7.4 Let $\gamma>0$ be given and let $\bar{W}$ be the optimal solution of the convex programming problem

$$
\begin{equation*}
\bar{\rho}_{Q}=\min \left\{\operatorname{trace}\left[R_{0} W\right]:\left(W, \gamma^{2}\right) \in \mathcal{C}_{\gamma Q}\right\} \tag{7.45}
\end{equation*}
$$

Then, $F=\bar{W}_{2}^{\prime} \bar{W}_{1}^{-1} \in \mathcal{K}_{Q}$ and $\bar{\rho}_{Q} \geq \rho_{Q}$.
Proof It follows immediately from the result of Theorem 7.2 together with the ones of Theorem 7.3. Indeed, from the last theorem $\left(\bar{W}, \gamma^{2}\right) \in \mathcal{C}_{\gamma Q}$ implies that $F=$ $\bar{W}_{2}^{\prime} \bar{W}_{1}^{-1} \in \mathcal{K}_{Q}$ and

$$
\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty} \leq \gamma, \forall M_{c} \in \mathcal{D}_{c}
$$

Furthermore, using the fact that

$$
\left(\bar{W}, \gamma^{2}\right) \in \mathcal{C}_{\gamma Q} \Longrightarrow \bar{W} \in \mathcal{C}_{Q}
$$

from Theorem 7.2 we conclude that

$$
\left\|T\left(z_{0}, w ; s\right)\right\|_{2}^{2} \leq \operatorname{trace}\left[R_{0} \bar{W}\right]=\bar{\rho}_{Q}, \forall M_{c} \in \mathcal{D}_{c}
$$

which proves the theorem proposed since $\left(F, \bar{\rho}_{Q}\right)$ is feasible for Problem (7.44).
Example 7.4 Let us consider the following uncertain linear system (recall Example 6.2)

$$
\begin{aligned}
\dot{x} & =A x+B_{1} w+B_{2} u \\
z_{0} & =C_{0} x+u \\
z_{1} & =C_{1} x+u \\
y & =x
\end{aligned}
$$

where

$$
\begin{aligned}
& A=\left[\begin{array}{cc}
0 & 1 \\
-1+\alpha & -1+2 \alpha
\end{array}\right], \quad B_{1}=\left[\begin{array}{l}
1 \\
0
\end{array}\right] \\
& B_{2}=\left[\begin{array}{c}
0 \\
1-\alpha / 2
\end{array}\right], \quad C_{0}=C_{1}=\left[\begin{array}{ll}
0 & -2
\end{array}\right]
\end{aligned}
$$

The only available information concerning the uncertain parameter $\alpha$ is that $0 \leq \alpha \leq 1$. We have first verified numerically that the set $\mathcal{C}_{\gamma Q}$ is not empty for $\gamma \geq 4.10$. Then we choose $\gamma=5.00$ for the calculations that follow. On the other hand, we observe that the dynamic system handled in Example 6.2 corresponds to the above one calculated for $\alpha=0$. This system is called nominal. For the nominal system we solved the Mixed $R H_{2} / R H_{\infty}$ control problem introduced in Theorem 6.8. The optimal solution provided the state feedback gain $F_{n o m i n a l}$ which has been used to calculate the closed-loop transfer function $\mathrm{RH}_{2}$ norm indicated in fig. 7.5. Simple calculation shows that the closed-loop system is internally stable only for $0 \leq \alpha<0.5$ which implies that $F_{\text {nominal }} \notin \mathcal{K}_{Q}$.

Moreover, solving the convex programming problem (7.45) we got

$$
\bar{\rho}_{Q}=4.83 \Longrightarrow F_{\text {guaranteed }}=\left[\begin{array}{ll}
-1.90 & -7.20
\end{array}\right]
$$



Figure 7.5: $R H_{2}$ norm and $\bar{\rho}_{Q}$

Now, the closed-loop is internally stable for all $\alpha$ in the prespecified interval. Figure 7.5 also shows the closed-loop system $\mathrm{RH}_{2}$ norm compared with the minimum upper bound available. It is interesting to observe that the behavior of the cost as a function of the uncertain parameter is almost constant in the interval considered. Furthermore, fig. 7.6 shows the magnitude of the transfer function $T\left(z_{1}, w ; j \omega\right)$ against frequency for several values of $\alpha$ in the given interval. From this figure we notice that the constraint $\left\|T\left(z_{1}, w ; s\right)\right\|_{\infty} \leq \gamma=5$ is binding and is always verified.

We now move our attention to the problem of uncertain system root clustering in a given circular region $\mathcal{R}$. It is stated as follows, given the uncertain linear system

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{7.46}\\
z & =C_{1} x+D_{12} u  \tag{7.47}\\
y & =x \tag{7.48}
\end{align*}
$$

satisfying the assumption i) $D_{12}^{\prime} D_{12}=I$, find (if one exists) an associated $R H_{2}$ guaranteed cost control such that in addition the closed-loop system poles with $u=$ $F x$ are all inside the circular region $\mathcal{R}$. More specifically, in the same framework considered before, we seek for the solution of the following optimization problem

$$
\begin{equation*}
\bar{\rho}_{Q}=\min \left\{\rho:(F, \rho) \in \mathcal{K}_{R Q}\right\} \tag{7.49}
\end{equation*}
$$

where $\mathcal{K}_{R Q}$ is a subset of the set of all pairs $(F, \rho)$ such that $F \in \mathcal{K}_{Q}$,

$$
\begin{equation*}
\|T(z, w ; s)\|_{2}^{2} \leq \rho, \quad \forall M_{\mathrm{c}} \in \mathcal{D}_{c} \tag{7.50}
\end{equation*}
$$

and the closed-loop system poles are all inside the circular region $\mathcal{R}$. As commented before, the key step towards the solution of problem (7.49) is the definition of the feasible set $\mathcal{K}_{R Q}$ (notice that now $\rho$ is a variable to be determined). To this end let us redefine the matrix

$$
M_{c}:=\left[\begin{array}{cc}
A+\alpha I & B_{2} \\
0 & 0
\end{array}\right]
$$



Figure 7.6: Bode diagram - magnitude
and accordingly the convex sets (recall Theorem 6.9) associated with each extreme matrix of the convex domain $\mathcal{D}_{c}$.

$$
\begin{equation*}
\mathcal{C}_{R i}:=\left\{W: W \geq 0, \quad v^{\prime} \Theta_{R i}(W) v \leq 0, \forall v \in \mathcal{N}_{c}\right\} \tag{7.51}
\end{equation*}
$$

where $\Theta_{R i}(W):=\Theta_{c i}(W)+r^{-1} M_{c i} W M_{c i}^{\prime}$ for each $i=1,2, \cdots, N$. In other words, each set above is the same as $\mathcal{C}_{R}$ with matrix $M_{c}$ replaced by the extreme matrix $M_{c i}$.

Theorem 7.5 Assume $B_{1} B_{1}^{\prime}$ is a positive definite matrix, let the circular region $\mathcal{R}$ be given and consider the set

$$
\begin{equation*}
\mathcal{C}_{R Q}:=\bigcap_{i=1}^{N} \mathcal{C}_{R i} \tag{7.52}
\end{equation*}
$$

The following hold
a) $\mathcal{C}_{R Q}$ is a convex set.
b) Each $W \in \mathcal{C}_{R Q}$ is such that $W_{1}>0$.
c) The subset $\mathcal{K}_{R Q}$ defined above can be chosen as

$$
\begin{equation*}
\mathcal{K}_{R Q}=\left\{\left(W_{2}^{\prime} W_{1}^{-1}, \operatorname{trace}\left[R_{c} W\right]\right): W \in \mathcal{C}_{R Q}\right\} \tag{7.53}
\end{equation*}
$$

Proof The first two points are immediate consequences of Theorem 6.9. Both points are valid for the intersection defining $\mathcal{C}_{R Q}$.

Point c) It is apparent that $\mathcal{C}_{R Q} \subset \mathcal{C}_{Q}$ which from part c) of Theorem 7.1 means that $F=W_{2}^{\prime} W_{1}^{-1} \in \mathcal{K}_{Q}$. Moreover, assume $\mathcal{C}_{R C} \neq \emptyset$ and take an arbitrary but fixed matrix $W \in \mathcal{C}_{R Q}$. For any $M_{c} \in \mathcal{D}_{c}$, from Theorem 6.9 the set $\mathcal{C}_{R}$ is well defined. Let us first prove that if $W \in \mathcal{C}_{Q R}$ then $W \in \mathcal{C}_{R}$. Indeed, taking into account that $W \geq 0$, then $W \in \mathcal{C}_{Q R}$ is equivalent to (by using Schur complement)

$$
\mathcal{A}_{i}(W)=\left[\begin{array}{cc}
U_{c}^{\prime} \Theta_{c i}(W) U_{c} & U_{c}^{\prime} M_{c i} W^{1 / 2} \\
W^{1 / 2} M_{c i}^{\prime} U_{c} & -r I
\end{array}\right] \leq 0
$$

for all $i=1,2, \cdots, N$ where $U_{c}=\left[\begin{array}{ll}I & 0\end{array}\right]^{\prime}$. Since $\mathcal{A}_{i}(W)$ is affine with respect to $M_{c i}$, writing $M_{c} \in \mathcal{D}_{c}$ as a convex combination of the extreme matrices we get

$$
\mathcal{A}(W):=\sum_{i=1}^{N} \xi_{i} \mathcal{A}_{i}(W) \leq 0
$$

which implies that $W \in \mathcal{C}_{R}$ indeed. From Theorem 6.9 the state feedback gain matrix $F=W_{2}^{\prime} W_{1}^{-1}$ places the closed-loop system poles inside the given circular region $\mathcal{R}$. Furthermore, from part d) of the same theorem we conclude that the linear function trace $\left[R_{c} W\right]$ is a valid upper bound to the $R H_{2}$ norm of the closed-loop transfer function $T(z, w ; s)$. Since $M_{c}$ is an arbitrary matrix in $\mathcal{D}_{c}$, the result follows. From this, the case on which one of the sets is empty trivially holds and so the theorem is proved.

From this theorem, we can now solve the guaranteed cost optimal control problem stated before, namely (7.49)

$$
\begin{align*}
\bar{\rho}_{Q} & =\min \left\{\rho:(F, \rho) \in \mathcal{K}_{R Q}\right\} \\
& =\min \left\{\operatorname{trace}\left[R_{c} W\right]: W \in \mathcal{C}_{R Q}\right\} \tag{7.54}
\end{align*}
$$

which is once again a convex programming problem. It is interesting to observe that, in general lines, the proof of this result follows the same pattern of the previous ones. The main difference is that the matrix function $\Theta_{R}(\cdot)$ depends nonlinearly on matrix $M_{c}$ but thanks to the Schur complement formula it can be converted to an affine and hence convex function from which the desired result is proved.

With this in mind it is possible to handle many other guaranteed cost optimal control problems with no big additional difficulty. An interesting and practically important case is to solve the guaranteed cost version of the Time-domain specification problem. Recalling Theorem 6.10, the optimal solution of the convex programming problem

$$
\begin{equation*}
\bar{\rho}_{Q}=\min \left\{g\left(\bar{R}_{c} W \bar{R}_{c}^{\prime}\right): W \in \mathcal{C}_{Q}\right\} \tag{7.55}
\end{equation*}
$$

provides $F=W_{2}^{\prime} W_{1}^{-1} \in \mathcal{K}_{Q}$ and the upper bound $\rho$ satisfying

$$
\mathcal{G}(T(z, w ; s)) \leq \rho, \quad \forall M_{c} \in \mathcal{D}_{c}
$$

is minimized. Of course this result follows from the fact that the function $g(\cdot)$ defined in Chapter 6 is convex. The same reasoning can be adopted if the designer wants to include convex structural constraints in the state feedback gain matrix like decentralization or static output feedback. For that, it suffices to impose the linear constraint

$$
W=W_{D} \text { or } W=W_{O}
$$

in the problem to be solved. Since these constraints involve only the matrix variable, they do not change any property of the closed-loop system stability and the proposed upper bounds of the transfer function norms meaning that all results obtained before remain valid.

### 7.3 Actuators failure

The practical implementation of a given control is done by means of certain adequate devices or actuators. Of course, in the real world these devices are possible to fail


Figure 7.7: Practical control implementation
during the system operation. In this section we want to analyze the impact of actuators failure in control systems stability and performance. The basic tools to be used are the ones related to robust stability and performance already introduced in the previous section. It will be shown that the polyhedral convex bounded domain $\mathcal{D}_{c}$ plays a central role in this design problem.

Let us consider the linear dynamic system described by the state space equations

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u  \tag{7.56}\\
z & =C_{1} x+D_{12} u  \tag{7.57}\\
y & =x \tag{7.58}
\end{align*}
$$

for which we seek a state feedback control of the form $u=F x$ where $F$ solves the Full information control problem in $\mathrm{RH}_{2}$

$$
\begin{equation*}
\min \left\{\|T(z, w ; s)\|_{2}^{2}: F \in \mathcal{K}_{c}\right\} \tag{7.59}
\end{equation*}
$$

Under the assumptions put in evidence in Chapter 4 this problem is feasible and admits an optimal solution given by

$$
\begin{equation*}
F=F_{2} \tag{7.60}
\end{equation*}
$$

With the optimal gain $F_{2}$ at hand, its practical implementation needs the use of $m$ actuators (recall that $m$ is the dimension of the control vector) as is indicated in fig. 7.7. Each actuator is ideally modeled as a a simple gain $\pi_{i}, i=1,2, \cdots, m$ such that

$$
\begin{equation*}
\pi_{i}:=\{0,1\} \tag{7.61}
\end{equation*}
$$

meaning that it has only two states, namely in operation corresponding to $\pi_{i}=1$ and out of operation corresponding to $\pi_{i}=0$. From this, the control action in the closed-loop system instead of $u=F_{2} x$ is

$$
\begin{equation*}
u(t)=\pi F_{2} x(t), \quad \pi:=\operatorname{diag}\left[\pi_{1}, \pi_{2}, \cdots, \pi_{m}\right] \tag{7.62}
\end{equation*}
$$

Since the optimal gain $F_{2}$ has been calculated using only informations of the nominal open-loop system, during normal operation characterized by $\pi=I$, the closed-loop system evolves following an optimal trajectory. However, if a fail occurs in some control channel, say $i$, the actuator state changes to zero. In this situation the closed-loop system evolves with the control $u=\tilde{F}_{2} x$ where $\tilde{F}_{2}$ is the same as $F_{2}$ but with the $i-t h$ row equal to zero. Of course not only the optimality is lost but even instability may be observed. A way to cope with this problem is to apply the guaranteed control design procedure introduced before. To this end, consider instead
of $F_{2}$ in (7.62) a matrix $F$ to be determined. The closed-loop system matrix can be written as

$$
A_{c c}=A+B_{2} \pi F=A+B_{2}(\pi) F
$$

where $B_{2}(\pi):=B_{2} \pi$. This makes clear that the input matrix $B_{2}$ can be interpreted as an uncertain matrix depending upon the occurrence of actuators failure. Hence, in this framework the dynamic system (7.56) - (7.58) is an uncertain system where

$$
\begin{equation*}
B_{2} \in\left\{B_{2}(\pi): \pi \in \Pi\right\} \tag{7.63}
\end{equation*}
$$

with $\Pi$ being a set of diagonal matrices which defines all combinations of actuators failure. We assume that $I \in \Pi$ and $0 \notin \Pi$. The first condition is clearly necessary because it imposes feasibility to the nominal system. The latter one is also necessary because if it is violated then $B_{2}=0$ is feasible and the system is stabilizable only if it is stable. Suppose the set $\Pi$ is composed by $N$ matrices $\pi_{i}, i=1,2, \cdots, N$ defined before representing all actuators failure we want to take into account. With matrices

$$
M_{c i}:=\left[\begin{array}{cc}
A & B_{2}\left(\pi_{i}\right) \\
0 & 0
\end{array}\right], \quad i=1,2, \cdots, N
$$

the convex domain

$$
\mathcal{D}_{c}:=\operatorname{co}\left\{M_{c i}, i=1,2, \cdots, N\right\}
$$

is the uncertain domain we have to consider since by construction it contains all input matrices $B_{2}(\pi)$ with $\pi \in \Pi$. Clearly, it contains many others, namely the ones generated by convex combination of the pairs $\left[A, B_{2}\left(\pi_{i}\right)\right]$. Fortunately, recalling Remark 7.4, the uncertain system defined by $M_{c} \in \mathcal{D}_{c}$ is quadratic stabilizable if and only if the collection of systems defined by all extreme matrices $M_{c i}$ is quadratic stabilizable. Hence, as far as quadratic stability is concerned, modeling actuators failure as $M_{c} \in \mathcal{D}_{c}$ instead of by means of (7.63) does not introduce any kind of conservativeness in the results.

From the above discussion, let us proceed by replacing problem (7.59) by the associated guaranteed cost control problem. It is given by (recall Theorem 7.2)

$$
\begin{equation*}
\bar{\rho}_{A}:=\min \left\{\operatorname{trace}\left[R_{c} W\right]: W \in \mathcal{C}_{Q}\right\} \tag{7.64}
\end{equation*}
$$

and the optimal solution (if one exists) provides $F_{A}=W_{2}^{\prime} W_{1}^{-1}$ such that with the state feedback $u=F_{A} x$ the closed-loop system is quadratically stable and $\|T(z, w ; s)\|_{2}^{2} \leq$ $\bar{\rho}_{A}$ for all $M_{c} \in \mathcal{D}_{c}$. In other words, we can say that the closed-loop system with the control $u=\pi F_{A} x$ is such that all eigenvalues of $A_{c c}=A+B_{2} \pi F_{A}$ lie in the left open part of the complex plane and

$$
\|T(z, w ; s)\|_{2}^{2} \leq \bar{\rho}_{A}, \quad \forall \pi \in \Pi
$$

that is all control requirements are met. The next example illustrates this design policy for a simple example.

Example 7.5 Let the dynamic system (7.56)-(7.58) be given by

$$
\begin{gathered}
A=\left[\begin{array}{lll}
0 & 1 & 1 \\
1 & 1 & 2 \\
1 & 1 & 1
\end{array}\right], B_{1}=\left[\begin{array}{lll}
1 & 0 & 0 \\
0 & 1 & 0 \\
0 & 0 & 1
\end{array}\right], \quad B_{2}=\left[\begin{array}{ll}
1 & 0 \\
0 & 0 \\
0 & 1
\end{array}\right] \\
C_{1}=\left[\begin{array}{lll}
2 & 2 & 0 \\
1 & 0 & 1
\end{array}\right], \quad D_{12}=\left[\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right]
\end{gathered}
$$

| Failure | Eigenvalues | $\\|T(z, w ; s)\\|_{2}^{2}$ |
| :---: | :---: | :---: |
| $\pi_{1}$ | $-1.91,-1.46+j 0.69,-1.46-j 0.69$ | 6.50 |
| $\pi_{2}$ | $0.43,-0.72,-2.37$ | $\infty$ |
| $\pi_{3}$ | $1.94,-1.99,-0.13$ | $\infty$ |

Table 7.1: Closed-loop system poles and performance

| Failure | Eigenvalues | $\\|T(z, w ; s)\\|_{2}^{2}$ |
| :---: | :---: | :---: |
| $\pi_{1}$ | $-10.56,-0.39+j 0.12,-0.39-j 0.12$ | 32.00 |
| $\pi_{2}$ | $-3.11,-0.63+j 0.18,-0.63-j 0.18$ | 61.49 |
| $\pi_{3}$ | $-2.88,-1.62,-0.47$ | 66.14 |

Table 7.2: Closed-loop system poles and performance

The optimal solution of Problem (7.59) provides

$$
F_{2}=\left[\begin{array}{lll}
-2.17 & -2.67 & -0.79 \\
-1.79 & -3.12 & -4.66
\end{array}\right] \Longrightarrow \min \|T(z, w ; s)\|_{2}^{2} \approx 6.50
$$

For the practical implementation of this state feedback gain we need two actuators. It is supposed that they can fail during operation. Since the open-loop system is unstable, it is assumed that the actuators can not fail simultaneously. Hence the set $\Pi$ is composed by matrices

$$
\pi_{1}=\left[\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right], \quad \pi_{2}=\left[\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right], \quad \pi_{3}=\left[\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right]
$$

corresponding to normal operation ( $\pi_{1}$ ), first control channel actuator failure ( $\pi_{2}$ ) and second control channel actuator failure ( $\pi_{3}$ ) respectively.

Table 7.1 gives the eigenvalues of matrix $A_{c c}=A+B_{2}(\pi) F_{2}$ for the three fault matrices. Obviously, for $\pi=\pi_{1}=I$ the closed-loop system is internally stable. However, if one of the two faults occurs then the closed-loop system becomes unstable. The optimal solution of the Full information problem is not robust for this kind of severe perturbation. On the contrary, the optimal solution of the guaranteed $\mathrm{RH}_{2}$ cost control problem (7.64) gives

$$
F_{A}=\left[\begin{array}{rrr}
-6.97 & -10.17 & -13.56 \\
-3.70 & -5.00 & -6.37
\end{array}\right] \Longrightarrow \bar{\rho}_{A} \approx 73.18
$$

which as indicated in Table 7.2 preserves internal stability in front of any prespecified failure occurrence. In this example, when compared with the nominal case, the price to be paid to keep internal stability under actuators failure is high. A possible interpretation of this fact is that the parametric perturbation caused by actuators failure is actually severe.

Remark 7.13 It is possible also to characterize actuators failure by means of norm bounded uncertainty. Actually, noticing that

$$
B_{2}(\pi)=B_{2} \pi=B_{2}+B_{2}(\pi-I)=B_{2}+B_{2} \Omega
$$

with $\Omega:=\pi-I$ and that $\pi$ is always a diagonal matrix, it is immediate to get

$$
\|\Omega\| \leq \max _{\pi \in \Pi}\|\pi-I\|=1
$$

Hence, with matrices

$$
M_{c n}=\left[\begin{array}{cc}
A & B_{2} \\
0 & 0
\end{array}\right], \quad B_{n}=\left[\begin{array}{c}
B_{2} \\
0
\end{array}\right], \quad C_{n}=\left[\begin{array}{ll}
0 & I
\end{array}\right]
$$

any actuators failure such that $\pi \in \Pi$ may be alternatively described by (recall Remark 7.2)

$$
M_{c} \in \mathcal{D}_{n}=\left\{M_{c}=M_{c n}+B_{n} \Omega C_{n}:\|\Omega\| \leq 1\right\}
$$

This characterization of the problem under consideration frequently leads, if any, to very conservative results and so worse than the use of the polyhedral convex bounded domain $\mathcal{D}_{c}$ as we have done before. Indeed, the particular choice $\Omega=-I$ which obviously is a diagonal matrix such that $\|\Omega\|=1$ produces

$$
\left[\begin{array}{ll}
A & 0 \\
0 & 0
\end{array}\right] \in \mathcal{D}_{n}
$$

which is stabilizable if and only if the open-loop system is stable. In other words, the use of the uncertain domain $\mathcal{D}_{n}$ includes by construction the spurious matrix $\Omega=-I$ being thus equivalent to have $\pi=0 \in \Pi$. This situation is avoided in the description of the set $\mathcal{D}_{c}$ as we have just shown in the solution of the previous example.

### 7.4 Nonlinear perturbations

In this section nonlinear robust stability and performance are considered. This topic is a generalization of the ones introduced before in the sense that advantages are taken into account from the a priori knowledgement of the class of nonlinear perturbations acting in the open-loop model. Two main classes are of importance, namely multiplicative and additive nonlinear perturbations, leading to what we call Persidiskii and Lur'e robust design procedures. In both cases, the control structure is assumed to be linear and the whole state vector is available for feedback.

### 7.4.1 Persidiskii design

Let us first consider the robust control design of a class of nonlinear systems subject to state dependent nonlinear perturbations called multiplicative perturbations. Once again the proposed design procedure will be expressed in terms of convex programming problems only. The block-scheme of the dynamic system to be dealt with is shown in fig. 7.8. The perturbed system $\Sigma_{p}$ is subject to the nonlinear perturbation $f(x)$ to be precisely defined in the sequel. For the moment, it is important to keep clear that the open-loop system is subject to a class of perturbations such that, when they occur, the whole state vector $x$ changes to $f(x)$. Then, the perturbation occurrence changes also the measured output $y$ accordingly.

Assuming the state vector has dimension $n$ and the nonlinear function $f(x)$ is not exactly known, the only available information is that it belongs to the uncertain domain $\mathcal{D}_{f}$ composed by all vector valued functions with the following properties :

1) Each component of $f(x)$ namely $f_{j}(x), j=1,2, \cdots, n$ is a real valued function such that

$$
\begin{equation*}
f_{j}(x)=f_{j}\left(x_{j}\right) \tag{7.65}
\end{equation*}
$$

where $x_{j} \in R$ denotes the $j-t h$ component of the vector $x \in R^{n}$.


Figure 7.8: The bock-scheme of the perturbed system
2) Each component $f_{j}\left(x_{j}\right), j=1,2, \cdots, n$ is such that

$$
\begin{align*}
f_{j}(0) & =0  \tag{7.66}\\
f_{j}(\xi) \xi & >0 \quad \forall \xi \neq 0 \in R  \tag{7.67}\\
\int_{0}^{\infty} f_{j}(\xi) d \xi & =\infty \tag{7.68}
\end{align*}
$$

The first condition defines the perturbation structure. Roughly speaking, the second condition says that the graph of $f(\cdot)$ must be contained in the first and third quadrants of the $(f, \xi)$ plane. Since $f(x)=x \in \mathcal{D}_{f}$ the corresponding linear system is called nominal system ( $\Sigma_{n}$ ) and it has the following state space representation

$$
\begin{align*}
\dot{x} & =A x+B_{2} u, \quad x(0)=x_{0}  \tag{7.69}\\
z & =C_{1} x+D_{12} u  \tag{7.70}\\
y & =x \tag{7.71}
\end{align*}
$$

Adapting the previous design goals to cope with nonlinear systems stability and performance, we proceed trying to determine (if one exists) a linear state feedback control law $u=F x$ for $\Sigma_{n}$ such that the origin $x=0$ of $\Sigma_{p}$ which has the state space representation

$$
\begin{align*}
& \dot{x}=\left(A+B_{2} F\right) f(x), \quad x(0)=x_{0}  \tag{7.72}\\
& z=\left(C_{1}+D_{12} F\right) f(x) \tag{7.73}
\end{align*}
$$

is globally stable for all $f \in \mathcal{D}_{f}$. Furthermore, among all state feedback gains with this property, find the one, namely $F_{f}$, which solves the associated guaranteed cost control problem

$$
\begin{equation*}
\bar{\rho}_{f}\left(x_{0}\right):=\min \bar{\rho}\left(F, x_{0}\right) \tag{7.74}
\end{equation*}
$$

where

$$
\int_{0}^{\infty} z(t)^{\prime} z(t) d t \leq \bar{\rho}\left(F, x_{0}\right), \quad \forall f \in \mathcal{D}_{f}
$$

Similarly the minimum value of $\bar{\rho}\left(F, x_{0}\right)$ with respect to all $F$ preserving stability is called the minimum guaranteed cost associated to the optimal feedback gain $F=$
$F_{f}$. To accomplish the first goal concerning the robust stability of $\Sigma_{n}$ we need to introduce the following important result. To ease notation, for any square matrix $P$, the subscript " d " indicates that $P=P_{d}$ is constrained to be a diagonal matrix.

Theorem 7.6 (Persidiskii theorem) For any given state feedback matrix $F$ suppose there exists a positive definite matrix $P_{d}$ such that

$$
\begin{equation*}
0 \geq\left(A+B_{2} F\right)^{\prime} P_{d}+P_{d}\left(A+B_{2} F\right)+Q \tag{7.75}
\end{equation*}
$$

for some matrix $Q=Q^{\prime}>0$. Then, the origin $x=0$ of the perturbed system $\Sigma_{p}$ is globally asymptotically stable for all $f \in \mathcal{D}_{f}$.

Proof Using the given properties of $f \in \mathcal{D}_{f}$, let us take

$$
v(x):=2 \sum_{j=1}^{n} P_{j j} \int_{0}^{x_{j}} f_{j}(\xi) d \xi
$$

as a Lyapunov function candidate associated to an arbitrary trajectory of $\Sigma_{p}$ such that $x(0)=x_{0}$. Its time derivative can be written as

$$
\begin{aligned}
\dot{v}(x) & =2 \sum_{j=1}^{n} P_{j j} f_{j}\left(x_{j}\right) \dot{x}_{j} \\
& =2 f(x)^{\prime} P_{d} \dot{x} \\
& =f(x)^{\prime}\left[\left(A+B_{2} F\right)^{\prime} P_{d}+P_{d}\left(A+B_{2} F\right)\right] f(x) \\
& \leq-f(x)^{\prime} Q f(x) \\
& <0, \quad \forall x \neq 0
\end{aligned}
$$

proving thus the theorem proposed.
In the above calculations, it is clear that matrix $Q$ must be positive definite but does not need to present any particular structure. Besides, the particular value of this matrix is immaterial to get $P=P_{d}$ satisfying inequality (7.75). In fact, if there exists $P=P_{d}$ satisfying (7.75) for a given $Q=\bar{Q}>0$ then the same is true for any other choice $Q>0$. This new degree of freedom is used in the next lemma to get the upper bound defined in (7.74).

Lemma 7.2 Assume for all $f \in \mathcal{D}_{f}$ there exist $n$ positive and finite parameters such that

$$
\begin{equation*}
\int_{0}^{x_{j}(0)} f_{j}(\xi) d \xi \leq \frac{r_{j}}{2}, \quad j=1,2, \cdots, n \tag{7.76}
\end{equation*}
$$

For any state feedback control gain $F$ such that there exists $P=P_{d}$ satisfying the matrix inequality (7.75), it is possible to choose $Q=Q^{\prime}>0$ such that the upper bound $\bar{\rho}\left(F, x_{0}\right)$ is given by

$$
\begin{equation*}
\bar{\rho}\left(F, x_{0}\right):=\sum_{j=1}^{n} P_{j j} r_{j} \tag{7.77}
\end{equation*}
$$

Proof Assume for $F$ given, there exists $P=P_{d}$ such that the matrix inequality (7.75) holds. In this case, we can choose matrix $Q>0$ as being

$$
Q=\left(C_{1}+D_{12} F\right)^{\prime}\left(C_{1}+D_{12} F\right)+\epsilon I
$$

where $\epsilon>0$ is an arbitrarily small parameter. From the previous theorem we already know that the origin of the perturbed system $\Sigma_{p}$ is globally stable. Hence, using (7.72) and (7.73) we immediately have

$$
\begin{aligned}
\dot{v}(x(t)) & \leq-f(x(t))^{\prime}\left[\left(C_{1}+D_{12} F\right)^{\prime}\left(C_{1}+D_{12} F\right)+\epsilon I\right] f(x(t)) \\
& \leq-z(t)^{\prime} z(t), \quad \forall t \geq 0
\end{aligned}
$$

Integrating both sides of this inequality from $t=0$ to $t=\infty$ we get

$$
\int_{0}^{\infty} z(t)^{\prime} z(t) d t \leq v(x(0))
$$

Using this inequality together with the definition of the Lyapunov equation introduced in Theorem 7.6 and the upper bound on each element of function $f(x)$ yield

$$
\begin{aligned}
\int_{0}^{\infty} z(t)^{\prime} z(t) d t & \leq 2 \sum_{j=1}^{n} P_{j j} \int_{0}^{x_{j}(0)} f_{j}(\xi) d \xi \\
& \leq \sum_{j=1}^{n} P_{j j} r_{j}
\end{aligned}
$$

which being true for all $f \in \mathcal{D}_{f}$ proves that the upper bound (7.77) is valid.
Up to now we have always worked with the dual version of inequality (7.75) where the system matrix transpose post-multiply the matrix variable. This fact was important to convert the associated optimal control problems to convex programming problems. Here, due to the nonlinearity of the perturbation $f(x)$ this is no longer possible. Even though, the guaranteed cost control problem (7.74) can be converted into a convex one by means of Schur complements. To this end, consider the affine matrix function which is defined by all pairs of matrices $(X, Y)$ of appropriate dimension with the first one being symmetric

$$
\mathcal{A}_{f}(X, Y):=\left[\begin{array}{cc}
A X+B_{2} Y+X A^{\prime}+Y^{\prime} B_{2}^{\prime} & X C_{1}^{\prime}+Y^{\prime} D_{12}^{\prime}  \tag{7.78}\\
C_{1} X+D_{12} Y & -I
\end{array}\right]
$$

The following preliminary result is of particular importance towards the complete solution of the guaranteed cost control problem stated before.

Theorem 7.7 Define the convex set

$$
\begin{equation*}
\mathcal{C}_{f}:=\left\{(X, Y): X=X_{d}>0, \quad \mathcal{A}_{f}(X, Y)<0\right\} \tag{7.79}
\end{equation*}
$$

The set of all state feedback matrices $F$ such that (7.75) holds for some $Q>0$, denoted as $\mathcal{K}_{f}$ is alternatively given by

$$
\begin{equation*}
\mathcal{K}_{f}:=\left\{Y X^{-1}:(X, Y) \in \mathcal{C}_{f}\right\} \tag{7.80}
\end{equation*}
$$

Proof It must be clear that the above defined set $\mathcal{C}_{f}$ is convex. Actually, it is defined by LMI's and the linear constraint $X=X_{d}$ which corresponds to impose that all off diagonal elements of matrix $X$ are zero. For the necessity, let us take an arbitrary $F \in \mathcal{K}_{f} \neq \emptyset$ and observe that in this case, there exists $P=P_{d}$ positive definite such that

$$
0>\left(A+B_{2} F\right)^{\prime} P_{d}+P_{d}\left(A+B_{2} F\right)+\left(C_{1}+D_{12} F\right)^{\prime}\left(C_{1}+D_{12} F\right)
$$

Multiplying this inequality to the left and to the right by $P_{d}^{-1}$ and using the Schur complement formula, it is simple to verify that

$$
(X, Y)=\left(P_{d}^{-1}, F P_{d}^{-1}\right) \in \mathcal{C}_{f}
$$

and $Y X^{-1}=F P_{d}^{-1} P_{d}=F$. Conversely, take any $(X, Y) \in \mathcal{C}_{f} \neq \emptyset$. Using again the Schur complement to the LMI $\mathcal{A}_{f}(X, Y)<0$ and taking into account that $X=X_{d}$ is positive definite and diagonal, we conclude that with $P_{d}=X_{d}^{-1}$ and $F=Y X^{-1}$ there exists $\epsilon>0$ sufficiently small such that inequality (7.75) holds for

$$
Q=\left(C_{1}+D_{12} F\right)^{\prime}\left(C_{1}+D_{12} F\right)+\epsilon I
$$

and so the proof of the theorem proposed is complete by noticing that if one of the sets in (7.80) is empty both are empty.

From this result, we are able to generate by means of a feasibility convex problem all gains belonging to the nonconvex set $\mathcal{K}_{f}$. The elements of this set assure robust stability of the nominal closed-loop system against all nonlinear perturbations $f \in \mathcal{D}_{f}$. Besides, using Lemma 7.2 and defining the matrix

$$
D:=\operatorname{diag}\left[\sqrt{r_{1}}, \sqrt{r_{2}}, \cdots, \sqrt{r_{n}}\right]
$$

the elements of the set $\mathcal{C}_{f}$ allow the determination of the upper bound $\bar{\rho}\left(F, x_{0}\right)$ for all $F \in \mathcal{K}_{f}$ as being

$$
\begin{align*}
\bar{\rho}\left(F, x_{0}\right) & =\sum_{j=1}^{n} P_{j j} r_{j} \\
& =\operatorname{trace}\left[D^{\prime} X^{-1} D\right] \tag{7.81}
\end{align*}
$$

valid for all $(X, Y) \in \mathcal{C}_{f}$ and $F=Y X^{-1}$. From this fact the minimum guaranteed cost is readily calculated from

$$
\begin{align*}
\bar{\rho}_{f}\left(x_{0}\right) & =\min _{F \in \mathcal{K}_{f}} \bar{\rho}\left(F, x_{0}\right) \\
& =\inf \left\{\operatorname{trace}\left[D^{\prime} X^{-1} D\right]:(X, Y) \in \mathcal{C}_{f}\right\} \tag{7.82}
\end{align*}
$$

which is a convex programming problem (recall Remark 6.3). Once the global solution of the right hand side of (7.82) is calculated, the corresponding state feedback gain, optimal solution of the left hand side of the same equation is provided simply by $F_{f}=Y X^{-1}$.

Remark 7.14 It follows that $\mathcal{K}_{f} \subset \mathcal{K}_{c}$ and $\mathcal{K}_{f} \neq \emptyset$ if and only if $\mathcal{C}_{f} \neq \emptyset$. The stabilizability of the pair $\left(A, B_{2}\right)$ is not sufficient to guarantee that $\mathcal{C}_{f} \neq \emptyset$ since besides internal stability, the linear constraint $X=X_{d}$ requires the inequality (7.75) to have a diagonal and positive definite feasible solution.

Remark 7.15 The effectiveness to solve problem (7.82) by means of convex programming methods, under the assumption that $D_{12}^{\prime} D_{12}=I$, is now addressed. The main point to be considered is the boundedness of the set $\mathcal{C}_{f}$. Notice that if there exists an unbounded feasible matrix $X>0$ then it is one of the possible global solutions of Problem (7.82) because

$$
0 \leq \operatorname{trace}\left[D^{\prime} X^{-1} D\right]
$$



Figure 7.9: The bock-scheme of the perturbed system
and the right hand side approaches to zero as $X$ increases. Using the same reasoning adopted in Remark 6.8 it is now determined under which conditions this occurrence is avoided. Let us assume that there exists $(X, Y) \in \mathcal{C}_{f}$ and $(\tilde{X}, \tilde{Y}) \neq 0$ such that $(X, Y)+\lambda(\tilde{X}, \tilde{Y}) \in \mathcal{C}_{f}$ for $\lambda>0$ arbitrarily large. Then

$$
\begin{gathered}
\tilde{X} \geq 0 \\
{\left[\begin{array}{cc}
A \tilde{X}+B_{2} \tilde{Y}+\tilde{X} A^{\prime}+\tilde{Y}^{\prime} B_{2}^{\prime} & \tilde{X} C_{1}^{\prime}+\tilde{Y}^{\prime} D_{12}^{\prime} \\
C_{1} \tilde{X}+D_{12} \tilde{Y} & 0
\end{array}\right] \leq 0}
\end{gathered}
$$

From Appendix C, these conditions are verified only if there exists a matrix $\tilde{X}$ such that

$$
\begin{aligned}
C_{1 c} \tilde{X} & =0 \\
\tilde{X} & \geq 0 \\
A_{c} \tilde{X}+\tilde{X} A_{c}^{\prime} & \leq 0
\end{aligned}
$$

where $A_{c}=A-B_{2} D_{12}^{\prime} C_{1}$ and $C_{1 c}=\left(I-D_{12} D_{12}^{\prime}\right) C_{1}$. As in Remark 6.8, this situation is completely avoided if the open-loop system is such that the pair $\left(-A_{c}, C_{1 c}\right)$ is detectable. Under this mild assumption the optimal value of the objective function of Problem (7.82) is strictly positive and finite.

### 7.4.2 Lur'e design

Let us consider now another important robust control design for the class of output dependent nonlinear additive perturbations. The basic block-scheme is given in fig. 7.9. It resembles the one used in the state feedback Mixed $R H_{2} / R H_{\infty}$ control design. The nominal system is denoted by $\Sigma_{n}$. The perturbed dynamic system, namely $\Sigma_{p}$ is subject to the nonlinear perturbation $h(\cdot)$ which is a vector valued function not a priori known. The available information is that it belongs to the uncertain domain $\mathcal{D}_{h}$ composed by all functions presenting the following properties:

1) The vector valued function $h(\cdot)$ is defined for all $\xi \in R^{r}$ and $h(\cdot) \in R^{r}$ where $r$ is a positive integer less or equal the dimension of the state vector $x \in R^{n}$.
2) It is such that

$$
\begin{align*}
h(0) & =0  \tag{7.83}\\
h(\xi)^{\prime} \xi & \leq 0 \quad \forall \xi \in R^{r} \tag{7.84}
\end{align*}
$$

The first condition imposes that, in fig. 7.9, the vectors $w$ and $z_{1}$ have the same dimension. The second one implies that the nonlinear function $-h(\cdot)$ belongs to the sector $[0, \infty)$. In fact, in the one-dimensional case the graph of $-h(\xi)$ in the plane $(h, \xi)$ is in the first and third quadrants. The state space equations of the nominal open-loop system $\Sigma_{n}$, corresponding to $h(\cdot)=0 \in \mathcal{D}_{h}$ is the standard one

$$
\begin{align*}
\dot{x} & =A x+B_{1} w+B_{2} u, \quad x(0)=x_{0}  \tag{7.85}\\
z_{0} & =C_{0} x+D_{0} u  \tag{7.86}\\
z_{1} & =C_{1} x+D_{12} u  \tag{7.87}\\
y & =x \tag{7.88}
\end{align*}
$$

As before, the goal is to design a state feedback control law, namely $u=F x$ such that the closed-loop perturbed system $\Sigma_{p}$ obtained from $\Sigma_{n}$ together with $w=h\left(z_{1}\right)$ presents the following properties associated with its state space representation

$$
\begin{align*}
\dot{x} & =\left(A+B_{2} F\right) x+B_{1} h\left(z_{1}\right), \quad x(0)=x_{0}  \tag{7.89}\\
z_{0} & =\left(C_{0}+D_{0} F\right) x  \tag{7.90}\\
z_{1} & =\left(C_{1}+D_{12} F\right) x \tag{7.91}
\end{align*}
$$

First, the origin $x=0$ must be globally stable for all $h \in \mathcal{D}_{h}$. From all state feedback gains with this property select (if possible) one, namely $F_{h}$ which solves the following guaranteed cost control problem

$$
\begin{equation*}
\bar{\rho}_{h}\left(x_{0}\right):=\min \bar{\rho}\left(F, x_{0}\right) \tag{7.92}
\end{equation*}
$$

where

$$
\int_{0}^{\infty} z_{0}(t)^{\prime} z_{0}(t) d t \leq \bar{\rho}\left(F, x_{0}\right), \quad \forall h \in \mathcal{D}_{h}
$$

The guaranteed cost control problem (7.92) is similar to the Mixed $R H_{2} / R H_{\infty}$ control problem. The existence of the nonlinear function $h \in \mathcal{D}_{h}$ does not allow us to express it in the frequency domain. Instead the guaranteed cost is given in terms of an upper bound to the above integral of the controlled output. Accordingly, the exogenous signal is replaced by an arbitrary initial condition $x(0) \neq 0$.

Theorem 7.8 (Passivity theorem) For any given state feedback matrix $F$ suppose there exists a symmetric and positive definite matrix $P$ such that

$$
\begin{align*}
0 & \geq\left(A+B_{2} F\right)^{\prime} P+P\left(A+B_{2} F\right)+Q  \tag{7.93}\\
B_{1}^{\prime} P & =\left(C_{1}+D_{12} F\right) \tag{7.94}
\end{align*}
$$

for some matrix $Q=Q^{\prime}>0$. Then, the origin $x=0$ of the perturbed system $\Sigma_{p}$ is globally asymptotically stable for all $h \in \mathcal{D}_{h}$.

Proof Consider the Lyapunov function candidate

$$
v(x):=x^{\prime} P x
$$

Its time derivate along any trajectory of $\Sigma_{p}$ is written as

$$
\begin{aligned}
\dot{v}(x) & =2 x^{\prime} P \dot{x} \\
& =x^{\prime}\left[\left(A+B_{2} F\right)^{\prime} P+P\left(A+B_{2} F\right)\right] x+2 x^{\prime} P B_{1} h\left(z_{1}\right)
\end{aligned}
$$

Making use of condition (7.94), for an arbitrary vector $x$ we have

$$
B_{1}^{\prime} P x=\left(C_{1}+D_{12} F\right) x=z_{1}
$$

which together with the first condition (7.93) yields

$$
\begin{aligned}
\dot{v}(x) & \leq-x^{\prime} Q x+z_{1}^{\prime} h\left(z_{1}\right) \\
& <0, \quad \forall x \neq 0
\end{aligned}
$$

where the last inequality holds due to (7.84). Hence, for any $h \in \mathcal{D}_{h}$ the origin $x=0$ of the perturbed system $\Sigma_{p}$ is asymptotically stable.

Remark 7.16 Theorem 7.8 also holds in a more general setting. Indeed the function $h(\cdot)$ can be considered time varying provided $h(\xi, t) \in \mathcal{D}_{h}$ for all $t \geq 0$. Furthermore, the assumption that $Q=Q^{\prime}>0$, yields

$$
\dot{v}(x) \leq-\alpha\|x\|^{2}, \quad \alpha:=\min _{i} \lambda_{i}(Q)>0
$$

That is, the origin is in fact globally exponentially stable.
Remark 7.17 Conditions (7.93) and (7.94) can be rewritten as follows

$$
\begin{aligned}
0 & \geq A_{c c}^{\prime} P+P A_{c c}+Q \\
B_{1}^{\prime} P & =C_{c c}
\end{aligned}
$$

where $A_{c c}:=A+B_{2} F$ and $C_{c c}:=C_{1}+D_{12} F$. Assume that for some $Q=Q^{\prime}>0$, there exist $F$ and $P=P^{\prime}>0$ satisfying them. In this case, summing and subtracting $j \omega P$, with $\omega \in[0, \infty)$, in the right hand side of the first condition, after simple algebraic manipulation we have

$$
0 \leq\left(-j \omega I-A_{c c}^{\prime}\right)^{-1} P+P\left(j \omega I-A_{c c}\right)^{-1}-\left(-j \omega I-A^{\prime}\right)^{-1} Q(j \omega I-A)^{-1}
$$

Finally, multiplying this inequality to the left by $B_{1}^{\prime}$ and to the right by $B_{1}$, making use of the second condition above and defining the closed-loop system transfer function

$$
G(s):=C_{c c}\left(s I-A_{c c}\right)^{-1} B_{1}
$$

we get

$$
\begin{aligned}
G^{\sim}(j \omega)+G(j \omega) & \geq B_{1}^{\prime}\left(-j \omega-A^{\prime}\right)^{-1} Q(j \omega-A)^{-1} B_{1} \\
& >0, \quad \forall \omega \in[0, \infty)
\end{aligned}
$$

A transfer function with this property is called strictly positive real ( $S P R$ ) and plays a central role in the stability analysis of nonlinear systems. In this remark we have shown that conditions (7.93) and (7.94) implies $G(s)$ is strictly positive real. Under certain mild additional assumptions the converse is also true and constitutes the well known KalmanYacubovitch lemma. Notice that the above frequency domain characterization of a $S P R$ transfer functions does not include the limit point $\omega=\infty$ for which it is not satisfied because $G(j \infty)=0$.

Remark 7.18 A transfer function $G_{E}(s)$ is called extended strictly positive real (ESPR) if it is $S P R$ and

$$
G_{E}^{\sim}(j \infty)+G_{E}(j \infty)>0
$$

From this definition, it is obvious that the transfer function $G(s)$ defined in the previous remark is $S P R$ but not $E S P R$ due to the fact that $G(j \infty)=0$. It is interesting to known
that there is a very simple relationship between $E S P R$ and $R H_{\infty}$ norm. Indeed, a transfer function $G_{E}(s)$ is $E S P R$ if and only if

$$
\left\|\left[\alpha I-G_{E}(s)\right]\left[\alpha I+G_{E}(s)\right]^{-1}\right\|_{\infty}<1
$$

where $\alpha$ is an arbitrary positive scalar. A proof of this property is not included here. We only mention that it is based on the fact that the above inequality holds if and only if

$$
\left[\alpha I+G_{E}^{\sim}(j \omega)\right]^{-1}\left[\alpha I-G_{E}^{\sim}(j \omega)\right]\left[\alpha I-G_{E}(j \omega)\right]\left[\alpha I+G_{E}(j \omega)\right]^{-1}<1, \quad \forall \omega
$$

or, equivalently

$$
\left[\alpha I-G_{E}^{\sim}(j \omega)\right]\left[\alpha I-G_{E}(j \omega)\right]<\left[\alpha I+G_{E}^{\sim}(j \omega)\right]\left[\alpha I+G_{E}(j \omega)\right], \quad \forall \omega
$$

After obvious simplifications, it provides

$$
2 \alpha\left[G_{E}^{\sim}(j \omega)+G_{E}(j \omega)\right]>0, \quad \forall \omega
$$

implying that $G_{E}(s)$ is $E S P R$. Indeed, for transfer functions which are only $S P R$, as for instance $G(s)$, it is violated at $\omega=\infty$. In this case, the consequence is

$$
\begin{aligned}
\left\|[\alpha I-G(s)][\alpha I+G(s)]^{-1}\right\|_{\infty} & =\sup _{\omega} \bar{\sigma}\left[[\alpha I-G(j \omega)][\alpha I+G(j \omega)]^{-1}\right] \\
& \geq \bar{\sigma}\left[[\alpha I-G(j \infty)][\alpha I+G(j \infty)]^{-1}\right] \\
& \geq 1
\end{aligned}
$$

which means that the previous constraint expressed in terms of a $R H_{\infty}$ norm is also violated.
It is interesting to observe what happen if we do not take care of the above results and try to convert our control design problem to a state feedback problem in $R H_{\infty}$ associated to the auxiliary transfer function

$$
H(s):=[\alpha I-G(s)][\alpha I+G(s)]^{-1}
$$

where $G(s)$ is the transfer function defined in the previous remark. With no loss of generality we adopt $\alpha=1$. Routine calculations (recall in Chapter 2 the definition of inverse system) show that the following is true

$$
\begin{gathered}
{[I+G(s)]^{-1}=\left[\begin{array}{c|c}
A_{c c}-B_{1} C_{c c} & B_{1} \\
\hline-C_{c c} & I
\end{array}\right]} \\
G(s)[I+G(s)]^{-1}=\left[\begin{array}{c|c}
A_{c c}-B_{1} C_{c c} & B_{1} \\
\hline C_{c c} & 0
\end{array}\right]
\end{gathered}
$$

from which it follows that

$$
H(s)=\left[\begin{array}{c|c}
A_{c c}-B_{1} C_{c c} & B_{1} \\
\hline-2 C_{c c} & I
\end{array}\right]
$$

Then, using the state space representation of the transfer function $H(s)$, the determination of a state feedback gain such that $\|H(s)\|_{\infty}<1$ is reduced to the solution of the state feedback problem in $R H_{\infty}$ for the auxiliary plant

$$
\begin{aligned}
\dot{x} & =\tilde{A} x+\tilde{B}_{1} w+\tilde{B}_{2} u \\
z & =\tilde{C}_{1} x+\tilde{D}_{11} w+\tilde{D}_{12} u \\
u & =F x
\end{aligned}
$$

where $\tilde{A}:=A-B_{1} C_{1}, \tilde{B}_{1}:=B_{1}, \tilde{B}_{2}:=B_{2}-B_{1} D_{12}, \tilde{C}_{1}:=-2 C_{1}, \tilde{D}_{11}:=I$ and $\tilde{D}_{12}:=$ $-2 D_{12}$. However, this is impossible because with $\bar{\sigma}\left[\tilde{D}_{11}\right]=1$, any state feedback gain such that $\tilde{A}+\tilde{B}_{2} F$ is stable implies $\|T(z, w ; s)\|_{\infty} \geq 1$.

In opposition to Theorem 7.6, the constraints (7.93) and (7.94) to be simultaneously satisfied, depend strongly on the particular choice of matrix $Q>0$. Even though, to be able to express the upper bound $\bar{\rho}\left(F, x_{0}\right)$ conveniently we need to impose

$$
\begin{equation*}
Q=\left(C_{0}+D_{0} F\right)^{\prime}\left(C_{0}+D_{0} F\right)+\epsilon I \tag{7.95}
\end{equation*}
$$

with $\epsilon>0$ being an arbitrarily small parameter. Indeed, with this particular choice, following the proof of Theorem 7.8 we have

$$
\begin{aligned}
\dot{v}(x(t)) & \leq-x(t)^{\prime} Q x(t) \\
& \leq-z_{0}(t)^{\prime} z_{0}(t), \quad \forall t \geq 0
\end{aligned}
$$

which after integration from $t=0$ to $t=\infty$ provides

$$
\int_{0}^{\infty} z_{0}(t)^{\prime} z_{0}(t) d t \leq v(x(0))=x_{0}^{\prime} P x_{0}
$$

Based on this, it is natural to define

$$
\begin{equation*}
\bar{\rho}\left(F, x_{0}\right):=x_{0}^{\prime} P x_{0} \tag{7.96}
\end{equation*}
$$

as a valid upper bound for all $h \in \mathcal{D}_{h}$. Furthermore, let us denote as $\mathcal{K}_{h}$ the set of all state feedback gains $F$ such that with $Q>0$ given in (7.95) both constraints (7.93) and (7.94) are simultaneously satisfied for some $P>0$ and introduce the affine matrix functions defined for all pairs of matrices $(X, Y)$ of appropriate dimension with the first one being symmetric

$$
\mathcal{A}_{h}(X, Y):=\left[\begin{array}{cc}
A X+B_{2} Y+X A^{\prime}+Y^{\prime} B_{2}^{\prime} & X C_{0}^{\prime}+Y^{\prime} D_{0}^{\prime} \\
C_{0} X+D_{0} Y & -I
\end{array}\right]
$$

and

$$
\mathcal{B}_{h}(X, Y):=C_{1} X+D_{12} Y-B_{1}^{\prime}
$$

The following theorem gives a complete parametrization of the set $\mathcal{K}_{h}$ in terms of a convex set. It is the basis for the solution of the associated optimal guaranteed cost control problem (7.92).

Theorem 7.9 Define the convex set

$$
\begin{equation*}
\mathcal{C}_{h}:=\left\{(X, Y): X>0, \quad \mathcal{A}_{h}(X, Y)<0, \mathcal{B}_{h}(X, Y)=0\right\} \tag{7.97}
\end{equation*}
$$

The set $\mathcal{K}_{h}$ is alternatively given by

$$
\begin{equation*}
\mathcal{K}_{h}:=\left\{Y X^{-1}:(X, Y) \in \mathcal{C}_{h}\right\} \tag{7.98}
\end{equation*}
$$

Proof Since the set $\mathcal{C}_{h}$ is defined by means of affine functions only, it is convex. To prove the necessity, assuming that $F \in \mathcal{K}_{h} \neq \emptyset$ then there exists $P>0$ such that

$$
0>\left(A+B_{2} F\right)^{\prime} P+P\left(A+B_{2} F\right)+\left(C_{0}+D_{0} F\right)^{\prime}\left(C_{1}+D_{0} F\right)
$$

Multiplying both sides of this inequality by $P^{-1}$ and using the Schur complement it is readily verified that

$$
\mathcal{A}_{h}\left(P^{-1}, F P^{-1}\right)<0
$$

On the other hand, multiplying (7.94) to the right by $P^{-1}$, it can be rewritten as

$$
\mathcal{B}_{h}\left(P^{-1}, F P^{-1}\right)=0
$$

hence, $(X, Y)=\left(P^{-1}, F P^{-1}\right) \in \mathcal{C}_{h}$ and $Y X^{-1}=F$ from which the necessity follows. The sufficiency for $\mathcal{C}_{h} \neq \emptyset$ is immediate. From this, it is also immediate to see that (7.98) holds in case one of the indicated sets is empty.

We have now all elements to face the optimal guaranteed cost control problem (7.92). From Theorem 7.9 and (7.96) it reduces to the problem

$$
\begin{align*}
\bar{\rho}_{h}\left(x_{0}\right) & =\min _{F \in \mathcal{K}_{h}} \bar{\rho}\left(F, x_{0}\right) \\
& =\inf \left\{x_{0}^{\prime} X^{-1} x_{0}:(X, Y) \in \mathcal{C}_{h}\right\} \tag{7.99}
\end{align*}
$$

which is a convex programming problem (recall Remark 6.3). Its global optimal solution provides both the minimum guaranteed cost $\bar{\rho}_{h}\left(x_{0}\right)$ and the associated optimal state feedback gain $F_{h}=Y X^{-1}$.

Remark 7.19 The inclusion $\mathcal{K}_{h} \subset \mathcal{K}_{c}$ is obvious. It is also clear that the stabilizability of the pair $\left(A, B_{2}\right)$ is necessary but not sufficient to $\mathcal{C}_{h} \neq \emptyset$. The necessary and sufficient condition is the existence of a state feedback gain $F \in \mathcal{K}_{c}$ such that the closed-loop system transfer function $G(s)$ is $S P R$ (recall Remark 7.17).

Remark 7.20 Following the same reasoning of Remark 7.15, the boundedness of the set $\mathcal{C}_{h}$ is assured whenever the pair $\left[-\left(A-B_{2} D_{0}^{\prime} C_{0}\right),\left(I-D_{0} D_{0}^{\prime}\right) C_{0}\right]$ is detectable.

Remark 7.21 Consider the convex programming problem

$$
\inf \left\{x_{0}^{\prime} X^{-1} x_{0}: X>0, \mathcal{A}_{h}(X, Y)<0\right\}
$$

which is the same as problem (7.99) but without the linear constraint $\mathcal{B}_{h}(\cdot)=0$. Let us search a solution of the above problem such that

$$
Y=-B_{2}^{\prime}-D_{0}^{\prime} C_{0} X
$$

Using this formula, taking into account the standard assumption $D_{0}^{\prime} D_{0}=I$ the Schur complement implies that $\mathcal{A}_{h}(X, Y)<0$ if and only if

$$
\mathcal{A}(X):=\left[\begin{array}{cc}
A_{c} X+X A_{c}^{\prime}-B_{2} B_{2}^{\prime} & X C_{0 c}^{\prime} \\
C_{0 c} X & -I
\end{array}\right]<0
$$

where $A_{c}:=A-B_{2} D_{0}^{\prime} C_{0}$ and $C_{0 c}:=\left(I-D_{0} D_{0}^{\prime}\right) C_{0}$. Hence, the problem under consideration can be solved in two steps. The convex programming problem

$$
\inf \left\{x_{0}^{\prime} X^{-1} x_{0}: X>0, \mathcal{A}(X)<0\right\}
$$

provides the optimal matrix $X>0$ and the previous formula gives matrix $Y$. The corresponding state feedback gain turns out to be $F=Y X^{-1}$. Comparing the last problem with respect to the determination of matrix $X$ and Remark 6.9 , it is seen that its optimal solution is exactly the associated Riccati equation and so, in this particular case, the guaranteed cost reduces exactly to the minimum $\mathrm{RH}_{2}$ optimal cost. This gives a measure of the quality of the upper bound $\bar{\rho}\left(F, x_{0}\right)$ proposed. In the general case, this solution is no longer feasible and the minimization must be done taking into account the linear constraint $\mathcal{B}_{h}(\cdot)=0$.

### 7.5 Notes and references

Uncertain systems design is nowadays a very wide topic in control theory. Concerning stability by means of linear state feedback control, the paper of Barmish [4] is important since, for the first time the author propose an effective and simple way to handle uncertainties acting on both $A$ and $B_{2}$ system matrices. The notion of guaranteed cost has been introduced by Chang and Peng in [11] related to a simple LQ problem. Differently of what we have done in this chapter, the main idea was to get an (nonlinear) upper bound to the associated Riccati equation. In section 7.2 the problems already solved in Chapter 6 have been revisited. They have been solved in a manner to cope with parameter uncertainty. Two of the most important classes of uncertainty have been considered and compared, namely polyhedral convex bounded and norm bounded uncertainty. For the first type, results from [65] have been used in Remark 7.15 while the numerical example fully described in Example 7.65 and analyzed throughout the section is also given in [50]. With this last paper the reader can go deeper on the comparison of these two types of parameter uncertainty models just mentioned. The main part of Section 7.2 related to $R H_{2}$ and $R H_{\infty}$ guaranteed cost control problems is based on [24] and [49] while the other results are the natural generalizations, to cope with parameter uncertainties, of problems introduced in Chapter 6. The stability and guaranteed cost control of dynamic linear systems subject to actuators failure has been analyzed in [23]. In Section 7.3 this problem is again solved but special attention is paid to the comparison and modeling this special kind of uncertainty by means of the domain $\mathcal{D}_{n}$. Once again, the convexity plays a central role and it is possible to verify that the uncertainty description by means of the convex domain $\mathcal{D}_{c}$ leads in many instances to better results. Section 7.4 is entirely devoted to control design problems involving nonlinear perturbations. The first one called Persidiskii design is based on papers [30], [25] and [26]. The former paper also provides many others and more general results and is an excellent reference on this topic. Finally the second control design procedure called Lur'e design is based on the classical results reported in the important book [61] where the notions of passivity and strictly positive real transfer functions are addressed in a general and complete setting.

## Appendix A

## Some Facts on Polynomials

The three results presented in this appendix can be found in many text of Algebra.
Lemma A. 1 Let $r_{0}(s)$ and $r_{1}(s)$ be two polynomials with $\operatorname{deg}\left[r_{0}(s)\right] \geq \operatorname{deg}\left[r_{1}(s)\right]$. Let $r(s)$ be a greatest common divisor of $r_{0}(s)$ and $r_{1}(s)$. Then, there exist two polynomials $\varphi_{0}(s)$ and $\varphi_{1}(s)$ such that

$$
\varphi_{0}(s) r_{0}(s)+\varphi_{1}(s) r_{1}(s)=r(s)
$$

Proof The sequence of polynomials

$$
\begin{array}{rlrl}
r_{0}(s) & =r_{1}(s) q_{1}(s)+r_{2}(s), & \operatorname{deg}\left[r_{2}(s)\right]<\operatorname{deg}\left[r_{1}(s)\right] \\
r_{1}(s) & =r_{2}(s) q_{2}(s)+r_{3}(s), & \operatorname{deg}\left[r_{3}(s)\right]<\operatorname{deg}\left[r_{2}(s)\right] \\
r_{2}(s) & =r_{3}(s) q_{3}(s)+r_{4}(s), & \operatorname{deg}\left[r_{4}(s)\right]<\operatorname{deg}\left[r_{3}(s)\right] \\
\vdots & =\quad \vdots \\
r_{p-1}(s) & =r_{p}(s) q_{p}(s)+r_{p+1}(s), \quad \operatorname{deg}\left[r_{p+1}(s)\right]<\operatorname{deg}\left[r_{p}(s)\right] \\
r_{p}(s) & =r_{p+1}(s) q_{p+1}(s)
\end{array}
$$

is well defined for a certain $p \leq \operatorname{deg}\left[r_{1}(s)\right]-1$.
From such a sequence, it is easily checked that the polynomial $r_{p+1}(s)$ divides $r_{p}(s)$ and then also $r_{p-1}(s)$ and so on and so forth till $r_{1}(s)$ and $r_{0}(s)$.

The single elements of the sequence can be written as

$$
\begin{aligned}
r_{2}(s) & =r_{0}(s)-r_{1}(s) q_{1}(s) \\
r_{3}(s) & =r_{1}(s)-r_{2}(s) q_{2}(s) \\
r_{4}(s) & =r_{2}(s)-r_{3}(s) q_{3}(s) \\
\vdots & =\quad \vdots
\end{aligned}
$$

Hence, $r_{i}(s)=r_{1}(s) \varphi_{p+2-i}(s)+\varphi_{p+1-i}(s) r_{0}(s), i>1$, where $\varphi_{p+1-i}(s)$ and $\varphi_{p-i+2}(s)$ are suitable polynomials, so that $r_{p+1}(s)=r_{0}(s) \varphi_{0}(s)+r_{1}(s) \varphi_{1}(s)$. Hence, any common factor of the two polynomials $r_{0}(s)$ and $r_{1}(s)$ is also a factor of $r_{p+1}(s)$ so that this last polynomial is a greatest common divisor of $r_{0}(s)$ and $r_{1}(s)$. Through a suitable choice of $q_{p+1}(s)$ such a polynomial can be made coincident with $r(s)$ (recall that two greatest common divisors of two assigned polynomials differ for a multiplicative factor).

Theorem A. 1 Two polynomials $r_{0}(s)$ and $r_{1}(s)$ are coprime if and only if there exist two polynomials $\varphi_{0}(s)$ and $\varphi_{1}(s)$ such that

$$
\begin{equation*}
r_{0}(s) \varphi_{0}(s)+r_{1}(s) \varphi_{1}(s)=1 \tag{A.1}
\end{equation*}
$$

Proof If the two polynomials are coprime, i.e. 1 is a greatest common divisor, then eq. (A.1) follows from Lemma A.1. Conversely, if eq. (A.1) holds, then $r_{0}(s)$ and $r_{1}(s)$ are obviously coprime, because any common factor should appear at the right hand side of this equation.

Lemma A. 2 Let $p(s), q(s)$ and $r(s)$ be three polynomials with $p(s)$ and $q(s)$ coprime. Then there exist two polynomials $\varphi(s)$ and $\psi(s)$ with $\operatorname{deg}[\varphi(s)]<\operatorname{deg}[q(s)]$ such that

$$
\varphi(s) p(s)+\psi(s) q(s)=r(s)
$$

Proof Being $p(s)$ and $q(s)$ coprime, there exist, in view of Theorem A.1, two polynomials $\bar{\varphi}(s)$ and $\bar{\psi}(s)$ such that $\bar{\varphi}(s) p(s)+\bar{\psi}(s) q(s)=1$, so that, letting $\hat{\varphi}(s):=$ $\bar{\varphi}(s) r(s), \hat{\psi}(s):=\bar{\psi}(s) r(s)$, it follows

$$
\hat{\varphi}(s) p(s)+\hat{\psi}(s) q(s)=r(s)
$$

If $\operatorname{deg}[\hat{\varphi}(s)]<\operatorname{deg}[q(s)]$, the two polynomials $\hat{\varphi}(s)$ and $\hat{\psi}(s)$ satisfy the conclusion of the theorem, otherwise for each polynomial $\vartheta(s)$ it results

$$
\begin{equation*}
p(s)[\hat{\varphi}(s)-\vartheta(s) q(s)]+q(s)[\hat{\psi}(s)+\vartheta(s) p(s)]=r(s) \tag{A.2}
\end{equation*}
$$

If one performs the division of $\hat{\varphi}(s)$ by $q(s)$ it follows

$$
\hat{\varphi}(s)=q(s) \beta(s)+\alpha(s)
$$

with $\operatorname{deg}[\alpha(s)]<\operatorname{deg}[q(s)]$. It is immediate to ascertain that the polynomials $\varphi(s):=$ $\hat{\varphi}(s)-\beta(s) q(s)$ and $\psi(s)):=\hat{\psi}(s)+\beta(s) p(s)(\operatorname{set} \vartheta(s)=\beta(s)$ in eq. (A.2)) verify the claim.

## Appendix B

## Singular Values of Matrices

In this section the proofs of the results presented in Section 2.6 are reported, along with a few useful matrix properties (Lemmas B.8-B.15) and hints on matrix manipulations. The book by Lawson and Hanson [38] is the main reference for the singular value decomposition. For the proof of Theorem 2.8 some preliminary results are needed (Lemmas B.1-B.7).

Lemma B. 1 Let $v:=\left[\begin{array}{llll}v_{1} & v_{2} & \cdots & v_{n}\end{array}\right]^{\prime} \neq 0$. Then there exists a unitary matrix $Q$ (Householder matrix transformation) such that

$$
Q v=-\sigma\|v\| e_{1}
$$

where $e_{1}:=\left[\begin{array}{llll}1 & 0 & \cdots & 0\end{array}\right]^{\prime} \in R^{n}, \quad \sigma:=e^{j \arg \left(v_{1}\right)}$
Proof Take $u:=v+\sigma\|v\| e_{1}$ and notice that $u:=\left[\begin{array}{llll}u_{1} & u_{2} & \cdots & u_{n}\end{array}\right]^{\prime} \neq 0$ since $u_{1}=$ $\sigma\left(\left|v_{1}\right|+\|v\|\right)$ and $v \neq 0$. Moreover, let

$$
Q:=I-2 \frac{u u^{\sim}}{u^{\sim} u}
$$

Then,

$$
Q^{\sim} Q=\left(I-2 \frac{u u^{\sim}}{u^{\sim} u}\right)\left(I-2 \frac{u u^{\sim}}{u^{\sim} u}\right)=I-4 \frac{u u^{\sim}}{u^{\sim} u}+4 \frac{u u^{\sim} u u^{\sim}}{u^{\sim} u u^{\sim} u}=I
$$

since $u u^{\sim} u u^{\sim}=u^{\sim} u u u^{\sim}$ (actually $u^{\sim} u$ is scalar). Hence, it has been shown that $Q$ is an unitary matrix. On the other hand,

$$
\begin{aligned}
u^{\sim} u & =\left(v^{\sim}+\sigma^{\sim}\|v\| e_{1}^{\prime}\right)\left(v+\sigma\|v\| e_{1}\right) \\
& =\|v\|^{2}+|\sigma|^{2}\|v\|^{2}+\sigma^{\sim}\|v\| v_{1}+v_{1}^{\sim} \sigma\|v\| \\
& =2\|v\|^{2}+2 \mid v_{1}\| \| v \|
\end{aligned}
$$

since $|\sigma|=1$ and $\sigma^{\sim} v_{1}=\sigma v_{1}^{\sim}=\left|v_{1}\right|$. Therefore,

$$
\begin{aligned}
Q v & =Q\left(u-\sigma\|v\| e_{1}\right) \\
& =\left(I-\frac{2 u u^{\sim}}{u^{\sim} u}\right)\left(u-\sigma\|v\| e_{1}\right) \\
& =u-\sigma\|v\| e_{1}-2 u+\frac{2 u}{u^{\sim} u} u_{1}^{\sim} \sigma\|v\| \\
& =-u-\sigma\|v\| e_{1}+\frac{2 u\left(v_{1}^{\sim}+\sigma^{\sim}\|v\|\right) \sigma\|v\|}{2\left(\|v\|^{2}+\mid v_{1}\| \| v \|\right)} \\
& =-\sigma\|v\| e_{1}
\end{aligned}
$$

The last equality follows from $v_{1}^{\sim}+\sigma^{\sim}\|v\|=\sigma^{\sim}\left(\left|v_{1}\right|+\|v\|\right)$ and $\sigma^{\sim} \sigma=1$.
Lemma B. 2 Let $v:=\left[\begin{array}{ll}v_{1} & v_{2}\end{array}\right]^{\prime} \neq 0$ be an element of $C^{2}$ and let $\vartheta_{i}:=\arg \left(v_{i}\right)$. Then there exists a unitary matrix $G$ (Givens transformation matrix) such that

$$
G v=\left[\begin{array}{c}
\|v\| e^{j \vartheta_{1}} \\
0
\end{array}\right]
$$

Proof Take

$$
G:=\left[\begin{array}{cc}
\left|v_{1}\right| & \left|v_{2}\right| e^{j\left(\vartheta_{1}-\vartheta_{2}\right)} \\
\left|v_{2}\right| e^{j\left(\vartheta_{2}-\vartheta_{1}\right)} & -\left|v_{1}\right|
\end{array}\right]\|v\|^{-1}
$$

It is readily seen that $G^{\sim} G=I$ and $G v$ has the form claimed in the statement.
Lemma B. 3 Let $B$ be a $n \times m$ matrix. Then there exists a unitary matrix $Q$ such that $R:=Q B$ is upper triangular.

Proof In view of Lemma B. 1 there exists a unitary matrix $Q_{1}$ such that

$$
Q_{1} B=\left[\begin{array}{cccc}
x & x & \cdots & x \\
0 & & & \\
\vdots & & B_{1} & \\
0 & & &
\end{array}\right]
$$

where the $x$ 's are generic scalar numbers. Again in view of Lemma B. 1 there exists a unitary matrix $P_{2}$ such that $P_{2} B_{1}$ has the same structure as $Q_{1} B$. Then, defining the unitary matrix

$$
Q_{2}:=\left[\begin{array}{cc}
I & 0 \\
0 & P_{2}
\end{array}\right]
$$

it follows

$$
Q_{2} Q_{1} B=\left[\begin{array}{ccccc}
x & x & x & \cdots & x \\
0 & x & x & \cdots & x \\
0 & 0 & & & \\
\vdots & \vdots & & B_{2} & \\
0 & 0 & & &
\end{array}\right]
$$

Iterating this procedure at most $n-1$ times, one can conclude that matrix $Q$, given by $Q:=Q_{n-1} Q_{n-2} \cdots Q_{2} Q_{1}$, is such that $Q B$ is upper triangular.

Lemma B. 4 Let $B$ be a $n \times m$ matrix with $\operatorname{rank}[B]=k$. Then, there exist two unitary matrices $Q$ and $P$ such that

$$
\left.Q B P=\left[\begin{array}{cc}
R & T \\
0 & 0
\end{array}\right]\right\} k \text { rows }
$$

Proof Let $P$ be the permutation matrix (which is obviously unitary) such that the first $k$ columns of $B P$ are linearly independent. In view of Lemma B. 3 there exists a unitary matrix $Q$ such that $S:=Q B P$ is upper triangular. Obviously, the first $k$ columns of $S$ are linearly independent. The $i$-th row of $S$ is zero for $i>k$, since, otherwise, the triangularity of $S$ would imply that $\operatorname{rank}[S]=\operatorname{rank}[B]>k$, which is a contradiction.

Lemma B.5 Let $\left[\begin{array}{ll}R & T\end{array}\right]$ be a $k \times m$ matrix with $\operatorname{rank}[R]=k$. Then there exists a unitary matrix $W$ such that $\left[\begin{array}{ll}R & T\end{array}\right] W=\left[\begin{array}{ll}\hat{R} & 0\end{array}\right]$, with $\hat{R}$ nonsingular and lower triangular.

Proof Lemma B. 4 implies that there exist two unitary matrices $Q$ and $P$ such that

$$
Q\left[\begin{array}{l}
R^{\sim} \\
T^{\sim}
\end{array}\right] P=\left[\begin{array}{c}
\hat{R}^{\sim} \\
0
\end{array}\right]
$$

with $\hat{R}^{\sim}$ upper triangular. Thanks to the assumption on the rank of $R$, the permutation matrix $P$ is actually the identity matrix and $\hat{R}^{\sim}$ is nonsingular. The thesis now immediately follows by letting $W=Q^{\sim}$.

Lemma B. 6 Let $B$ be a $n \times m$ matrix with $\operatorname{rank}[B]=k$. Then there exist two unitary matrices $H$ and $K$ such that $H^{\sim} B K=R$ with

$$
R=\left[\begin{array}{cc}
R_{11} & 0 \\
0 & 0
\end{array}\right]
$$

where $R_{11}$ is triangular and nonsingular.
Proof The proof is straightforward in view of Lemmas B. 4 and B.5.
Lemma B. 7 Let $A$ be a square nonsingular matrix. Then there exist two unitary matrices $U$ and $V$ such that $U^{\sim} A V=S$, where $S$ is diagonal with real positive entries.

Proof Let $V$ be a unitary matrix such that $A^{\sim} A=V D V^{\sim}$ with $D$ real diagonal and define $S$ in such a way that its generic $(i, j)$ element is the square root of the $(i, j)$ element of $D$, so that $S^{\sim} S=S^{2}=D$. Let now $U:=A V S^{-1}$ (recall that, being $A^{\sim} A>0$, matrices $D$ and $S$ are nonsingular). It follows that $U^{\sim} U=S^{-1} V^{\sim} A^{\sim} A V S^{-1}=S^{-1} V^{\sim} V D V^{\sim} V S^{-1}=S^{-1} D S^{-1}=I$, so that $U$ is unitary. Finally, $U S V^{\sim}=A V S^{-1} S V^{\sim}=A V V^{\sim}=A$. Notice that, by suitably choosing matrix $V$, it is possible to arrange the elements on the diagonal of $D$ so that they are nonincreasing.

Proof of Theorem 2.8 Thanks to Lemma B.6, there exist two unitary matrices $H$ and $K$ such that $A=H R K^{\sim}$ with

$$
R=\left[\begin{array}{cc}
R_{11} & 0 \\
0 & 0
\end{array}\right]
$$

where $R_{11}$ is $k$-dimensional and nonsingular. In view of Lemma B. 7 there exist two unitary matrices $\hat{U}$ and $\hat{V}$ such that $R_{11}=\hat{U} \hat{S} \hat{V}^{\sim}$ where matrix $\hat{S}$ is diagonal with positive nonincreasing elements. Letting

$$
\begin{aligned}
& \bar{U}:= {\left.\left.\left[\begin{array}{ll}
\hat{U} & 0 \\
0 & I
\end{array}\right]\right\} n-k \text { rows } \quad, \quad \bar{V}:=\left[\begin{array}{cc}
\hat{V} & 0 \\
0 & I
\end{array}\right]\right\} m-k \text { rows } } \\
& S:= {\left.\left[\begin{array}{ll}
\hat{S} & 0 \\
0 & 0
\end{array}\right]\right\} n-k \text { rows } } \\
& \underbrace{}_{m-k \text { columns }}
\end{aligned}
$$

it follows

$$
\bar{U} S \bar{V}^{\sim}=\left[\begin{array}{cc}
R_{11} & 0 \\
0 & 0
\end{array}\right]=R
$$

so that $A=H R K^{\sim}=H \bar{U} S \bar{V}^{\sim} K^{\sim}=U S V^{\sim}$, with $U:=H \bar{U}$ and $V:=K \bar{V}$.
Proof of Lemma 2.16 Point 1) Let $T$ be a unitary matrix such that

$$
T A^{\sim} A T^{\sim}=D=\operatorname{diag}\left\{\lambda_{i}^{2}\right\}
$$

Hence

$$
\begin{aligned}
\max _{x \neq 0} \frac{\|A x\|^{2}}{\|x\|^{2}} & =\max _{x \neq 0} \frac{x^{\sim} A^{\sim} A x}{x^{\sim} T^{\sim} T x} \\
& =\max _{x \neq 0} \frac{x^{\sim} T^{\sim} T A^{\sim} A T^{\sim} T x}{x^{\sim} T^{\sim} T x} \\
& =\max _{z=T x \neq 0} \frac{z^{\sim} D z}{z^{\sim} z} \\
& =\max _{z \neq 0} \frac{\sum_{i} z_{i}^{\sim} z_{i} \lambda_{i}^{2}}{\sum_{i} z_{i}^{\sim} z_{i}} \\
& \leq \max _{z \neq 0} \frac{\bar{\sigma}^{2}(A) \sum_{i} z_{i}^{\sim} z_{i}}{\sum_{i} z_{i}^{\sim} z_{i}}=\bar{\sigma}^{2}(A)
\end{aligned}
$$

On the other hand, if $A^{\sim} A \xi=\bar{\sigma}^{2}(A) \xi, \xi \neq 0$, it results

$$
\bar{\sigma}^{2}(A)=\frac{\xi^{\sim} A^{\sim} A \xi}{\xi^{\sim} \xi} \leq \max _{x \neq 0} \frac{x^{\sim} A^{\sim} A x}{x^{\sim} x}
$$

Point 2) Analogous considerations as in Point 1) lead to the conclusion.
Proof of Lemma 2.17 The conclusion is straightforward if $\operatorname{rank}[A]=0$. Then, let $\operatorname{rank}[A] \neq 0$ so that $\bar{\sigma}(A)>0$ and there exists $x \neq 0$ such that $A^{\sim} A x=\bar{\sigma}^{2}(A) x$. From this relation it follows that $A A^{\sim} y=\bar{\sigma}^{2}(A) y$ with $y:=A x$. Vector $y$ is nonzero otherwise $\bar{\sigma}^{2}(A) x=0$. Hence, $\sigma^{2}(A)$ is an eigenvalue of $A A^{\sim}$ and then $\bar{\sigma}^{2}\left(A^{\sim}\right) \geq$ $\bar{\sigma}^{2}(A)$. The same line of reasoning applied to $A^{\sim}$ instead of $A$ leads to the conclusion that $\bar{\sigma}(A) \geq \bar{\sigma}\left(A^{\sim}\right)$. Hence, the thesis follows.

Proof of Lemma 2.18 Points 1) and 2) Let $x^{i}$ be an eigenvector associated with $\lambda_{i}(A)$. Then, recalling Lemma 2.16

$$
\underline{\sigma}^{2}(A)=\min _{x \neq 0} \frac{x^{\sim} A^{\sim} A x}{x^{\sim} x} \leq \frac{x^{i \sim} A^{\sim} A x^{i}}{x^{i \sim} x^{i}} \leq \max _{x \neq 0} \frac{x^{\sim} A^{\sim} A x}{x^{\sim} x}=\bar{\sigma}^{2}(A)
$$

Finally

$$
\frac{x^{i \sim A^{\sim} A x^{i}}}{x^{i \sim} x^{i}}=\frac{x^{i \sim} \lambda_{i}^{\sim}(A) \lambda_{i}(A) x^{i}}{x^{i \sim} x^{i}}=\left|\lambda_{i}(A)\right|^{2} \leq r_{s}^{2}(A)
$$

leads to the conclusion.
Point 3) Recall that the eigenvalues of an inverse matrix are the inverses of the eigenvalues of the matrix. Then,

$$
\underline{\sigma}^{2}(A)=\min _{i} \lambda_{i}\left(A^{\sim} A\right)
$$

$$
\begin{aligned}
& =\min _{i} \frac{1}{\lambda_{i}\left(\left(A^{\sim} A\right)^{-1}\right)} \\
& =\frac{1}{\max _{i} \lambda_{i}\left(\left(A^{\sim} A\right)^{-1}\right)}=\frac{1}{\bar{\sigma}^{2}\left(A^{-1}\right)}
\end{aligned}
$$

where the last equality follows from Lemma 2.17.
Point 4) The correctness of the claim directly follows by interchanging $A$ and $A^{-1}$ in Point 3).

Point 5) By exploiting the assumption that $A$ is hermitian and recalling that $\lambda_{i}\left(A^{2}\right)=\lambda_{i}^{2}(A)$, it follows

$$
\bar{\sigma}^{2}(A)=\max _{i} \lambda_{i}\left(A^{\sim} A\right)=\max _{i} \lambda_{i}\left(A^{2}\right)=\max _{i} \lambda_{i}^{2}(A)=r_{s}^{2}(A)
$$

Proof of Lemma 2.19 It turns out that

$$
\sigma_{i}^{2}(\alpha A)=\lambda_{i}\left(\alpha^{\sim} A^{\sim} A \alpha\right)=|\alpha|^{2} \lambda_{i}\left(A^{\sim} A\right)=|\alpha|^{2} \sigma_{i}^{2}(A)
$$

so that the thesis directly follows.
Proof of Lemma 2.20 By recalling Lemma 2.16 and well known properties on the norm of a vector, it results

$$
\begin{aligned}
\bar{\sigma}(A+B) & =\max _{x \neq 0} \frac{\|(A+B) x\|}{\|x\|} \\
& \leq \max _{x \neq 0} \frac{\|A x\|+\|B x\|}{\|x\|} \\
& \leq \max _{x \neq 0} \frac{\|A x\|}{\|x\|}+\max _{x \neq 0} \frac{\|B x\|}{\|x\|}=\bar{\sigma}(A)+\bar{\sigma}(B)
\end{aligned}
$$

Proof of Lemma 2.21 Preliminarily, by recalling Lemma 2.16 and the fact that in general $\operatorname{Im}[B]$ may not coincide with $C^{m}$, it follows

$$
\bar{\sigma}(A)=\max _{x \neq 0} \frac{\|A x\|}{\|x\|} \geq \max _{z \neq 0} \frac{\|A B z\|}{\|B z\|} \geq \frac{\|A B z\|}{\|B z\|}, \forall z
$$

Hence

$$
\bar{\sigma}(A B)=\max _{z \neq 0} \frac{\|A B z\|}{\|z\|} \leq \max _{z \neq 0} \frac{\bar{\sigma}(A)\|B z\|}{\|z\|}=\bar{\sigma}(A) \bar{\sigma}(B)
$$

Proof of Lemma 2.22 Preliminarily, observe that, $\forall x$

$$
\frac{\|A x\|+\|B x\|}{\|x\|} \leq \frac{\|A x\|}{\|x\|}+\max _{x \neq 0} \frac{\|B x\|}{\|x\|}=\frac{\|A x\|}{\|x\|}+\bar{\sigma}(B)
$$

and that

$$
\frac{\|A x\|-\|B x\|}{\|x\|} \geq \frac{\|A x\|}{\|x\|}-\max _{x \neq 0} \frac{\|B x\|}{\|x\|}=\frac{\|A x\|}{\|x\|}-\bar{\sigma}(B)
$$

From these expressions it follows that

$$
\begin{aligned}
& \min _{x \neq 0} \frac{\|A x\|+\|B x\|}{\|x\|} \leq \underline{\sigma}(A)+\bar{\sigma}(B) \\
& \min _{x \neq 0} \frac{\|A x\|-\|B x\|}{\|x\|} \geq \underline{\sigma}(A)-\bar{\sigma}(B)
\end{aligned}
$$

As for the second inequality, it turns out

$$
\begin{aligned}
\underline{\sigma}(A+B) & =\min _{x \neq 0} \frac{\|(A+B) x\|}{\|x\|} \\
& =\min _{x \neq 0} \frac{\|A x+B x\|}{\|x\|} \\
& \leq \min _{x \neq 0} \frac{\|A x\|+\|B x\|}{\|x\|} \leq \underline{\sigma}(A)+\bar{\sigma}(B)
\end{aligned}
$$

Analogously, as for the first inequality it turns out

$$
\begin{aligned}
\underline{\sigma}(A+B) & =\min _{x \neq 0} \frac{\|(A+B) x\|}{\|x\|} \\
& =\min _{x \neq 0} \frac{\|A x+B x\|}{\|x\|} \\
& \geq \min _{x \neq 0} \frac{\|A x\|-\|B x\|}{\|x\|} \geq \underline{\sigma}(A)-\bar{\sigma}(B)
\end{aligned}
$$

Proof of Lemma 2.23 As for the first inequality, Lemma 2.16 implies that

$$
\begin{aligned}
\bar{\sigma}([A B]) & =\max _{\left[\begin{array}{l}
x \\
y
\end{array}\right] \neq 0} \frac{\|A x+B y\|}{\left\|\left[\begin{array}{l}
x \\
y
\end{array}\right]\right\|} \geq \max _{x \neq 0} \frac{\|A x\|}{\|x\|}=\bar{\sigma}(A) \\
\bar{\sigma}([A B])= & \max \frac{\|A x+B y\|}{\left\|\left[\begin{array}{l}
x \\
y
\end{array}\right]\right\|} \geq \max _{y \neq 0} \frac{\|B y\|}{\|y\|}=\bar{\sigma}(B)
\end{aligned}
$$

As for the second inequality, from Lemma 2.16 it follows that

$$
\begin{aligned}
\bar{\sigma}\left(\left[\begin{array}{ll}
A B
\end{array}\right]\right) & \max _{\left[\begin{array}{l}
x \\
y
\end{array}\right] \neq 0} \frac{\|A x+B y\|}{\left\|\left[\begin{array}{l}
x \\
y
\end{array}\right]\right\|} \\
& \leq \max _{\left[\begin{array}{l}
x \\
y
\end{array}\right] \neq 0} \frac{\|A x\|+\|B y\|}{\left\|\left[\begin{array}{l}
x \\
y
\end{array}\right]\right\|} \\
& \leq \max \frac{\bar{\sigma}(A)\|x\|+\bar{\sigma}(B)\|y\|}{\left\|\left[\begin{array}{l}
x \\
y
\end{array}\right]\right\|} \\
& {\left[\begin{array}{l}
x \\
y
\end{array}\right] \neq 0 } \\
& \max [\bar{\sigma}(A), \bar{\sigma}(B)] \max ^{x} \frac{\|x\|+\|y\|}{\|\left[\begin{array}{l}
x \\
y
\end{array}\right] \neq 0} \|
\end{aligned}
$$

The thesis then follows since, $\forall\left[x^{\prime} y^{\prime}\right]^{\prime}$ it results

$$
\frac{\|x\|+\|y\|}{\left\|\left[\begin{array}{l}
x \\
y
\end{array}\right]\right\|} \leq \sqrt{2}
$$

Actually,

$$
\begin{aligned}
0 \leq(\|x\|-\|y\|)^{2} & =\|x\|^{2}+\|y\|^{2}-2\|x\|\|y\| \\
& =2\left(\|x\|^{2}+\|y\|^{2}\right)-\|x\|^{2}-\|y\|^{2}-2\|x\|\|y\| \\
& =2\left\|\left[\begin{array}{l}
x \\
y
\end{array}\right]\right\|^{2}-(\|x\|+\|y\|)^{2}
\end{aligned}
$$

Proof of Lemma 2.24 It turns out that

$$
\sum_{i} \sigma_{i}^{2}(A)=\sum_{i} \lambda_{i}\left(A^{\sim} A\right)=\operatorname{trace}\left[A^{\sim} A\right]
$$

Proof of Lemma 2.25 As for the first inequality, from Lemma 2.16 it follows

$$
\bar{\sigma}(A)=\max _{x \neq 0} \frac{\|A x\|}{\|x\|} \geq\left\|A e_{i}\right\|, \quad \forall i
$$

where $e_{i}$ is the $i$-th column of the identity matrix. But,

$$
\left\|A e_{i}\right\|=\sqrt{\sum_{j}\left|\{A\}_{j, i}\right|^{2}} \geq \max _{j}\left|\{A\}_{j, i}\right|
$$

As for the second inequality, recalling Lemma 2.24 it results

$$
\begin{aligned}
\bar{\sigma}^{2}(A) & \leq \sum_{h} \sigma_{h}^{2}(A) \\
& \leq \operatorname{trace}\left[A^{\sim} A\right] \\
& \leq \sum_{i}\left\{A^{\sim} A\right\}_{i, i} \\
& \leq m \max _{i}\left\{A^{\sim} A\right\}_{i, i} \\
& \leq m \max _{i} \sum_{j}\left|\{A\}_{i, j}\right|^{2} \\
& \leq m \max _{i} m \max _{j}\left|\{A\}_{i, j}\right|^{2}=m^{2} \max _{i, j}\left|\{A\}_{i, j}\right|^{2}
\end{aligned}
$$

Lemma B. 8 Let $A$ and $B$ be two matrices with dimensions $n \times m$ and $m \times n$, respectively. If $\lambda \neq 0$ is an eigenvalue of $A B$, then it is also an eigenvalue of $B A$.

Proof Let $x \neq 0$ such that $A B x=\lambda x \neq 0$. Hence $y:=B x \neq 0$, otherwise $\lambda x=0$. Finally, $B A y=B A B x=\lambda B x=\lambda y$, that is $\lambda$ is an eigenvalue of $B A$.

Lemma B. 9 Let $A$ and $B$ be two matrices with dimensions $n \times m$ and $m \times n$, respectively. Suppose also that 1 is not an eigenvalue of $A B$. Then
i) $I+A B(I-A B)^{-1}=(I-A B)^{-1}$
ii) $A B(I-A B)^{-1}=(I-A B)^{-1} A B=A(I-B A)^{-1} B$

Proof Preliminarily observe that $(I-A B)$ and $(I-B A)$ are nonsingular thanks to the assumption on the eigenvalue of $A B$ and Lemma B.8. Hence, the formulas are well defined.

Point i) It is

$$
\begin{aligned}
I+A B(I-A B)^{-1} & =(I-A B)(I-A B)^{-1}+A B(I-A B)^{-1} \\
& =(I-A B+A B)(I-A B)^{-1}=(I-A B)^{-1}
\end{aligned}
$$

Point ii) As for the first equality, it is

$$
\begin{aligned}
A B(I-A B)^{-1} & =-I+(I-A B)^{-1} \\
& =-(I-A B)^{-1}(I-A B)+(I-A B)^{-1} \\
& =(I-A B)^{-1}(-I+A B+I)=(I-A B)^{-1} A B
\end{aligned}
$$

As for the second equality, it is $B-B A B=(I-B A) B=B(I-A B)$, so that $B(I-A B)^{-1}=(I-B A)^{-1} B$, which implies $A B(I-A B)^{-1}=A(I-B A)^{-1} B$.

Lemma B. 10 Let $A=A^{\prime} \geq 0$ and $B=B^{\prime} \geq 0$ be two matrices with the same dimensions. Then the eigenvalues of $C:=A B$ are real and nonnegative.

Proof Let $\lambda$ be an eigenvalue of $C$ and $\xi \neq 0$ an associated eigenvector, i.e. $C \xi=$ $A B \xi=\lambda \xi, \quad \xi \neq 0$. From these expressions it follows that $B A B \xi=\lambda B \xi$ so that $\xi^{\sim} B A B \xi=\lambda \xi^{\sim} B \xi$. Being $A$ and $B$ positive semidefinite, the quantities $\alpha:=$ $\xi^{\sim} B A B \xi$ and $\beta:=\xi^{\sim} B \xi$ are both real and nonnegative. If $\beta \neq 0$ then $\lambda=\frac{\alpha}{\beta}$ is real and nonnegative. If $\beta=0$, then $B \xi=0$ and $\lambda=0$.

Lemma B. 11 Let $A=A^{\prime} \geq 0$ and $B=B^{\prime} \geq 0$ be two matrices with the same dimensions and $\gamma$ a positive scalar. If $r_{s}\left(\gamma^{-2} A B\right)<1$, then
i) $C:=\left(I-\gamma^{-2} A B\right)^{-1} A \geq 0$
ii) If $A>0$, then $C>0$ and $A^{-1}>\gamma^{-2} B$

Conversely, if $C \geq 0$, then

$$
\text { iii) } r_{s}\left(\gamma^{-2} A B\right)<1
$$

Proof Preliminarily observe that $I-\gamma^{-2} A B$ is singular if and only if $\gamma^{-2} A B$ has (at least) one eigenvalue equal to one. This is not possible thanks to the assumption. Hence, matrix $C$ is well defined. It is also symmetric in view of what has been shown in the proof of Lemma B. 9 .

Point i) By exploiting Lemma B. 9 and the definition of $C$, it is easy to check that

$$
\left(-I+\gamma^{-2} A B\right) C+C\left(-I+\gamma^{-2} B A\right)+2 A=0
$$

Thanks to the assumption on the spectral radius of $\gamma^{-2} A B$, it turns out that matrix $-I+\gamma^{-2} A B$ has all its eigenvalue in the open left half plane, i.e. it is stable. The
equation above is therefore a Lyapunov equation in the unknown $C$, with stable coefficient matrix and positive semidefinite known term $2 A$. In view of Lemma C. 1 one can conclude that $C$ is positive semidefinite.

Point ii) If $A>0$ then, recalling again Lemma C.1, it turns out that $C>0$. Moreover,

$$
C=\left(I-\gamma^{-2} A B\right)^{-1} A=\left[A\left(A^{-1}-\gamma^{-2} B\right)\right]^{-1} A=\left(A^{-1}-\gamma^{-2} B\right)^{-1}
$$

so that $\left(A^{-1}-\gamma^{-2} B\right)>0$, which is the thesis.
Point iii) Assume by contradiction that $\lambda \geq 1$ is an eigenvalue of $\gamma^{-2} A B$ (recall that, thanks to Lemma B. 10 the eigenvalues of $A B$ are real and nonnegative). Since $C=\left(I-\gamma^{-2} A B\right)^{-1} A$ exists, such an eigenvalue can not be equal to 1 . In view of Lemma B.8, $\lambda$ is also an eigenvalue of $\gamma^{-2} B A$ so that there exists a vector $\xi \neq 0$ such that $\gamma^{-2} B A \xi=\lambda \xi$, i.e. $\left(I-\gamma^{-2} B A\right) \xi=(1-\lambda) \xi$ or, alternatively, $\left(I-\gamma^{-2} B A\right)^{-1} \xi=$ $(1-\lambda)^{-1} \xi$. Hence $\xi^{\prime}(1-\lambda)^{-1}=\xi^{\prime}\left(I-\gamma^{-2} A B\right)^{-1}$. It follow that

$$
\frac{\xi^{\prime} A \xi}{(1-\lambda)}=\xi^{\prime}\left(I-\gamma^{-2} A B\right)^{-1} A \xi
$$

The left hand side of this equation is a nonpositive number (actually, the numerator is nonnegative and the denominator is negative). The right hand side is nonnegative in view of the assumption on $C$. Hence, it must be $A \xi=0$ and hence, from ( $I$ -$\left.\gamma^{-2} B A\right) \xi=(1-\lambda) \xi=\xi$ and $\lambda \neq 0$ it follows that $\xi=0$, a contradiction.

Lemma B. 12 Let $A=A^{\prime}>0$ and $B=B^{\prime} \geq 0$ be two matrices with the same dimensions and $\alpha$ a positive scalar. Then $r_{s}(B A) \geq \alpha$ if and only if there exists $x \neq 0$ such that $x^{\prime}\left(B-\alpha A^{-1}\right) x \geq 0$.

Proof If $r_{s}(B A) \geq \alpha$, it must exist (recall Lemma B.10) a vector $\xi \neq 0$ such that $B A \xi=\lambda \xi, \lambda \geq \alpha$. Letting $x:=A \xi$, it follows that $B x=\lambda A^{-1} x$ and also $x^{\prime} B x=$ $\lambda x^{\prime} A^{-1} x \geq \alpha x^{\prime} A^{-1} x$, which yields $x^{\prime}\left(B-\alpha A^{-1}\right) x \geq 0$.

Conversely, if $r_{s}(B A)<\alpha$, then, from Lemma B. 11 it follows $A^{-1}>\alpha^{-1} B$ so that $x^{\prime}\left(B-\alpha A^{-1}\right) x<0, \forall x \neq 0$ and hence it does not exist $x \neq 0$ such that $x^{\prime}\left(B-\alpha A^{-1}\right) x \geq 0$.

Lemma B. 13 Consider the block matrix

$$
\Phi=\left[\begin{array}{ll}
\Phi_{1} & \Phi_{2} \\
\Phi_{3} & \Phi_{4}
\end{array}\right]
$$

where the submatrices $\Phi_{1}$ and $\Phi_{4}$ are square. Then
i) If $\Phi_{1}$ is nonsingular, $\operatorname{det}[\Phi]=\operatorname{det}\left[\Phi_{1}\right] \operatorname{det}\left[\Phi_{4}-\Phi_{3} \Phi_{1}^{-1} \Phi_{2}\right]$
ii) If $\Phi_{4}$ is nonsingular, $\operatorname{det}[\Phi]=\operatorname{det}\left[\Phi_{4}\right] \operatorname{det}\left[\Phi_{1}-\Phi_{2} \Phi_{4}^{-1} \Phi_{3}\right]$

Proof Point i) The identity

$$
\left[\begin{array}{ll}
\Phi_{1} & \Phi_{2} \\
\Phi_{3} & \Phi_{4}
\end{array}\right]=\left[\begin{array}{cc}
I & 0 \\
\Phi_{3} \Phi_{1}^{-1} & I
\end{array}\right]\left[\begin{array}{cc}
\Phi_{1} & \Phi_{2} \\
0 & \Phi_{4}-\Phi_{3} \Phi_{1}^{-1} \Phi_{2}
\end{array}\right]
$$

is straightforward. By considering the determinant of the matrices on the left and right hand sides, point $i$ ) follows.

Point ii) The identity

$$
\left[\begin{array}{ll}
\Phi_{1} & \Phi_{2} \\
\Phi_{3} & \Phi_{4}
\end{array}\right]=\left[\begin{array}{cc}
I & \Phi_{2} \Phi_{4}^{-1} \\
0 & I
\end{array}\right]\left[\begin{array}{cc}
\Phi_{1}-\Phi_{2} \Phi_{4}^{-1} \Phi_{3} & 0 \\
\Phi_{3} & \Phi_{4}
\end{array}\right]
$$

is straightforward. By considering the determinant of the matrices on the left and right hand sides, point $i i$ ) follows.

Lemma B. 14 (Schur complements) Consider the block symmetric matrix

$$
\Phi=\left[\begin{array}{ll}
\Phi_{1} & \Phi_{2} \\
\Phi_{2}^{\prime} & \Phi_{3}
\end{array}\right]
$$

where the submatrices $\Phi_{1}$ and $\Phi_{3}$ are square. Then
i) If $\Phi_{1}$ is positive definite, $\Phi \geq 0 \Longleftrightarrow \Phi_{3} \geq \Phi_{2}^{\prime} \Phi_{1}^{-1} \Phi_{2}$
ii) If $\Phi_{3}$ is positive definite, $\Phi \geq 0 \Longleftrightarrow \Phi_{1} \geq \Phi_{2} \Phi_{3}^{-1} \Phi_{2}^{\prime}$

Proof Point i) Considering the nonsingular matrix

$$
T=\left[\begin{array}{cc}
I & 0 \\
\Phi_{2}^{\prime} \Phi_{1}^{-1} & I
\end{array}\right]
$$

it is straightforward to verify the identity

$$
\Phi=T\left[\begin{array}{cc}
\Phi_{1} & 0 \\
0 & \Phi_{3}-\Phi_{2}^{\prime} \Phi_{1}^{-1} \Phi_{2}
\end{array}\right] T^{\prime}
$$

from which point $i$ ) follows.
Point ii) Considering the nonsingular matrix

$$
T=\left[\begin{array}{cc}
I & \Phi_{2} \Phi_{3}^{-1} \\
0 & I
\end{array}\right]
$$

it is straightforward to verify the identity

$$
\Phi=T\left[\begin{array}{cc}
\Phi_{1}-\Phi_{2} \Phi_{3}^{-1} \Phi_{2}^{\prime} & 0 \\
0 & \Phi_{3}
\end{array}\right] T^{\prime}
$$

from which point $i i$ ) follows.
Lemma B. 15 Consider $A$ and $D$ real matrices such that $D^{\prime} D=I$. Then

$$
\min _{Z}\|A-D Z\|=\left\|\left(I-D D^{\prime}\right) A\right\|
$$

and the optimal solution is $Z^{\circ}=D^{\prime} A$.
Proof Define $\Delta:=Z-Z^{o}=Z-D^{\prime} A$. From the fact that matrices $D$ and $\left(I-D D^{\prime}\right)$ are orthogonal we have

$$
\begin{aligned}
\|A-D Z\|^{2} & =\left\|\left(I-D D^{\prime}\right) A-D \Delta\right\|^{2} \\
& =\max _{\|x\|=1}\left\{\left\|\left(I-D D^{\prime}\right) A x\right\|^{2}+\|\Delta x\|^{2}\right\}
\end{aligned}
$$

However, due to

$$
\begin{aligned}
\min _{Z}\|A-D Z\|^{2} & =\min _{\Delta} \max _{\|x\|=1}\left\{\left\|\left(I-D D^{\prime}\right) A x\right\|^{2}+\|\Delta x\|^{2}\right\} \\
& \geq \max _{\|x\|=1} \min _{\Delta}\left\{\left\|\left(I-D D^{\prime}\right) A x\right\|^{2}+\|\Delta x\|^{2}\right\} \\
& \geq \max _{\|x\|=1}\left\|\left(I-D D^{\prime}\right) A x\right\|^{2} \\
& \geq\left\|\left(I-D D^{\prime}\right) A\right\|^{2}
\end{aligned}
$$

the lemma is proved by simple verification that the equality holds for $Z=Z^{\circ}$.
Remark B. 1 From Lemma B. 15 it follows that the linear equation $A-D Z=0$, with $D^{\prime} D=I$, admits a solution namely $Z=D^{\prime} A$ if and only if $\left(I-D D^{\prime}\right) A=0$.

## Appendix C

## Riccati Equation

The Lyapunov and Riccati equations play an important role in the analysis and control of linear time-invariant systems. Here, for a given $n$-th dimensional system $\Sigma(A, B, C, D)$, a few basic results on these equations are reported. Standard references are the book by Bittanti et al. [7] and the paper by Doyle et al. [17].

Lemma C. 1 (Extended Lyapunov lemma) Consider the Lyapunov equation

$$
0=P A+A^{\prime} P+C^{\prime} C
$$

Then
i) If $A$ is stable, there exists a unique solution. Such a solution is symmetric and positive semidefinite.
ii) If the pair $(A, C)$ is detectable and there exists a symmetric and positive semidefinite solution, $A$ is stable.
iii) If the pair $(A, C)$ is observable, the solution at point i) is actually positive definite.

Proof Points i) and iii) If $A$ is stable, then

$$
P:=\int_{0}^{\infty} e^{A^{\prime} t} C^{\prime} C e^{A t} d t
$$

is well defined, symmetric and positive semidefinite. If $(A, C)$ is observable, then $P$ is positive definite. In fact, if, by contradiction, $P x_{0}=0, x_{0} \neq 0$, then

$$
x_{0}^{\prime} P x_{0}:=\int_{0}^{\infty} x_{0}^{\prime} e^{A^{\prime} t} C^{\prime} C e^{A t} x_{0} d t=\int_{0}^{\infty} y^{\prime} y d t=0
$$

where $y$ is the free output of $\Sigma$ with initial state $x(0)=x_{0}$. Hence, $y(t)=0, \forall t \geq 0$ contradicts the observability assumption.

Based on the definition of $P$, it follows

$$
A^{\prime} P+P A=\int_{0}^{\infty} \frac{d}{d t}\left(e^{A^{\prime} t} C^{\prime} C e^{A t}\right) d t=-C^{\prime} C
$$

Matrix $P$ is the unique solution of the Lyapunov equation. In fact, if $\hat{P}$ is any other solution, it results

$$
0=(P-\hat{P}) A+A^{\prime}(P-\hat{P})
$$

Thanks to the stability of $A$ and a well known result of linear algebra, this equation admits the unique solution $P-\hat{P}=0$.

Point ii) Assume that $\bar{P}=\bar{P}^{\prime} \geq 0$ is a solution of the Lyapunov equation and, by contradiction, that $A$ is not stable, i.e. $A x=\lambda x, x \neq 0, \operatorname{Re}(\lambda) \geq 0$. From the equation it then follows

$$
0=x^{\sim} \bar{P} A x+x^{\sim} A^{\prime} \bar{P} x+x^{\sim} C^{\prime} C x=2 \operatorname{Re}(\lambda) x^{\sim} \bar{P} x+x^{\sim} C^{\prime} C x
$$

Since the last term of this equation is the sum of two nonnegative elements, it turns out that $C x=0$. But $A x=\lambda x, C x=0, x \neq 0, R e(\lambda) \geq 0$ contradicts the assumed detectability of ( $A, C$ ) (recall Lemma D.2).

Lemma C. 2 (Stabilizing solution of the Riccati equation - 1) Consider the Riccati equation

$$
0=P A+A^{\prime} P+P R P+Q
$$

with $R$ and $Q$ real and assume that
a) The matrix

$$
Z:=\left[\begin{array}{rr}
A & R \\
-Q & -A^{\prime}
\end{array}\right]
$$

does not have eigenvalues lying on the imaginary axis.
b) Matrices $R$ and $Q$ are symmetric and $n$-dimensional.

If the subspace

$$
\mathcal{X}:=\operatorname{Im}[L], \quad L:=\left[\begin{array}{l}
X  \tag{C.1}\\
Y
\end{array}\right]
$$

generated by the (generalized) eigenvectors associated with the negative real part eigenvalues of $Z$ is complementary to the $n$-dimensional subspace

$$
\mathcal{I}:=\operatorname{Im}[\bar{I}], \quad \bar{I}:=\left[\begin{array}{l}
0  \tag{C.2}\\
I
\end{array}\right]
$$

then
i) The matrix $P_{s}:=Y X^{-1}$ is a real, symmetric and stabilizing solution, namely is such that $\left(A+R P_{s}\right)$ is stable.
ii) $P_{s}$ is the unique stabilizing solution.

Proof $Z$ is an Hamiltonian matrix, i.e. it satisfies

$$
\begin{equation*}
J Z+Z^{\prime} J=0 \tag{C.3}
\end{equation*}
$$

where

$$
J:=\left[\begin{array}{rr}
0 & I \\
-I & 0
\end{array}\right]
$$

is the so called sympletic matrix. Matrix $Z$ has eigenvalues symmetric with respect to the imaginary axis, since $Z=-J^{-1} Z^{\prime} J$. Assumption a) assures the existence of $n$
eigenvalues of $Z$ in the open left half plane. It is well known that the subspace $\mathcal{X}$ is $Z$-invariant so that there exists a real matrix $T$ with eigenvalues in the open left half plane (stable) such that

$$
\begin{equation*}
Z L=L T \tag{C.4}
\end{equation*}
$$

where $L$ satisfies eq. (C.1). The proof is divided into four main steps.
Step 1 It is here proved that $X^{\sim} Y=Y^{\sim} X$. Actually, letting $W:=X^{\sim} Y-Y^{\sim} X$, it is immediate to check that $W=L^{\sim} J L$ so that, exploiting eqs. (C.3), (C.4), it follows $W T=L^{\sim} J Z L=-L^{\sim} Z^{\prime} J L=-T^{\prime} W$, i.e.

$$
W T+T^{\prime} W=0
$$

Since $T$ is stable, this Lyapunov equation admits the unique solution $W=0$, which is the claim.

Step 2 Matrix $X$ is invertible since the subspace $\mathcal{X}$ is complementary to the $n$ dimensional subspace $\mathcal{I}$ (see eq. (C.2)), so that $P_{s}=Y X^{-1}$ is well defined.

Step 3 It is here proved that $P_{s}$ is a real, symmetric and stabilizing solution of the Riccati equation.

As for reality, observe that the columns of $L=\left[X^{\sim} Y^{\sim}\right]^{\sim}$ can be chosen complex conjugate in pair, so that, if $X_{c}$ and $Y_{c}$ are the complex conjugate of $X$ and $Y$, respectively, then

$$
\left[\begin{array}{ll}
X_{c}^{\sim} & Y_{c}^{\sim}
\end{array}\right]^{\sim} \Gamma=\left[\begin{array}{ll}
X^{\sim} & Y^{\sim}
\end{array}\right]^{\sim}
$$

where $\Gamma$ is a permutation matrix. Hence $P_{s}=Y X^{-1}=Y_{c} \Gamma \Gamma^{-1} X_{c}^{-1}=Y_{c} X_{c}^{-1}$ so that $P_{s}$ coincides with its conjugate, i.e. $P_{s}$ is real.

As for symmetry, observe that, thanks to the identity (proved in Step 1) $X^{\sim} Y=$ $Y^{\sim} X$ it follows that $P_{s}^{\prime}=P_{s}^{\sim}=\left(X^{-1}\right)^{\sim} Y^{\sim}=\left(X^{-1}\right)^{\sim} X^{\sim} Y X^{-1}=Y X^{-1}=P_{s}$.

As for the stabilizing property, from eq. (C.4) it follows that $A X+R Y=A X+$ $R Y X^{-1} X=\left(A+R P_{s}\right) X=X T$ so that $A+R P_{s}=X T X^{-1}$. Therefore, $A+R P_{s}$ and $T$ are similar. Since $T$ is stable, matrix $A+R P_{s}$ is stable as well.

It is now shown that $P_{s}$ is a solution of the equation. In fact, eq. (C.4) is equivalent to

$$
\begin{aligned}
A X+R Y & =X T \\
-Q X-A^{\prime} Y & =Y T
\end{aligned}
$$

Premultiplying the first equation by $Y^{\sim}$ yields $Y^{\sim} A X+Y^{\sim} R Y=Y^{\sim} X T$. Moreover, premultiplying the second equation by $X^{\sim}$ gives $-X^{\sim} Q X-X^{\sim} A^{\prime} Y=X^{\sim} Y T$. Recalling that $X^{\sim} Y=Y^{\sim} X$ it then follows that

$$
0=X^{\sim} Q X+X^{\sim} A^{\prime} Y+Y^{\sim} A X+Y^{\sim} R Y
$$

Postmultiplying this equation by $X^{-1}$ and premultiplying it by $\left(X^{-1}\right)^{\sim}$ it follows

$$
\begin{aligned}
0 & =Q+A^{\prime}\left(Y X^{-1}\right)+\left(Y X^{-1}\right)^{\sim} A+\left(Y X^{-1}\right)^{\sim} R\left(Y X^{-1}\right) \\
& =Q+A^{\prime} P_{s}+P_{s}^{\sim} A+P_{s}^{\sim} R P_{s}
\end{aligned}
$$

Since it has been already shown that $P_{s}$ is hermitian $\left(P_{s}=P_{s}^{\sim}\right)$, the conclusion is straightforwardly derived.

Step 4 As for the uniqueness of the stabilizing solution, assume, by contradiction, that there exist two stabilizing solutions, namely $P_{1}$ and $P_{2}$. Then

$$
\begin{aligned}
& 0=P_{1} A+A^{\prime} P_{1}+P_{1} R P_{1}+Q \\
& 0=P_{2} A+A^{\prime} P_{2}+P_{2} R P_{2}+Q
\end{aligned}
$$

Subtracting term by term, it is easy to see that

$$
0=\left(P_{1}-P_{2}\right)\left(A+R P_{1}\right)+\left(A+R P_{2}\right)^{\prime}\left(P_{1}-P_{2}\right)
$$

This equation can be considered as a linear equation in the unknown $P_{1}-P_{2}$. Being both $A+R P_{1}$ and $A+R P_{2}$ stable, a well known result of linear algebra implies that there exists only one solution, precisely $P_{1}-P_{2}=0$.

Lemma C. 3 (Stabilizing solution of the Riccati equation - 2) Consider the Riccati equation

$$
0=P A+A^{\prime} P+P R P+Q
$$

and suppose that matrix $R$ is square, $n$-dimensional and either equal to $B B^{\prime}$ or $-B B^{\prime}$. Moreover, assume that
a) The pair $(A, B)$ is stabilizable
b) The matrix

$$
Z:=\left[\begin{array}{rr}
A & R \\
-Q & -A^{\prime}
\end{array}\right]
$$

does not have eigenvalues lying on the imaginary axis
c) The matrix $Q$ is real and symmetric

Then,
i) The subspace

$$
\mathcal{X}:=\operatorname{Im}[L], \quad L:=\left[\begin{array}{l}
X \\
Y
\end{array}\right]
$$

generated by the (generalized) eigenvectors of $Z$ associated with the negative real part is complementary to the $n$-dimensional subspace

$$
\operatorname{Im}[\bar{I}], \quad \bar{I}:=\left[\begin{array}{l}
0 \\
I
\end{array}\right]
$$

ii) $P_{s}:=Y X^{-1}$ is a solution of the Riccati equation. It is real, symmetric and stabilizing, i.e. such that $\left(A+R P_{s}\right)$ is stable.
iii) $P_{s}$ is the unique stabilizing solution of the Riccati equation.

Proof In view of Lemma C.2, the proof of points $i i$ ) and $i i i$ ) is straightforward once point $i$ ) has been proved, which, in turn, derives once the invertibility of $X$ is proved. To this aim, consider eq. (C.4) and assume, by contradiction, that there exists $\xi \neq 0$ such that $X \xi=0$, i.e. $\xi \in \operatorname{Ker}[X]$. It follows in particular that

$$
\begin{equation*}
A X+R Y=X T \tag{C.5}
\end{equation*}
$$

Premultiplying this equation by $Y^{\sim}$ and recalling that $Y^{\sim} X=X^{\sim} Y$, it follows $Y^{\sim} A X+Y^{\sim} R Y=X^{\sim} Y T$ so that $\xi^{\sim}\left(Y^{\sim} A X+Y^{\sim} R Y-X^{\sim} Y T\right) \xi=\xi^{\sim} Y^{\sim} R Y \xi=0$, namely (recall that $R= \pm B B^{\prime}$ )

$$
\begin{equation*}
R Y \xi=0, \quad \forall \xi \in \operatorname{Ker}[X] \tag{C.6}
\end{equation*}
$$

This equation together with eq. (C.5) yields

$$
\begin{equation*}
X T \xi=0, \quad \forall \xi \in \operatorname{Ker}[X] \tag{C.7}
\end{equation*}
$$

On the other hand, eq. (C.4) entails that $-Q X-A^{\prime} Y=Y T$ so that

$$
\begin{equation*}
-A^{\prime} Y \xi=Y T \xi, \quad \forall \xi \in \operatorname{Ker}[X] \tag{C.8}
\end{equation*}
$$

For each $\xi \in \operatorname{Ker}[X]$ eqs. (C.6)-(C.8) hold. In particular, eq. (C.7) implies $T \xi \in$ $\operatorname{Ker}[X]$ so that $T^{i} \xi \in \operatorname{Ker}[X], \forall i \geq 0$. Now it is proved, by induction, that, from eqs. (C.7),(C.8), it follows

$$
\begin{equation*}
R\left(A^{\prime}\right)^{k} Y \xi=0, \quad \forall k \geq 0, \quad \forall \xi \in \operatorname{Ker}[X] \tag{C.9}
\end{equation*}
$$

Actually, eq. (C.9) is true for $k=0$ (eq. (C.6)) and, moreover, it results (see eq. (C.8))

$$
R\left(A^{\prime}\right)^{k+1} Y \xi=R\left(A^{\prime}\right)^{k} A^{\prime} Y \xi=-R\left(A^{\prime}\right)^{k} Y T \xi=0
$$

since $T \xi \in \operatorname{Ker}[X]$. Recalling that $T^{i} \xi \in \operatorname{Ker}[X], \forall i \geq 0$, from eq. (C.8) it follows $Y T^{k} \xi=Y T T^{k-1} \xi=-A^{\prime} Y T^{k-1} \xi$. By repeatedly using this last equation, one obtains

$$
\begin{equation*}
Y T^{k} \xi=\left(-A^{\prime}\right)^{k} Y \xi, \quad \forall k \geq 0, \quad \forall \xi \in \operatorname{Ker}[X] \tag{C.10}
\end{equation*}
$$

For each polynomial $\nu(s)$ it then follows $\forall \xi \in \operatorname{Ker}[X]$

$$
\begin{align*}
\nu\left(-A^{\prime}\right) Y \xi & =Y \nu(T) \xi  \tag{C.11}\\
R \nu\left(-A^{\prime}\right) Y \xi & =0  \tag{C.12}\\
\nu(T) \xi & \in \operatorname{Ker}[X] \tag{C.13}
\end{align*}
$$

Let now $\nu_{m}(s)$ be the monic polynomial of minimum degree such that $\nu_{m}(T) \xi=0$. Notice that such a polynomial actually exists and its degree is not greater than the degree of the minimal polynomial $\varphi(s)$ of $T$ (in fact $\varphi(T)=0$ ). If $\lambda$ is a root of $\nu_{m}(s)$ it follows that $(\lambda-s) \mu(s)=\nu_{m}(s)$, with $\operatorname{deg}[\mu(s)]<\operatorname{deg}\left[\nu_{m}(s)\right]$. Then,

$$
\begin{equation*}
\nu_{m}(T) \xi=(\lambda I-T) \mu(T) \xi=0 \tag{C.14}
\end{equation*}
$$

Observe that $\beta:=\mu(T) \xi \neq 0$ due to the minimallity of $\nu_{m}(s)$, so that $\beta$ is an eigenvector of $T$ and $\lambda$, being an associated eigenvalue, must have negative real part (recall that $T$ is stable). By exploiting eq. (C.14) and letting $\nu(s)=\nu_{m}(s)$ in eq. (C.11), it turns out that, for $\xi \in \operatorname{Ker}[X]$

$$
\nu_{m}\left(-A^{\prime}\right) Y \xi=\left(\lambda I+A^{\prime}\right) \mu\left(-A^{\prime}\right) Y \xi=Y \nu_{m}(T) \xi=0
$$

so that

$$
A^{\prime} \mu\left(-A^{\prime}\right) Y \xi=-\lambda \mu\left(-A^{\prime}\right) Y \xi
$$

with $\operatorname{Re}(-\lambda)>0$, whereas eq. (C.12), letting $\nu(s)=\mu(s)$, becomes

$$
R \mu\left(-A^{\prime}\right) Y \xi= \pm B B^{\prime} \mu\left(-A^{\prime}\right) Y \xi=0
$$

which, in turn, implies $B^{\prime} \mu\left(-A^{\prime}\right) Y \xi=0$. Letting $\eta:=\mu\left(-A^{\prime}\right) Y \xi$, it then follows

$$
\begin{aligned}
& A^{\prime} \eta=-\lambda \eta \\
& B^{\prime} \eta=0
\end{aligned}
$$

The assumption of stabilizability of the pair $(A, B)$ and $\operatorname{Re}(-\lambda)>0$, implies, in view of the $P B H$ test (Lemma D.4), that $\eta=0$. From eq. (C.11) with $\nu(s)=\mu(s)$ it then follows that $Y \mu(T) \xi=0$ so that $\mu(T) \xi \in \operatorname{Ker}[Y]$. On the other hand, $\mu(T) \xi \in \operatorname{Ker}[X]$, so that $L \beta=0$ contradicts the fact that $\operatorname{Im}[L]$ is a $n$ dimensional subspace. Hence $X$ is invertible and point $i$ ) is proved.

The preceding lemma allows proving the following result which settles in the context of the theory of optimal control with quadratic cost functionals.
Lemma C. 4 Consider the Riccati equation

$$
0=P A+A^{\prime} P+P R P+C^{\prime} C, R:=-B B^{\prime}
$$

If $\bar{P}$ is a solution, then

$$
\operatorname{Ker}[\bar{P}] \subset \mathcal{X}_{n o}:=\operatorname{Ker}[C] \cap \operatorname{Ker}[C A] \cdots \cap \operatorname{Ker}\left[C A^{n-1}\right]
$$

Moreover, there exists a solution $P_{s}$ which is real, symmetric, stabilizing and positive semidefinite if and only if
a) There do not exist eigenvalues of the unobservable part of the pair $(A, C)$ lying on the imaginary axis.
b) The pair $(A, B)$ is stabilizable

Proof It is first proved that the kernel of $\bar{P}$ is contained in the unobservable subspace $\mathcal{X}_{n o}$ of $(A, C)$. Let $\xi \in \operatorname{Ker}[\bar{P}]$, so that $\bar{P} \xi=0$. From the Riccati equation it follows that $0=\xi^{\sim}\left(\bar{P} A+A^{\prime} \bar{P}+\bar{P} R \bar{P}+C^{\prime} C\right) \xi=\xi^{\sim} C^{\prime} C \xi$ so that $C \xi=0$. Taking this into account, again from the Riccati equation one obtains $0=\left(\bar{P} A+A^{\prime} \bar{P}+\bar{P} R \bar{P}+C^{\prime} C\right) \xi=$ $\bar{P} A \xi$ so that $A \xi \in \operatorname{Ker}[\bar{P}]$. Hence it has been proved that $\xi \in \operatorname{Ker}[\bar{P}] \Rightarrow \xi \in \operatorname{Ker}[C]$ and $A \xi \in \operatorname{Ker}[\bar{P}]$. By repeating again this argument for $A \xi$ one obtains $A^{2} \xi \in \operatorname{Ker}[\bar{P}]$ and $\xi \in \operatorname{Ker}[C A]$ so that, in conclusion

$$
\operatorname{Ker}[\bar{P}] \subset:=\operatorname{Ker}[C] \cap \operatorname{Ker}[C A] \cdots \cap \operatorname{Ker}\left[C A^{n-1}\right]=\mathcal{X}_{n o}
$$

It is now proved that assumptions $a$ ), $b$ ) imply condition $b$ ) of Lemma C.3. Assume, by contradiction, that the Hamiltonian matrix $Z$ associated with the Riccati equation has an eigenvalue $\lambda=j \omega$ (on the imaginary axis), i.e. $Z z=j \omega z$, with $z:=\left[x^{\sim} y^{\sim}\right]^{\sim} \neq 0$. Hence,

$$
\begin{align*}
A x+R y & =j \omega x  \tag{C.15}\\
-Q x-A^{\prime} y & =j \omega y \tag{C.16}
\end{align*}
$$

where $Q=C^{\prime} C$. Now add to the first equation, premultiplied by $y^{\sim}$, the second, premultiplied by $x^{\sim}$. It follows

$$
\begin{equation*}
j \omega\left(x^{\sim} y+y^{\sim} x\right)-y^{\sim} A x+x^{\sim} A^{\prime} y=y^{\sim} R y-x^{\sim} Q x \tag{C.17}
\end{equation*}
$$

The left hand side is a purely imaginary number, whereas the right hand side is real. Hence, both sides must be equal to zero. In particular, being $Q \geq 0$ and $R \leq 0$, it follows that $Q x=0$ and $R y=0$, which in turn imply that $C x=0$ and $B^{\prime} y=0$. These last equations, together with eqs. (C.15)-(C.17) entail

$$
\begin{align*}
A x & =j \omega x  \tag{C.18}\\
C x & =0  \tag{C.19}\\
A^{\prime} y & =-j \omega y  \tag{C.20}\\
B^{\prime} y & =0 \tag{C.21}
\end{align*}
$$

In conclusion, since either $x$ or $y$ is different from zero, the above eqs. (C.18)-(C.19) violate either assumption $a$ ) (recall Lemma D.1) if $x \neq 0$ or the stabilizability assumption b) (recall Lemma D.4) if $y \neq 0$. Assumptions $a$ )-c) of Lemma C. 3 are so verified. Such lemma assures the existence of a real and symmetric solution $P_{s}$ such that $A+R P_{s}$ is stable. Obviously, $P_{s}$ is a solution of the Lyapunov equation (in the unknown $P$ )

$$
0=P A+A^{\prime} P+Q+P_{s} R P_{s}
$$

which can be also rewritten as

$$
0=P\left(A+R P_{s}\right)+\left(A+R P_{s}\right)^{\prime} P-P_{s} R P_{s}+Q
$$

This is a Lyapunov equation whose coefficient matrix $A+R P_{s}$ is stable and the known term $Q-P_{s} R P_{s}$ is positive semidefinite. Thanks to Lemma C. 1 this equation admits a unique solution, $P_{s}$, which is positive semidefinite.

Conversely, assume that there exists a real, symmetric and stabilizing solution $P_{s}$ of the Riccati equation. Necessity of condition b) is then obvious. Now suppose by contradiction that condition a) does not hold, i.e. $A x=j \omega x, C x=0$ and $x \neq 0$. Hence $Z z=j \omega z$ where $z:=\left[x^{\prime} 0\right] \neq 0$ and

$$
Z:=\left[\begin{array}{rr}
A & -B B^{\prime} \\
-C^{\prime} C & -A^{\prime}
\end{array}\right]
$$

Hence, the Hamiltonian matrix $Z$ has an eigenvalue on the imaginary axis, so contradicting the existence of a stabilizing solution of the Riccati equation.

Lemma C. 5 Let $A, B, C$ and $D$ be four matrices with dimensions $n \times n, n \times m$, $p \times n$ and $p \times m$, respectively. Assume also that $C^{\prime} D=0, D^{\prime} D=I$ and there exists the symmetric and stabilizing solution $P_{s}$ of the Riccati equation

$$
0=P A+A^{\prime} P-P B B^{\prime} P+C^{\prime} C
$$

Let $A_{c c}:=A-B B^{\prime} P_{s}$ and $C_{1 c}:=C-D B^{\prime} P_{s}$ and define the three systems

$$
\begin{aligned}
G(s) & :=\Sigma\left(A_{c c}, I, C_{1 c}, 0\right) \\
U(s) & :=\Sigma\left(A_{c c}, B, C_{1 c}, D\right) \\
U^{\perp}(s) & :=\Sigma\left(A_{c c},-P_{s}^{\dagger} C^{\prime} D^{\perp}, C_{1 c}, D^{\perp}\right)
\end{aligned}
$$

where $P_{s}^{\dagger}$ is the Moore Penrose pseudo-inverse of $P_{s}$, i.e. the matrix such that $P_{s} P_{s}^{\dagger} P_{s}=P_{s}$ and $P_{s}^{\dagger} P_{s} P_{s}^{\dagger}=P_{s}$, and $D^{\perp}$ is a matrix such that

$$
\left[\begin{array}{c}
D^{\prime} \\
D^{\perp^{\prime}}
\end{array}\right]\left[\begin{array}{ll}
D & D^{\perp^{\prime}}
\end{array}\right]=I
$$

Then the system with transfer function

$$
F(s):=\left[\begin{array}{cc}
U(s) & U^{\perp}(s)
\end{array}\right]
$$

is square, inner and

$$
H(s):=G^{\sim}(s) F(s)=\Sigma\left(A_{c c},\left[\begin{array}{cc}
B & -P_{s}^{\dagger} C^{\prime} D^{\perp}
\end{array}\right], P_{s},\left[\begin{array}{ll}
0 & 0
\end{array}\right]\right)
$$

Proof Preliminarily, observe that matrix $D^{\perp}$ with the requested properties actually exists since $D^{\prime} D=I$. Systems $U(s)$ and $U^{\perp}(s)$ have the same dynamical matrix and output transformation. Hence, a realization of $F(s)$ is simply

$$
F(s)=\left[\begin{array}{c|cc}
A_{c c} & B & -P_{s}^{\dagger} C^{\prime} D^{\perp} \\
\hline C_{1 c} & D & D^{\perp}
\end{array}\right]
$$

The system with transfer function $F(s)$ is square since | $D \quad D^{\perp} \mid$ is square. Now let $\eta$ be the state of the system with transfer function $F(s)$ and $\xi$ the state of the adjoint system, namely of the system with transfer function

$$
F^{\sim}(s)=\left[\begin{array}{c|c}
-A_{c c}^{\prime} & -C_{1 c}^{\prime} \\
\hline B^{\prime} & D^{\prime} \\
-D^{\perp^{\prime}} C P_{s}^{\dagger} & D^{\perp^{\prime}}
\end{array}\right]
$$

Consider the series connection $F^{\sim}(s) F(s)$. By exploiting the fact that $C_{1 c}^{\prime} C_{1 c}=$ $C^{\prime} C+P_{s} B B^{\prime} P_{s}=-A_{c c}^{\prime} P_{s}-P_{s} A_{c c}$ and the properties of matrices $C, D$ and $D^{\perp}$, it follows

$$
\begin{align*}
\left(\dot{\xi}-P_{s} \dot{\eta}\right) & =-A_{c c}^{\prime}\left(\xi-P_{s} \eta\right)-\left(I-P_{s} P_{s}^{\dagger}\right) C^{\prime} D^{\perp} \psi_{2}  \tag{C.22}\\
l_{1} & =B^{\prime}\left(\xi-P_{s} \eta\right)+\psi_{1}  \tag{C.23}\\
l_{2} & =-D^{\perp^{\prime}} C\left(P_{s}^{\dagger} \xi-\eta\right)+\psi_{2} \tag{C.24}
\end{align*}
$$

where $\psi_{1}$ and $\psi_{2}$ are the inputs of $F(s)$ whereas $l_{1}$ and $l_{2}$ are the outputs of $F^{\sim}(s)$. Lemma C. 4 entails that $\operatorname{Ker}\left[P_{s}\right] \subset \operatorname{Ker}[C]$ so that $C\left(I-P_{s}^{\dagger} P_{s}\right) \beta=0$, if $\beta \in \operatorname{Ker}\left[P_{s}\right]$. On the other hand, $\left(I-P_{s}^{\dagger} P_{s}\right) \beta=0$, if $\beta \in \operatorname{Ker}\left[P_{s}\right]^{\perp}$. In conclusion $C\left(I-P_{s}^{\dagger} P_{s}\right) \beta=$ $0, \forall \beta$. Hence, eqs. (C.22)-(C.24) can be rewritten as

$$
\begin{aligned}
\left(\dot{\xi}-P_{s} \dot{\eta}\right) & =-A_{c c}^{\prime}\left(\xi-P_{s} \eta\right) \\
l_{1} & =B^{\prime}\left(\xi-P_{s} \eta\right)+\psi_{1} \\
l_{2} & =-D^{\perp^{\prime}} C P_{s}^{\dagger}\left(\xi-P_{s} \eta\right)+\psi_{2}
\end{aligned}
$$

These equations show that the transfer function from $\left[\psi_{1}^{\prime} \psi_{2}^{\prime}\right]^{\prime}$ to $\left[l_{1}^{\prime} l_{2}^{\prime}\right]^{\prime}$, i.e. $F^{\sim}(s) F(s)$, is the identity. Hence, $F(s)$ is inner.

Let now $\vartheta$ be the state of a realization of $G^{\sim}(s)$. A realization of $G^{\sim}(s) F(s)$ is given by

$$
\begin{aligned}
& \dot{\vartheta}=-A_{c c}^{\prime} \vartheta-C_{1 c}^{\prime} \nu \\
& \dot{\eta}=A_{c c} \eta+B \psi_{1}-P_{s}^{\dagger} C^{\prime} D^{\perp} \psi_{2} \\
& \nu=C_{1 c} \eta+D \psi_{1}+D^{\perp} \psi_{2} \\
& \sigma=\vartheta
\end{aligned}
$$

where $\nu$ and $\sigma$ are the input and output of $G^{\sim}(s)$, respectively. It is easy to check that this realization is equivalent to

$$
\begin{aligned}
\left(\dot{\vartheta}-P_{s} \dot{\eta}\right) & =-A_{c c}^{\prime}\left(\vartheta-P_{s} \eta\right) \\
\dot{\eta} & =A_{c c} \eta+B \psi_{1}-P_{s}^{\dagger} C^{\prime} D^{\perp} \psi_{2} \\
\sigma & =\left(\vartheta-P_{s} \eta\right)+P_{s} \eta
\end{aligned}
$$

so that a possible realization of $G^{\sim}(s) F(s)$ is

$$
G^{\sim}(s) F(s)=\left[\begin{array}{c|cc}
A_{c c} & B & -P_{s}^{\dagger} C^{\prime} D^{\perp} \\
\hline P_{s} & 0 & 0
\end{array}\right]
$$

as claimed in the statement.
Lemma C. 6 Let $A$ and $C$ be matrices with dimensions $n \times n$ and $p \times n$ respectively and assume the pair $(-A, C)$ is detectable. Then the unique symmetric solution of

$$
\begin{aligned}
C P & =0 \\
P & \geq 0 \\
A P+P A^{\prime} & \leq 0
\end{aligned}
$$

is the trivial solution $P=0$.
Proof By contradiction, suppose $P=P^{\prime} \neq 0$ is a solution. With no loss of generality, one can assume that

$$
P=\left[\begin{array}{cc}
P_{1} & 0 \\
0 & 0
\end{array}\right], \quad A=\left[\begin{array}{ll}
A_{1} & A_{2} \\
A_{3} & A_{4}
\end{array}\right], \quad C=\left[\begin{array}{ll}
C_{1} & C_{2}
\end{array}\right]
$$

with $P_{1}>0$. Actually, matrix $P$ can be put in above form by means of a suitable orthogonal transformation. Then, $C P=0$ implies $C_{1}=0$ and $A P+P A^{\prime} \leq 0$ yields

$$
\left[\begin{array}{cc}
A_{1} P_{1}+P_{1} A_{1}^{\prime} & P_{1} A_{3}^{\prime} \\
A_{3} P_{1} & 0
\end{array}\right] \leq 0
$$

so that $A_{3}=0$ and $A_{1} P_{1}+P_{1} A_{1}^{\prime} \leq 0$. Finally, the detectability of the pair

$$
(-A, C)=\left(\left[\begin{array}{cc}
-A_{1} & -A_{2} \\
0 & -A_{4}
\end{array}\right],\left[\begin{array}{ll}
0 & C_{2}
\end{array}\right]\right)
$$

entails that matrix $-A_{1}$ is stable. In this case, all possible solutions of the inequality

$$
A_{1} P_{1}+P_{1} A_{1}^{\prime} \leq 0
$$

are such that $-P_{1} \geq 0$ proving thus that $P_{1}>0$ is an impossibility. The final conclusion is that the trivial solution $P=0$ is unique.

## Appendix D

## Structural Properties

This section is devoted to the properties of reachability, observability, stabilizability and detectability of a linear time invariant continuous time $n$-dimensional system $\Sigma(A, B, C, D)$. All the material presented here can be found in any text of system theory, see, e.g., the book by Kailath [29].

Controllability and reconstructability will not be treated in the present context since they enjoy the same properties as reachability and observability, respectively.

In the forthcoming lemmas, whose proofs are easily available in specialized texts, the basic characterizations of the above properties will be provided.

Preliminarily, define the matrices $P_{C}(s), P_{B}(s), K_{o}$ and $K_{r}$ as

$$
\begin{aligned}
P_{C}(s) & :=\left[\begin{array}{c}
s I-A \\
C
\end{array}\right], P_{B}(s):=\left[\begin{array}{lll}
s I-A & B
\end{array}\right] \\
K_{o} & :=\left[\begin{array}{llll}
C^{\prime} & A^{\prime} C^{\prime} & \cdots & \left(A^{\prime}\right)^{n-1} C^{\prime}
\end{array}\right] \\
K_{r} & :=\left[\begin{array}{llll}
B & A B & \cdots & A^{n-1} B
\end{array}\right]
\end{aligned}
$$

Lemma D. 1 (Observability) System $\Sigma$ or, equivalently, the pair $(A, C)$, is observable if and only if the following equivalent conditions hold:
a) PBH test :

$$
\operatorname{rank}\left[P_{C}(s)\right]=n, \quad \forall s
$$

The set of eigenvalues of the unobservable part of $(A, C)$ coincides with the set of values of $s$ in correspondence of which matrix $P_{C}(s)$ looses rank.
b) Kalman test :

$$
\operatorname{rank}\left[K_{o}\right]=n
$$

c) Wonham test : Given an arbitrary symmetric set $\Lambda$ of $n$ complex number, there exists a matrix $L$ such that the spectrum of $A+L C$ coincides with $\Lambda$

Lemma D. 2 (Detectability) System $\Sigma$ or, equivalently, the pair $(A, C)$ is detectable if and only if the following equivalent conditions hold:
a) PBH test :

$$
\operatorname{rank}\left[P_{C}(s)\right]=n, \quad \operatorname{Re}(s) \geq 0
$$

b) Kalman test : The unobservable part of the system is stable.
c) Wonham test : There exists a matrix $L$ such that $A+L C$ is stable.

Lemma D. 3 (Reachability) System $\Sigma$ or, equivalently, the pair $(A, B)$ is reachable if and only if the following equivalent conditions hold:
a) PBH test :

$$
\operatorname{rank}\left[P_{B}(s)\right]=n, \quad \forall s
$$

The set of eigenvalues of the unreachable part of $(A, B)$ coincides with the set of values of $s$ in correspondence of which matrix $P_{B}(s)$ looses rank.
b) Kalman test :

$$
\operatorname{rank}\left[K_{r}\right]=n
$$

c) Wonham test : Given an arbitrary symmetric set $\Lambda$ of $n$ complex number, there exists a matrix $K$ such that the spectrum of $A+B K$ coincides with $\Lambda$

Lemma D. 4 (Stabilizability) System $\Sigma$ or, equivalently, the pair $(A, B)$ is stabilizable if and only if the following equivalent conditions hold :
a) PBH test:

$$
\operatorname{rank}\left[P_{B}(s)\right]=n, \quad \operatorname{Re}(s) \geq 0
$$

b) Kalman test : The unreachable part of the system is stable.
c) Wonham test : There exists a matrix $K$ such that $A+B K$ is stable.

At the light of what said in Section 2.5 for the zeros of $\Sigma$, one can now stress that $\Sigma$ is reachable (resp. observable) if and only if it does not possess input (resp. output) decoupling zeros, and, analogously, system $\Sigma$ is stabilizable (resp. detectable) if and only if it does not possess input (resp. output) decoupling zeros in the closed right half plane.

## Appendix E

## The Standard 2-Block Scheme

The material presented in this appendix is largely taken from the paper by Doyle et al. [17]. However, the proof of Lemma E. 3 is partially original.

Lemma E. 1 Consider the systems $P(s)$ and $\hat{P}(s):=P^{\prime}(s)$, where

$$
P(s):=\left[\begin{array}{c|cc}
A & B_{1} & B_{2} \\
\hline C_{1} & D_{11} & D_{12} \\
C_{2} & D_{21} & D_{22}
\end{array}\right]
$$

These two systems are connected in a feedback configuration to $K(s)$ and $\hat{K}(s)$, respectively, according to the block schemes of fig. E.1. Correspondingly, let $T(a, b ; s)$ denote the transfer function from the generic input $b$ to the generic output $a$. Then, given $K(s)$ and letting $\hat{K}(s):=K^{\prime}(s)$ it follows $T(\hat{z}, \hat{w} ; s)=T^{\prime}(z, w ; s)$. Conversely, given $\hat{K}(s)$ and letting $K(s):=\hat{K}^{\prime}(s)$, it follows $T(z, w ; s)=T^{\prime}(\hat{z}, \hat{w} ; s)$. Moreover, the eigenvalues of the two systems in fig. E. 1 coincide.

Proof Let $P_{i j}(s), i=j=1,2$, be the four transfer functions of $P(s)$. It results $T(\hat{z}, \hat{w} ; s)=P_{11}^{\prime}(s)+P_{21}^{\prime}(s) \hat{K}(s)\left[I-P_{22}^{\prime}(s) \hat{K}(s)\right]^{-1} P_{12}^{\prime}(s)$. If $\hat{K}(s)=K^{\prime}(s)$ it follows that $T(\hat{z}, \hat{w} ; s)=T^{\prime}(z, w ; s)$. Conversely, it results $T(z, w ; s)=P_{11}(s)+$ $P_{12}(s) K(s)\left[I-P_{22}(s) K(s)\right]^{-1} P_{21}(s)$. If $K(s)=\hat{K}^{\prime}(s)$ it follows that $T(z, w ; s)=$ $T^{\prime}(\hat{z}, \hat{w} ; s)$.


Figure E.1: A feedback system and its transpose


Figure E.2: A feedback system with full information

Now, let

$$
K(s)=\left[\begin{array}{c|c}
F & G \\
\hline H & E
\end{array}\right]
$$

be a controller such that the system of fig. E.1(a) is well defined, i.e. $\operatorname{det}\left[I-E D_{22}\right]$ is nonzero. It is easy to check that the dynamical matrix of the system of fig. E.1(a) is

$$
A_{a}=\left[\begin{array}{cc}
A+B_{2}\left(I-E D_{22}\right)^{-1} E C_{2} & B_{2}\left(I-E D_{22}\right)^{-1} H \\
G\left[I+D_{22}\left(I-E D_{22}\right)^{-1} E\right] C_{2} & F+G D_{22}\left(I-E D_{22}\right)^{-1} H
\end{array}\right]
$$

whereas that of system of fig. E.1(b) where $\hat{K}(s)=K^{\prime}(s)$ is

$$
A_{b}=\left[\begin{array}{cc}
A^{\prime}+C_{2}^{\prime}\left(I-E^{\prime} D_{22}^{\prime}\right)^{-1} E^{\prime} B_{2}^{\prime} & C_{2}^{\prime}\left(I-E^{\prime} D_{22}^{\prime}\right)^{-1} G^{\prime} \\
H^{\prime}\left[I+D_{22}^{\prime}\left(I-E^{\prime} D_{22}^{\prime}\right)^{-1} E^{\prime}\right] B_{2}^{\prime} & F^{\prime}+H^{\prime} D_{22}^{\prime}\left(I-E^{\prime} D_{22}^{\prime}\right)^{-1} G^{\prime}
\end{array}\right]
$$

Recalling Lemma B.9, it is easy to verify that $A_{b}^{\prime}=A_{a}$ so that the two systems have the same eigenvalues.

Lemma E. 2 Consider the system

$$
\bar{P}(s):=\left[\begin{array}{c|cc}
A & B_{1} & B_{2} \\
\hline C_{1} & 0 & D_{12} \\
I & 0 & 0 \\
0 & I & 0
\end{array}\right]
$$

connected in feedback configuration with $\bar{K}(s)$ according to the block-scheme of fig. E.2. Moreover, consider system

$$
\hat{P}(s):=\left[\begin{array}{c|cc}
A & B_{1} & B_{2} \\
\hline C_{1} & 0 & D_{12} \\
C_{2} & I & 0
\end{array}\right]
$$

connected in feedback configuration with $\hat{K}(s)$ as shown in fig. E.3. Finally suppose that matrix $A-B_{1} C_{2}$ is stable. Then,
i) Given $\hat{K}(s)$ and letting $\bar{K}(s):=\hat{K}(s)\left[C_{2} I\right]$ it follows
i1) $\bar{K}(s)$ stabilizes $\bar{P}(s)$ if and only if $\hat{K}(s)$ stabilizes $\hat{P}(s)$
i2) $T(\hat{z}, \hat{w} ; s)=T(\bar{z}, \bar{w} ; s)$


Figure E.3: A feedback system with partial information
ii) Given $\bar{K}(s)$ and defining $\hat{K}(s)$ through the block-scheme of fig. E. 4 where

$$
K_{v}(s):=\left[\begin{array}{c|cc}
A-B_{1} C_{2} & B_{1} & B_{2} \\
\hline 0 & 0 & I \\
I & 0 & 0 \\
-C_{2} & I & 0
\end{array}\right]
$$

it follows
ii1) $\hat{K}(s)$ stabilizes $\hat{P}(s)$ if and only if $\bar{K}(s)$ stabilizes $\bar{P}(s)$
ii2) $T(\hat{z}, \hat{w} ; s)=T(z, w ; s)$

Proof Point i) Let

$$
\hat{K}(s):=\left[\begin{array}{c|c}
\hat{A} & \hat{B} \\
\hline \hat{C} & \hat{D}
\end{array}\right]
$$

Then, by keeping in mind fig. E. 3 and the definition of $\hat{P}(s)$ it follows

$$
T(\hat{z}, \hat{w} ; s)=\left[\begin{array}{cc|c}
A+B_{2} \hat{D} C_{2} & B_{2} \hat{C} & B_{1}+B_{2} \hat{D} \\
\hat{B} C_{2} & \hat{A} & \hat{B} \\
\hline C_{1}+D_{12} \hat{D} C_{2} & D_{12} \hat{C} & D_{12} \hat{D}
\end{array}\right]
$$

while $\bar{K}(s)=\hat{K}(s)\left[C_{2} I\right]$ implies

$$
\bar{K}(s)=\left[\begin{array}{c|cc}
\hat{A} & \hat{B} C_{2} & \hat{B} \\
\hline \hat{C} & \hat{D} C_{2} & \hat{D}
\end{array}\right]
$$

It is easy to see (for instance by inspecting the state equations of the relevant systems with fig. E. 2 and E. 3 in mind) that the state representations of $T(\hat{z}, \hat{w} ; s)$ and $T(\bar{z}, \bar{w} ; s)$ are equal, so that point i) is proved.

Point ii) A state space representation of $T(\hat{z}, \hat{w} ; s)$ in fig. E. 3 is sought, when $\hat{K}(s)$ is built according to fig. E.4. Hence, let $\hat{x}, \bar{\mu}$ and $x_{v}$ be the state variables of


Figure E.4: A controller structure for the partial information case


Figure E.5: A feedback system
$\hat{P}(s), \bar{K}(s)$ and $K_{v}(s)$, respectively. A state description of $T(\hat{z}, \hat{w} ; s)$ is then given by

$$
\begin{aligned}
& \hat{P}(s):=\left\{\begin{array}{l}
\dot{\hat{x}}=A \hat{x}+B_{1} \hat{w}+B_{2} \hat{u} \\
\hat{z}=C_{1} \hat{x}+D_{12} \hat{u} \\
\hat{y}=C_{2} \hat{x}+\hat{w}
\end{array}\right. \\
& K_{v}(s):= \begin{cases}\dot{x}_{v}= & \left(A-B_{1} C_{2}\right) x_{v}+B_{1} \hat{y}+B_{2} \hat{u} \\
\eta_{v}= & -C_{2} x_{v}+\hat{y}\end{cases} \\
& \bar{K}(s):=\left\{\begin{array}{l}
\dot{\bar{\mu}}=\bar{A} \mu+\bar{B}_{1} x_{v}+\bar{B}_{2} \eta_{v} \\
\hat{u}=\bar{C} \mu+\bar{D}_{1} x_{v}+\bar{D}_{2} \eta_{v}
\end{array}\right.
\end{aligned}
$$

Letting $\varepsilon:=x_{v}-\hat{x}$ it results $\dot{\varepsilon}=\left(A-B_{1} C_{2}\right) \varepsilon$, which is by assumption a stable system. Then, in order to evaluate the stability of system of fig. E. 3 and to compute the transfer function $T(\hat{z}, \hat{w} ; s)$ one can, without loss of generality, put $\varepsilon=0$, i.e. $x_{v}=\hat{x}$ in the relevant equations. It follows

$$
T(\hat{z}, \hat{w} ; s)=\left[\begin{array}{cc|c}
A+B_{2} \bar{D}_{1} & B_{2} \bar{C} & B_{1}+B_{2} \bar{D}_{2} \\
\bar{B}_{1} & \bar{A} & \bar{B}_{2} \\
\hline C_{1}+D_{12} \bar{D}_{1} & D_{12} \bar{C} & D_{12} \bar{D}_{2}
\end{array}\right]
$$

On the other hand, it is simple to verify that the feedback system of fig. E. 2 is described by the very same equations. Hence Point ii) immediately follows.

Lemma E. 3 Let $\gamma$ be a positive scalar and assume that there exists the symmetric, positive semidefinite solution $P_{\infty}$ of the Riccati equation (in the unknown $P$ )

$$
\begin{equation*}
0=P A_{c}+A_{c}^{\prime} P+P\left(\gamma^{-2} B_{1} B_{1}^{\prime}-B_{2} B_{2}^{\prime}\right) P+C_{1 c}^{\prime} C_{1 c} \tag{E.1}
\end{equation*}
$$

that is such that

$$
\begin{equation*}
A_{c c}:=A_{c}+\left(\gamma^{-2} B_{1} B_{1}^{\prime}-B_{2} B_{2}^{\prime}\right) P_{\infty} \tag{E.2}
\end{equation*}
$$

is stable. Further, consider the block-scheme of fig. E. 5 where

$$
\begin{gathered}
P_{q}(s):=\left[\begin{array}{c|cc}
A_{c}-B_{2} B_{2}^{\prime} P_{\infty} & B_{1} & B_{2} \\
\hline C_{1 c}-D_{12} B_{2}^{\prime} P_{\infty} & 0 & D_{12} \\
-\gamma^{-2} B_{1}^{\prime} P_{\infty} & I & 0
\end{array}\right] \\
Q(s):=\left[\begin{array}{c|c}
A_{q} & B_{q} \\
\hline C_{q} & D_{q}
\end{array}\right]
\end{gathered}
$$

with

$$
\begin{align*}
D_{12}^{\prime} C_{1 c} & =0  \tag{E.3}\\
D_{12}^{\prime} D_{12} & =I \tag{E.4}
\end{align*}
$$

and $T(z, w ; s):=\Sigma\left(A_{z}, B_{z}, C_{z}, D_{z}\right)$ is the transfer function from $w$ to $z$. Then, $A_{z}$ is stable and $\|T(z, w ; s)\|_{\infty}<\gamma$ if and only if $A_{q}$ is stable and $\|Q(s)\|_{\infty}<\gamma$.

Proof Sufficiency Let $A_{q}$ be stable and $\|Q(s)\|_{\infty}<\gamma$. The system in fig. E. 5 with transfer function $T(z, w ; s)$ is described by

$$
T(z, w ; s)=\left[\begin{array}{c|c}
A_{z} & B_{z} \\
\hline C_{z} & D_{z}
\end{array}\right]
$$

where

$$
\begin{align*}
A_{z} & :=\left[\begin{array}{cc}
A_{c}-B_{2}\left(B_{2}^{\prime}+\gamma^{-2} D_{q} B_{1}^{\prime}\right) P_{\infty} & B_{2} C_{q} \\
-\gamma^{-2} B_{q} B_{1}^{\prime} P_{\infty} & A_{q}
\end{array}\right]  \tag{E.5}\\
B_{z} & :=\left[\begin{array}{cc}
B_{1}+B_{2} D_{q} \\
B_{q}
\end{array}\right]  \tag{E.6}\\
C_{z} & :=\left[\begin{array}{ll}
C_{1 c}-D_{12}\left(B_{2}^{\prime}+\gamma^{-2} D_{q} B_{1}^{\prime}\right) P_{\infty} & D_{12} C_{q}
\end{array}\right]  \tag{E.7}\\
D_{z} & :=D_{12} D_{q} \tag{E.8}
\end{align*}
$$

Since $\|Q(s)\|_{\infty}<\gamma$, it results $\bar{\sigma}\left(D_{q}\right)<\gamma$. Then, by recalling that $A_{q}$ is stable, Theorem 2.13, applied to system $\Sigma\left(A_{q}, B_{q}, C_{q}, D_{q}\right)$ entails that there exists the symmetric, positive semidefinite solution $S_{\infty}$ of the Riccati equation (in the unknown $S$ )

$$
\begin{equation*}
0=S \hat{A}_{q}+\hat{A}_{q}^{\prime} S+S B_{q} \Delta_{1 q} B_{q}^{\prime} S+C_{q}^{\prime} \Delta_{2 q} C_{q} \tag{E.9}
\end{equation*}
$$

where

$$
\begin{equation*}
\Delta_{1 q}:=\left(\gamma^{2} I-D_{q}^{\prime} D_{q}\right)^{-1}, \quad \Delta_{2 q}:=\left(I-\gamma^{-2} D_{q} D_{q}^{\prime}\right)^{-1} \tag{E.10}
\end{equation*}
$$

and

$$
\begin{equation*}
\hat{A}_{q}:=A_{q}+B_{q} \Delta_{1 q} D_{q}^{\prime} C_{q} \tag{E.11}
\end{equation*}
$$

Therefore, matrix

$$
\begin{equation*}
\tilde{A}_{q}:=A_{q}+B_{q} \Delta_{1 q}\left(D_{q}^{\prime} C_{q}+B_{q}^{\prime} S_{\infty}\right) \tag{E.12}
\end{equation*}
$$

is stable. From the definition of $D_{z}$ (eq. (E.8)) and eq. (E.4) it follows that

$$
\begin{equation*}
\bar{\sigma}\left(D_{z}\right)<\gamma \tag{E.13}
\end{equation*}
$$

which is a necessary condition for $\|T(z, w ; s)\|_{\infty}<\gamma$. It is now shown that

$$
V_{\infty}:=\left[\begin{array}{cc}
P_{\infty} & 0 \\
0 & S_{\infty}
\end{array}\right]
$$

is a symmetric and positive semidefinite solution of the Riccati equation (in the unknown $V$ )

$$
\begin{equation*}
0=V \hat{A}_{z}+\hat{A}_{z}^{\prime} V+V B_{z} \Delta_{1 z} B_{z}^{\prime} V+C_{z}^{\prime} \Delta_{2 z} C_{z} \tag{E.14}
\end{equation*}
$$

where

$$
\begin{equation*}
\Delta_{1 z}:=\left(\gamma^{2} I-D_{z}^{\prime} D_{z}\right)^{-1}, \quad \Delta_{2 z}:=\left(I-\gamma^{-2} D_{z} D_{z}^{\prime}\right)^{-1} \tag{E.15}
\end{equation*}
$$

and

$$
\begin{equation*}
\hat{A}_{z}:=A_{z}+B_{z} \Delta_{1 z} D_{z}^{\prime} C_{z} \tag{E.16}
\end{equation*}
$$

In fact, preliminarily observe that the right hand side of the Riccati equation (E.14) can be rewritten (if eqs.(E.15),(E.16) and Lemma B. 9 are taken into account) as

$$
\begin{aligned}
& V A_{z}+A_{z}^{\prime} V+\left(V B_{z}+C_{z}^{\prime} D_{z}\right) \Delta_{1 z}\left(V B_{z}+C_{z}^{\prime} D_{z}\right)^{\prime}+C_{z}^{\prime} C_{z}:= \\
& \quad:=f(V):=\left[\begin{array}{ll}
f_{11}(V) & f_{12}(V) \\
f_{21}(V) & f_{22}(V)
\end{array}\right]
\end{aligned}
$$

In view of eqs. (E.5)-(E.8) it follows

$$
V_{\infty} A_{z}=\left[\begin{array}{cc}
P_{\infty} A_{c}-P_{\infty} B_{2} B_{2}^{\prime} P_{\infty}-\gamma^{-2} P_{\infty} B_{2} D_{q} B_{1}^{\prime} P_{\infty} & P_{\infty} B_{2} C_{q} \\
-\gamma^{-2} S_{\infty} B_{q} B_{1}^{\prime} P_{\infty} & S_{\infty} A_{q}
\end{array}\right]
$$

and, by recalling eqs. (E.3),(E.4),

$$
\begin{align*}
V_{\infty} B_{z}+C_{z}^{\prime} D_{z} & =\left[\begin{array}{c}
\gamma^{-2} P_{\infty} B_{1} \Delta_{1 q}^{-1} \\
S_{\infty} B_{q}+C_{q}^{\prime} D_{q}
\end{array}\right]  \tag{E.17}\\
C_{z}^{\prime} C_{z} & =\left[\begin{array}{ll}
C_{z 1} & C_{z 2} \\
C_{z 2}^{\prime} & C_{z 3}
\end{array}\right]
\end{align*}
$$

where

$$
\begin{aligned}
C_{z 1}:= & C_{1 c}^{\prime} C_{1 c}+P_{\infty} B_{2} B_{2}^{\prime} P_{\infty}+\gamma^{-4} P_{\infty} B_{1} D_{q}^{\prime} D_{q} B_{1}^{\prime} P_{\infty}+ \\
& +\gamma^{-2} P_{\infty} B_{2} D_{q} B_{1}^{\prime} P_{\infty}+\gamma^{-2} P_{\infty} B_{1} D_{q}^{\prime} B_{2}^{\prime} P_{\infty} \\
C_{z 2}:= & -P_{\infty} B_{2} C_{q}-\gamma^{-2} P_{\infty} B_{1} D_{q}^{\prime} C_{q} \\
C_{z 3}:= & C_{q}^{\prime} C_{q}
\end{aligned}
$$

Thanks to eqs. (E.4),(E.8),(E.10) and (E.15) it is

$$
\begin{equation*}
\Delta_{1 z}=\Delta_{1 q} \tag{E.18}
\end{equation*}
$$

so that

$$
\begin{aligned}
f_{11}\left(V_{\infty}\right)= & P_{\infty} A_{c}+A_{c}^{\prime} P_{\infty}-P_{\infty}\left(B_{2} B_{2}^{\prime}-\gamma^{-2} B_{1} B_{1}^{\prime}\right) P_{\infty}+ \\
& +C_{1 c}^{\prime} C_{1 c}=0
\end{aligned}
$$

since matrix $P_{\infty}$ is a solution of eq. (E.1). By direct substitution one finds $f_{12}\left(V_{\infty}\right)=$ 0 . Finally,

$$
\begin{aligned}
f_{22}\left(V_{\infty}\right)= & S_{\infty} A_{q}+A_{q}^{\prime} S_{\infty}+\left(S_{\infty} B_{q}+C_{q}^{\prime} D_{q}\right) \Delta_{1 q} \\
& \cdot\left(S_{\infty} B_{q}+C_{q}^{\prime} D_{q}\right)^{\prime}+C_{q}^{\prime} C_{q}
\end{aligned}
$$

By recalling Lemma B. 9 and eqs. (E.10),(E.11) it is easy to check that $f_{22}\left(V_{\infty}\right)$ coincides with the right hand side of eq. (E.9) and is therefore zero.

Now, let

$$
\begin{aligned}
\tilde{A}_{z} & :=A_{z}+B_{z} \Delta_{1 z}\left(D_{z}^{\prime} C_{z}+B_{z}^{\prime} V_{\infty}\right) \\
& =\left[\begin{array}{cc}
A_{c}-B_{2} B_{2}^{\prime} P_{\infty}+\gamma^{-2} B_{1} B_{1}^{\prime} P_{\infty} & A_{q}+B_{q} \Delta_{1 z}\left(B_{q}^{\prime} S_{\infty}+D_{q}^{\prime} C_{q}\right)
\end{array}\right]
\end{aligned}
$$

where $\Phi$ denotes a matrix of no interest in this context. In view of eqs. (E.2),(E.12) and (E.18), the conclusion is drawn that the eigenvalues of $\tilde{A}_{z}$ are those of $A_{c c}$ and $\tilde{A}_{q}$. The solution $V_{\infty}$ of eq. (E.14) is therefore the stabilizing one since matrices $A_{c c}$ and $\tilde{A}_{q}$ are stable. Thus, by recalling eq. (E.13), it can be said that, thanks to Theorem 2.14, $\|T(z, w ; s)\|_{\infty}<\gamma$ and $A_{z}$ is stable.

Necessity Assume that $\|T(z, w ; s)\|_{\infty}<\gamma$ and $A_{z}$ is stable. It has to be proved that $A_{q}$ is stable and $\|Q(s)\|_{\infty}<\gamma$. To this aim, consider the Hamiltonian matrix

$$
Z_{z}:=\left[\begin{array}{cc}
\hat{A}_{z} & B_{z} \Delta_{1 z} B_{z}^{\prime} \\
-C_{z}^{\prime} \Delta_{2 z} C_{z} & -\hat{A}_{z}^{\prime}
\end{array}\right]
$$

and the relevant Riccati equation (in the unknown $V$ )

$$
\begin{equation*}
0=V \hat{A}_{z}+\hat{A}_{z}^{\prime} V+V B_{z} \Delta_{1 z} B_{z}^{\prime} V+C_{z}^{\prime} \Delta_{2 z} C_{z} \tag{E.19}
\end{equation*}
$$

In view of eqs. (E.15) and (E.16) matrix $Z_{z}$ is associated with system $T(z, w ; s)=$ $\Sigma\left(A_{z}, B_{z}, C_{z}, D_{z}\right)$ defined by eqs. (E.5)-(E.8). By taking into account the equations (E.3),(E.4),(E.1),(E.10) and (E.11), it will be now verified that, chosen

$$
\Psi:=\left[\begin{array}{rr}
I & 0 \\
-\Omega & I
\end{array}\right], \quad \Omega:=\left[\begin{array}{cc}
P_{\infty} & 0 \\
0 & 0
\end{array}\right]
$$

it is $\Psi Z_{z} \Psi^{-1}=\hat{Z}_{z}$, where

$$
\hat{Z}_{z}=\left[\begin{array}{cccc}
\Phi_{1} & \Phi_{2} & \Phi_{3} & \Phi_{4}  \tag{E.20}\\
0 & \hat{A}_{q} & \Phi_{5} & B_{q} \Delta_{1 q} B_{q}^{\prime} \\
0 & 0 & \Phi_{6} & 0 \\
0 & -C_{q}^{\prime} \Delta_{2 q} C_{q} & \Phi_{7} & -\hat{A}_{q}^{\prime}
\end{array}\right]
$$

The explicit expressions of the matrices $\Phi_{i}, i=1, \cdots, 7$ appearing in $\hat{Z}_{z}$ are not given since of no interest in the subsequent discussion. In order to verify eq. (E.20) observe that by recalling eqs. (E.15),(E.16) and Lemma B.9, it is

$$
\Psi Z_{z} \Psi^{-1}=\left[\begin{array}{cc}
\zeta_{11} & \zeta_{12} \\
\zeta_{21} & -\zeta_{11}^{\prime}
\end{array}\right]
$$

where

$$
\begin{aligned}
& \zeta_{11}:=A_{z}+B_{z} \Delta_{1 z}\left(D_{z}^{\prime} C_{z}+B_{z}^{\prime} \Omega\right), \quad \zeta_{12}:=B_{z} \Delta_{1 z} B_{z}^{\prime} \\
& \zeta_{21}:=-\Omega A_{z}-A_{z}^{\prime} \Omega-\left(\Omega B_{z}+C_{z}^{\prime} D_{z}\right) \Delta_{1 z}\left(\Omega B_{z}+C_{z}^{\prime} D_{z}\right)^{\prime}-C_{z}^{\prime} C_{z}
\end{aligned}
$$

Further, by taking into account eqs. (E.3)-(E.8) and (E.10), one has

$$
\Omega B_{z}+C_{z}^{\prime} D_{z}=\left[\begin{array}{c}
\gamma^{-2} P_{\infty} B_{1} \Delta_{1 q}^{-1} \\
C_{q}^{\prime} D_{q}
\end{array}\right]
$$

so that, thanks to eq. (E.18),

$$
\zeta_{11}=\left[\begin{array}{cc}
\Phi_{1} & \Phi_{2} \\
0 & A_{q}+B_{q} \Delta_{1 q} D_{q}^{\prime} C_{q}
\end{array}\right]
$$

In view of eqs. (E.11) and (E.20) it follows that $\zeta_{11}$ coincides with the submatrix of $\hat{Z}_{z}$ made up of the elements belonging to the first two blocks of rows and columns. In a similar way, thanks to eqs. (E.18) and (E.6),

$$
\zeta_{12}=\left[\begin{array}{cc}
\Phi_{3} & \Phi_{4} \\
\Phi_{5} & B_{q} \Delta_{1 q} B_{q}^{\prime}
\end{array}\right]
$$

which coincides with the submatrix of $\hat{Z}_{z}$ made up of the elements belonging to the first two blocks of rows and the last two blocks of columns. Finally, one obtains,

$$
\zeta_{21}=\left[\begin{array}{cc}
g_{1} & 0 \\
0 & g_{2}
\end{array}\right]
$$

with

$$
\begin{aligned}
& g_{1}:=-\left[P_{\infty} A_{c}+A_{c}^{\prime} P_{\infty}-P_{\infty}\left(B_{2} B_{2}^{\prime}-\gamma^{-2} B_{1} B_{1}^{\prime}\right) P_{\infty}+C_{1 c}^{\prime} C_{1 c}\right] \\
& g_{2}:=-C_{q}^{\prime} C_{q}-C_{q}^{\prime} D_{q} \Delta_{1 q} D_{q}^{\prime} C_{q}
\end{aligned}
$$

By recalling eq. (E.1) it results $g_{1}=0$, whereas from Lemma B. 9 it follows $g_{2}=$ $-C_{q}^{\prime} \Delta_{2 q} C_{q}$. Therefore, $\zeta_{21}$ coincides with the submatrix of $\hat{Z}_{z}$ made up of the elements belonging to the last two blocks of rows and the first two blocks of columns.

In conclusion eq. (E.20) has been proved. Now notice that the eigenvalues of $\hat{Z}_{z}$ are those of the matrices $\Phi_{1}, \Phi_{6}$ and

$$
Z_{q}:=\left[\begin{array}{cc}
\hat{A}_{q} & B_{q} \Delta_{1 q} B_{q}^{\prime} \\
-C_{q}^{\prime} \Delta_{2 q} C_{q} & -\hat{A}_{q}^{\prime}
\end{array}\right]
$$

Condition (E.13) is satisfied since $\|T(z, w ; s)\|_{\infty}<\gamma$ and matrix $A_{z}$ is stable so that from Theorem 2.13 it follows that the eigenvalues of $Z_{z}$ do not lie on the imaginary axis. Thus also the eigenvalues of $\hat{Z}_{z}$ do not have zero real part and the same can be said of the eigenvalues of $Z_{q}$, thanks to the above discussion.

Since the resulting system is stable ( $A_{z}$ is stable), the pair ( $A_{q}, B_{q}$ ) is stabilizable and such is the pair $\left(\hat{A}_{q},\left(B_{q} \Delta_{1 q} B_{q}^{\prime}\right)^{1 / 2}\right)$ too. In fact, first recall that in view of eq. (E.11) the stabilizability of the pair $\left(A_{q}, B_{q}\right)$ implies that also the pair $\left(\hat{A}_{q}, B_{q}\right)$ is such (state feedback does not affect this property). Second, notice that if $\left(B_{q} \Delta_{1 q} B_{q}^{\prime}\right)^{1 / 2} x=$ 0 , then also $B_{q} \Delta_{1 q} B_{q}^{\prime} x=0$ and $B_{q}^{\prime} x=0\left(\Delta_{1 q}\right.$ is positive definite), so that, in view of the $P B H$ test (Lemma D.4), should the pair $\left(\hat{A}_{q},\left(B_{q} \Delta_{1 q} B_{q}^{\prime}\right)^{1 / 2}\right)$ not be stabilizable then also the pair ( $\hat{A}_{q}, B_{q}$ ) would not be such.

Therefore Lemma C. 3 can be applied to system $\Sigma\left(A_{q}, B_{q}, C_{q}, D_{q}\right)$ yielding the existence of the symmetric, positive semidefinite and stabilizing solution $S_{\infty}$ of the Riccati equation (in the unknown $S$ )

$$
\begin{equation*}
0=S \hat{A}_{q}+\hat{A}_{q}^{\prime} S+S B_{q} \Delta_{1 q} B_{q}^{\prime} S+C_{q}^{\prime} \Delta_{2 q} C_{q} \tag{E.21}
\end{equation*}
$$



Figure E.6: A feedback system

The same lemma can be applied to eq. (E.19) since: $(i)$ the pair $\left(\hat{A}_{z}, B_{z}\right)$ is stabilizable (indeed, matrix $A_{z}$ is stable); (ii) the eigenvalues of matrix $Z_{z}$ do not have zero real part; (iii) matrix $C_{z}^{\prime} \Delta_{2 z} C_{z}$ is symmetric; (iv) matrix $B_{z} \Delta_{1 z} B_{z}^{\prime}$ is positive semidefinite $\left(\Delta_{1 z}>0\right.$ thanks to Lemma B. 11 and $\bar{\sigma}\left(D_{z}\right)<\gamma\left(\|T(z, w ; s)\|_{\infty}<\gamma\right)$. Therefore, the stabilizing solution of eq. (E.19), which is unique in view of Lemma C.3, is given (see the sufficiency part of the proof of the lemma) by

$$
V_{\infty}=\left[\begin{array}{cc}
P_{\infty} & 0 \\
0 & S_{\infty}
\end{array}\right]
$$

Matrix $V_{\infty}$ is positive semidefinite (see Theorem 2.13), hence also $S_{\infty}$ is such.
Finally, Theorem 2.14 applied to system $Q(s)$ (notice that eqs. (E.4) and (E.8) imply $\bar{\sigma}\left(D_{q}\right)=\bar{\sigma}\left(D_{z}\right)$ ) and eq. (E.21) lead to the desired conclusion, namely $A_{q}$ stable and $\|Q(s)\|_{\infty}<\gamma$.

Lemma E. 4 Consider the system

$$
P(s):=\left[\begin{array}{c|cc}
A & B_{1} & B_{2} \\
\hline C_{1} & 0 & D_{12} \\
C_{2} & D_{21} & 0
\end{array}\right]
$$

with $D_{12}^{\prime} D_{12}=I$ and, for a given $\gamma>0$, assume that there exists the symmetric, positive semidefinite and stabilizing solution $P_{\infty}$ of the Riccati equation (in the unknown P)

$$
\begin{equation*}
0=P A_{c}+A_{c}^{\prime} P+P\left(\gamma^{-2} B_{1} B_{1}^{\prime}-B_{2} B_{2}^{\prime}\right) P+C_{1 c}^{\prime} C_{1 c} \tag{E.22}
\end{equation*}
$$

that is such that

$$
\begin{equation*}
A_{c c}:=A_{c}+\left(\gamma^{-2} B_{1} B_{1}^{\prime}-B_{2} B_{2}^{\prime}\right) P_{\infty} \tag{E.23}
\end{equation*}
$$

is stable, where

$$
\begin{equation*}
A_{c}:=A-B_{2} D_{12}^{\prime} C_{1}, C_{1 c}:=\left(I-D_{12} D_{12}^{\prime}\right) C_{1} \tag{E.24}
\end{equation*}
$$

Further, consider the block-scheme in fig. E. 6 and E.7, where

$$
P_{t}(s):=\left[\begin{array}{c|cc}
A+\gamma^{-2} B_{1} B_{1}^{\prime} P_{\infty} & B_{1} & B_{2} \\
\hline B_{2}^{\prime} P_{\infty}+D_{12}^{\prime} C_{1} & 0 & I \\
C_{2}+\gamma^{-2} D_{21} B_{1}^{\prime} P_{\infty} & D_{21} & 0
\end{array}\right]
$$

and let $T(z, w ; s)$ and $T(q, r ; s)$ be the transfer functions from $w$ to $z$ and from $r$ to $q$, respectively. Then, a controller $K(s)$ is admissible in $R H_{\infty}$ for $P(s)$ and such that $\|T(z, w ; s)\|_{\infty}<\gamma$ if and only if it is admissible in $R H_{\infty}$ for $P_{t}(s)$ and such that $\|T(q, r ; s)\|_{\infty}<\gamma$.


Figure E.7: An auxiliary feedback system


Figure E.8: A second auxiliary feedback system

Proof Consider the block-scheme in fig. E. 8 where

$$
P_{q}(s):=\left[\begin{array}{c|cc}
A_{c}-B_{2} B_{2}^{\prime} P_{\infty} & B_{1} & B_{2} \\
\hline C_{1 c}-D_{12} B_{2}^{\prime} P_{\infty} & 0 & D_{12} \\
-\gamma^{-2} B_{1}^{\prime} P_{\infty} & I & 0
\end{array}\right]
$$

and let $x_{q}$ and $x_{t}$ be the state variables of $P_{q}(s)$ and $P_{t}(s)$, respectively. It is easy to check that, by letting $\varepsilon:=x_{q}-x_{t}$, it results

$$
\begin{aligned}
\dot{\varepsilon} & =A_{c c} \varepsilon \\
\dot{x}_{q} & =A x_{q}+B_{1} w_{q}+B_{2} u-B_{2}\left(B_{2}^{\prime} P_{\infty}+D_{12}^{\prime} C_{1}\right) \varepsilon \\
z_{q} & =C_{1} x_{q}+D_{12} u-D_{12}\left(B_{2}^{\prime} P_{\infty}+D_{12}^{\prime} C_{1}\right) \varepsilon \\
y & =C_{2} x_{q}+D_{21} w_{q}-\left(C_{2}+\gamma^{-2} D_{21} B_{1}^{\prime} P_{\infty}\right) \varepsilon
\end{aligned}
$$

Being matrix $A_{c c}$ stable, the system in fig. E. 8 is internally stable if and only if the system in fig. E. 6 is such, since the equations for system $P(s)$ coincide with the last ones with $\varepsilon=0$. For the same reason $T(z, w ; s)=T\left(z_{q}, w_{q} ; s\right)$, where $T\left(z_{q}, w_{q} ; s\right)$ is the transfer function from $w_{q}$ to $z_{q}$ of the system in fig. E.8.

Thus the lemma is proved in view of Lemma E.3. Indeed, the system in fig. E. 8 is equal to the one depicted in fig. E. 5 if in the former figure the system with input $r$ and output $q$ is denoted by $Q(s)$. Further, the assumptions required by Lemma E. 3 are satisfied.

## Appendix F

## Loop Shifting

This appendix presents some results which are (mainly) exploited in Remark 5.24. In particular, Lemmas F. 1 and F. 2 are almost standard in matrix algebra, while the remaining material has been taken from the paper by Parrott [46].

Lemma F. 1 Let $A$ be a $m \times n$ dimensional matrix with $\|A\| \leq 1$. Then

$$
A^{\sim}\left(I-A A^{\sim}\right)^{1 / 2}=\left(I-A^{\sim} A\right)^{1 / 2} A^{\sim}
$$

Proof Preliminarily observe that both the matrices $I-A A^{\sim}$ and $I-A^{\sim} A$ have real square roots since they are positive semidefinite thanks to the assumption on $\|A\|$ which implies the nonnegativity of their eigenvalues. In fact, if, for instance, $\left(I-A A^{\sim}\right) x=\lambda x$ with $\lambda<0$ and $x \neq 0$, then $x^{\sim} A A^{\sim} x=x^{\sim} x(1-\lambda)>x^{\sim} x$ which implies $\|A\|>1$.

The proof of the lemma is carried out by assuming $m \geq n$, the reverse case being dealt with by modifying the subsequent discussion in an obvious way. By recalling Remark 2.11 it is possible to write $A=U_{1} \Delta V^{\sim}$ with $U_{1}=U[I 0]^{\prime}, U^{\sim} U=U U^{\sim}=I$, $U_{1}^{\sim} U_{1}=I, V^{\sim} V=V V^{\sim}=I$ and $\Delta$ diagonal and real. Therefore,

$$
I-A^{\sim} A=I-V \Delta U_{1}^{\sim} U_{1} \Delta V^{\sim}=I-V \Delta^{2} V^{\sim}=V\left(I-\Delta^{2}\right) V^{\sim}=V D^{2} V^{\sim}
$$

where $D^{2}:=I-\Delta^{2}$ is diagonal. Then

$$
\left(I-A^{\sim} A\right)^{1 / 2} A^{\sim}=V D V^{\sim} V \Delta U_{1}^{\sim}=V D \Delta U_{1}^{\sim}=V \Delta D U_{1}^{\sim}
$$

having exploited the fact that $V D^{2} V^{\sim}=V D V^{\sim} V D V^{\sim}$ (first equality sign) and the fact that the product of two diagonal matrices commutes (last equality sign). Being $U_{1}^{\sim}=\left[\begin{array}{ll}I & 0\end{array}\right] U^{\sim}$ it follows

$$
\left(I-A^{\sim} A\right)^{1 / 2} A^{\sim}=V \Delta\left[\begin{array}{ll}
D & 0 \tag{F.1}
\end{array}\right] U^{\sim}
$$

On the other hand, by again exploiting the relation between $U$ and $U_{1}$, it is

$$
\begin{aligned}
I-A A^{\sim} & =I-U_{1} \Delta V^{\sim} V \Delta U_{1}^{\sim} \\
& =I-U_{1} \Delta^{2} U_{1}^{\sim} \\
& =I-U\left[\begin{array}{l}
I \\
0
\end{array}\right] \Delta^{2}\left[\begin{array}{ll}
I & 0
\end{array}\right] U^{\sim}
\end{aligned}
$$

$$
\begin{aligned}
& =I-U\left[\begin{array}{cc}
\Delta^{2} & 0 \\
0 & 0
\end{array}\right] U^{\sim} \\
& =U\left[I-\left[\begin{array}{cc}
\Delta^{2} & 0 \\
0 & 0
\end{array}\right]\right] U^{\sim} \\
& =U\left[\begin{array}{cc}
D^{2} & 0 \\
0 & I
\end{array}\right] U^{\sim}
\end{aligned}
$$

so that

$$
\begin{align*}
A^{\sim}\left(I-A A^{\sim}\right)^{1 / 2} & =V \Delta U_{1}^{\sim} U\left[\begin{array}{cc}
D & 0 \\
0 & I
\end{array}\right] U^{\sim} \\
& =V \Delta\left[\begin{array}{ll}
I & 0
\end{array}\right] U^{\sim} U\left[\begin{array}{cc}
D & 0 \\
0 & I
\end{array}\right] U^{\sim} \\
& =V \Delta\left[\begin{array}{ll}
I & 0
\end{array}\right]\left[\begin{array}{cc}
D & 0 \\
0 & I
\end{array}\right] U^{\sim} \\
& =V \Delta\left[\begin{array}{ll}
D & 0
\end{array}\right] U^{\sim} \tag{F.2}
\end{align*}
$$

From eqs. (F.1),(F.2) the lemma follows.
Lemma F. 2 Let $A$ be a $m \times n$ dimensional matrix with $\operatorname{rank}[A]=r$. Further, let $A_{R}:=\left(A^{\sim} A\right)^{1 / 2}$ and $A_{L}:=\left(A A^{\sim}\right)^{1 / 2}$. Then there exists a $m \times n$ dimensional matrix $T$ such that
a) $A=T A_{R}=A_{L} T$.
b) $A_{R} T^{\sim} T=A_{R}$.
c) $T^{\sim} T \leq I$.

Proof Consider the singular value decomposition (see Section 2.6) of $A$, namely $A=U S V^{\sim}$ with

$$
U:=\left[\begin{array}{ll}
U_{1} & U_{2}
\end{array}\right], \quad S:=\left[\begin{array}{cc}
D & 0 \\
0 & 0
\end{array}\right]
$$

where $U_{1}$ is $m \times r$ dimensional and $D$ is $r \times r$ dimensional. Define

$$
S_{R}:=\left[\begin{array}{cc}
D & 0 \\
0 & 0
\end{array}\right], S_{L}:=\left[\begin{array}{cc}
D & 0 \\
0 & 0
\end{array}\right]
$$

where $S_{R}$ and $S_{L}$ are square and $n \times n$ and $m \times m$ dimensional, respectively, so that $S_{R}^{2}=S^{\prime} S$ and $S_{L}^{2}=S S^{\prime}$. Then it is now proved that

$$
A_{R}=V S_{R} V^{\sim}, A_{L}=U S_{L} U^{\sim}, T=\left[\begin{array}{ll}
U_{1} & 0 \tag{F.3}
\end{array}\right] V^{\sim}
$$

verify the theorem. First observe that

$$
\begin{aligned}
A_{R}^{\sim} A_{R} & =V S_{R} V^{\sim} V S_{R} V^{\sim} \\
& =V S_{R}^{2} V^{\sim} \\
& =V S^{\prime} S V^{\sim} \\
& =V S^{\prime} U^{\sim} U S V^{\sim}=A^{\sim} A
\end{aligned}
$$

Thus $A_{R}$ is a square root of $A^{\sim} A$. Second, notice that

$$
\begin{aligned}
A_{L}^{\sim} A_{L} & =U S_{L} U^{\sim} U S_{L} U^{\sim} \\
& =U S_{L}^{2} U^{\sim} \\
& =U S S^{\prime} U^{\sim} \\
& =U S V^{\sim} V S^{\prime} U^{\sim}=A A^{\sim}
\end{aligned}
$$

Thus, $A_{L}$ is a square root of $A A^{\sim}$.
Point a) From eq. (F.3) it follows that

$$
\begin{aligned}
T A_{R} & =\left[\begin{array}{ll}
U_{1} & 0
\end{array}\right] V^{\sim} V S_{R} V^{\sim} \\
& =\left[\begin{array}{ll}
U_{1} & 0
\end{array}\right] S_{R} V^{\sim} \\
& =\left[\begin{array}{ll}
U_{1} & U_{2}
\end{array}\right] S V^{\sim}=A
\end{aligned}
$$

and

$$
A_{L} T=U S_{L} U^{\sim}\left[\begin{array}{ll}
U_{1} & 0
\end{array}\right] V^{\sim}=U S_{L}\left[\begin{array}{ll}
I & 0 \\
0 & 0
\end{array}\right] V^{\sim}=A
$$

Point b) It is

$$
\begin{aligned}
A_{R} T^{\sim} T & =V S_{R} V^{\sim} V\left[\begin{array}{c}
U_{1}^{\sim} \\
0
\end{array}\right]\left[\begin{array}{ll}
U_{1} & 0
\end{array}\right] V^{\sim} \\
& =V S_{R}\left[\begin{array}{cc}
I & 0 \\
0 & 0
\end{array}\right] V^{\sim} \\
& =V S_{R} V^{\sim}=A_{R}
\end{aligned}
$$

Point c) It is

$$
\begin{aligned}
T^{\sim} T & =V\left[\begin{array}{c}
U_{1}^{\sim} \\
0
\end{array}\right]\left[\begin{array}{ll}
U_{1} & 0
\end{array}\right] V^{\sim} \\
& =V\left[\begin{array}{cc}
I & 0 \\
0 & 0
\end{array}\right] V^{\sim} \\
& \leq V\left[\begin{array}{ll}
I & 0 \\
0 & I
\end{array}\right] V^{\sim}=I
\end{aligned}
$$

Lemma F. 3 Consider the real matrix

$$
M(Z):=\left[\begin{array}{ll}
P & Q \\
R & Z
\end{array}\right]
$$

where the submatrices $P, Q$ and $R$ satisfy the equations

$$
\begin{aligned}
& P^{\prime} P+R^{\prime} R=\alpha^{2} I \\
& P P^{\prime}+Q Q^{\prime}=\alpha^{2} I
\end{aligned}
$$

with $\alpha$ a nonnegative scalar. Then
i)

$$
\min _{Z}\|M(Z)\|=\alpha
$$

ii) Letting $Z^{o}:=-T_{R} P^{\prime} T_{Q}$, where $T_{R}$ and $T_{Q}$ are matrices such that the equalities $R=T_{R}\left(R^{\prime} R\right)^{1 / 2}$ and $Q=\left(Q Q^{\prime}\right)^{1 / 2} T_{Q}$ hold, it is $\left\|M\left(Z^{\circ}\right)\right\|=\alpha$.

Proof Preliminarily observe that matrices $T_{R}$ and $T_{Q}$ exist thanks to Lemma F.2.
If $\alpha=0$, then, necessarily, $P=0, Q=0, R=0$, hence $Z^{o}=0$ and $\left\|M\left(Z^{o}\right)\right\|=0$.
If $\alpha>0$, set $A:=P / \alpha, B:=Q / \alpha, C:=R / \alpha$, so that

$$
\begin{align*}
& A^{\prime} A+C^{\prime} C=I  \tag{F.4}\\
& A A^{\prime}+B B^{\prime}=I \tag{F.5}
\end{align*}
$$

Moreover, for any real $D$ set

$$
N(D):=\left[\begin{array}{ll}
A & B \\
C & D
\end{array}\right]
$$

By recalling the definition of the norm of a matrix and the relevant properties (see Section 2.7), it follows, for each $D$,

$$
\begin{aligned}
\|N(D)\|^{2} & =\max _{\left[\begin{array}{ll}
x \\
y
\end{array}\right] \neq 0} \frac{\left[\begin{array}{ll}
x^{\prime} & y^{\prime}
\end{array}\right]\left[\begin{array}{cc}
A^{\prime} A+C^{\prime} C & A^{\prime} B+C^{\prime} D \\
B^{\prime} A+D^{\prime} C & B^{\prime} B+D^{\prime} D
\end{array}\right]\left[\begin{array}{l}
x \\
y
\end{array}\right]}{\left\|\left[\begin{array}{l}
x \\
y
\end{array}\right]\right\|^{2}} \\
& \geq \max _{x \neq 0} \frac{x^{\prime}\left(A^{\prime} A+C^{\prime} C\right) x}{\|x\|^{2}}=\left\|\left[\begin{array}{l}
A \\
C
\end{array}\right]\right\|^{2}
\end{aligned}
$$

Analogously,

$$
\|N(D)\|^{2}=\left\|N^{\prime}(D)\right\|^{2} \geq\left\|\left[\begin{array}{c}
A^{\prime} \\
B^{\prime}
\end{array}\right]\right\|^{2}
$$

Hence the conclusion is drawn that

$$
\|N(D)\| \geq \max \left[\left\|\left[\begin{array}{c}
A^{\prime}  \tag{F.6}\\
B^{\prime}
\end{array}\right]\right\|,\left\|\left[\begin{array}{c}
A \\
C
\end{array}\right]\right\|\right]=1, \quad \forall D
$$

It is now shown that, chosen

$$
\begin{equation*}
D^{o}:=-T_{C} A^{\prime} T_{B} \tag{F.7}
\end{equation*}
$$

where $T_{C}$ and $T_{B}$ are matrices (whose existence is guaranteed by Lemma F.2) such that

$$
\begin{align*}
& C=T_{C}\left(C^{\prime} C\right)^{1 / 2}  \tag{F.8}\\
& B=\left(B B^{\prime}\right)^{1 / 2} T_{B} \tag{F.9}
\end{align*}
$$

it is $\left\|N\left(D^{o}\right)\right\| \leq 1$, so that, in view of eq. (F.6), $\left\|N\left(D^{o}\right)\right\|=1$ and $D^{o}$ is optimal.
First notice that from eqs. (F.4),(F.5), (F.7)-(F.9), Lemmas F.1, F. 2 (points a) and $b$ ), it follows

$$
\begin{aligned}
C^{\prime} D^{o} & =-\left(C^{\prime} C\right)^{1 / 2} T_{C}^{\prime} T_{C} A^{\prime} T_{B} \\
& =-\left(C^{\prime} C\right)^{1 / 2} A^{\prime} T_{B} \\
& =-\left(I-A^{\prime} A\right)^{1 / 2} A^{\prime} T_{B} \\
& =-A^{\prime}\left(I-A A^{\prime}\right)^{1 / 2} T_{B} \\
& =-A^{\prime}\left(B B^{\prime}\right)^{1 / 2} T_{B}=-A^{\prime} B
\end{aligned}
$$

Therefore, $C^{\prime} D^{o}+A^{\prime} B=0$ and, by recalling eq. (F.4), one obtains

$$
N^{\prime}\left(D^{o}\right) N\left(D^{o}\right)=\left[\begin{array}{cc}
I & 0 \\
0 & B^{\prime} B+\left(D^{o}\right)^{\prime} D^{o}
\end{array}\right]
$$

However, by exploiting Lemma F. 2 (points $a$ ) and $c$ ), eqs. (F.5), (F.7) it is

$$
\begin{aligned}
B^{\prime} B+\left(D^{o}\right)^{\prime} D^{o} & =T_{B}^{\prime}\left(B B^{\prime}\right)^{1 / 2}\left(B B^{\prime}\right)^{1 / 2} T_{B}+T_{B}^{\prime} A T_{C}^{\prime} T_{C} A^{\prime} T_{B} \\
& \leq T_{B}^{\prime} B B^{\prime} T_{B}+T_{B}^{\prime} A A^{\prime} T_{B} \\
& \leq T_{B}^{\prime}\left(B B^{\prime}+A A^{\prime}\right) T_{B} \\
& \leq T_{B}^{\prime} T_{B} \leq I
\end{aligned}
$$

This implies that the maximum eigenvalue of $B^{\prime} B+\left(D^{o}\right)^{\prime} D^{o}$ is not greater than 1: consequently the maximum eigenvalue of $N^{\prime}\left(D^{o}\right) N\left(D^{o}\right)$ is equal to 1. Therefore, $\left\|N\left(D^{\circ}\right)\right\|=1$.

Now observe that for all $D$, letting $Z:=\alpha D$, it is $\alpha N(D)=M(\alpha D)=M(Z)$, so that

$$
\min _{Z}\|M(Z)\|=\min _{D}\|M(\alpha D)\|=\alpha \min _{D}\|N(D)\|
$$

and

$$
\begin{equation*}
Z^{o}=\alpha D^{o} \tag{F.10}
\end{equation*}
$$

On the other side, $R=\alpha C=\alpha T_{C}\left(C^{\prime} C\right)^{1 / 2}=T_{C}\left(R^{\prime} R\right)^{1 / 2}$ and, analogously, $Q=$ $\alpha B=\alpha\left(B B^{\prime}\right)^{1 / 2} T_{B}=\left(Q Q^{\prime}\right)^{1 / 2} T_{B}$ which imply $T_{C}=T_{R}$ and $T_{B}=T_{Q}$. From eqs. (F.10) and (F.7) and the definition of $A$ the expression of $Z^{o}$ follows.

Theorem F. 1 (Parrott's theorem) Let $B, C, D, P, X$ be real matrices such that $N(X):=\left[\begin{array}{cc}P & C \\ B & D+X\end{array}\right]$ and $\alpha:=\max \left[\left\|\left[\begin{array}{c}P \\ B\end{array}\right]\right\|,\left\|\left[\begin{array}{l}P^{\prime} \\ C^{\prime}\end{array}\right]\right\|\right]$. Then
i) There exists $X^{o}$ such that $\alpha=\left\|N\left(X^{o}\right)\right\|=\min _{X}\|N(X)\|$.
ii) If $\|P\|<\alpha, \quad X^{o}:=-D-B\left(\alpha^{2} I-P^{\prime} P\right)^{-1} P^{\prime} C$.

Proof Preliminarily, notice that, in view of the definition of $\alpha$, the matrices

$$
\begin{aligned}
R & :=\left[\begin{array}{c}
B \\
\left(\alpha^{2} I-P^{\prime} P-B^{\prime} B\right)^{1 / 2}
\end{array}\right] \\
Q & :=\left[\begin{array}{ll}
C & \left(\alpha^{2} I-P P^{\prime}-C C^{\prime}\right)^{1 / 2}
\end{array}\right]
\end{aligned}
$$

are real.
Point i) Letting

$$
Z:=\left[\begin{array}{ll}
Z_{1} & Z_{2}  \tag{F.11}\\
Z_{3} & Z_{4}
\end{array}\right]
$$

with

$$
\begin{equation*}
Z_{1}:=D+X \tag{F.12}
\end{equation*}
$$

define the matrix

$$
M(Z):=\left[\begin{array}{ll}
P & Q \\
R & Z
\end{array}\right]
$$

Matrix $M(Z)$ satisfies the assumptions of Lemma F. 3 since

$$
\begin{align*}
& P^{\prime} P+R^{\prime} R=\alpha^{2} I  \tag{F.13}\\
& P P^{\prime}+Q Q^{\prime}=\alpha^{2} I \tag{F.14}
\end{align*}
$$

so that there exists a matrix

$$
Z^{o}:=\left[\begin{array}{ll}
Z_{1}^{o} & Z_{2}^{o} \\
Z_{3}^{o} & Z_{4}^{o}
\end{array}\right]
$$

such that $\left\|M\left(Z^{o}\right)\right\|=\alpha$.
Note that for any matrix $Z$ of the form (F.11),(F.12) it is $\left\|N\left(Z_{1}-D\right)\right\| \leq\|M(Z)\|$. Hence, by recalling that $\|N(X)\| \geq \alpha$, one obtains

$$
\alpha \leq\left\|N\left(Z_{1}^{o}-D\right)\right\| \leq\left\|M\left(Z^{o}\right)\right\|=\alpha
$$

which proves point $i$ ).
Point ii) Preliminarily observe that, thanks to the assumption on the norm of $P$, matrix $\alpha^{2} I-P^{\prime} P$ is nonsingular, so that $X^{\circ}$ is well defined.

By exploiting Lemma F. 3 one has

$$
\begin{align*}
Z_{1}^{o} & =\left[\begin{array}{ll}
I & 0
\end{array}\right] Z^{o}\left[\begin{array}{l}
I \\
0
\end{array}\right] \\
& =-\left[\begin{array}{ll}
I & 0
\end{array}\right] T_{R} P^{\prime} T_{Q}\left[\begin{array}{l}
I \\
0
\end{array}\right] \\
& =\left[\begin{array}{ll}
I & 0
\end{array}\right] T_{R} P^{\prime}\left(Q Q^{\prime}\right)^{-1 / 2} Q\left[\begin{array}{l}
I \\
0
\end{array}\right] \tag{F.15}
\end{align*}
$$

In writing down the last part of eq. (F.15) the nonsingularity of matrix $\left(Q Q^{\prime}\right)^{1 / 2}=$ $\left(\alpha^{2} I-P P^{\prime}\right)^{1 / 2}(\|P\|<\alpha)$ and the equality $Q=\left(Q Q^{\prime}\right)^{1 / 2} T_{Q}$ have been taken into account. In view of eqs. (F.13),(F.14) and Lemma F.1, from eq. (F.15) it follows

$$
\begin{aligned}
Z_{1}^{o} & =-\left[\begin{array}{ll}
I & 0
\end{array}\right] T_{R} P^{\prime}\left(\alpha^{2} I-P P^{\prime}\right)^{-1 / 2} C \\
& =-\left[\begin{array}{ll}
I & 0
\end{array}\right] T_{R}\left(\alpha^{2} I-P^{\prime} P\right)^{-1 / 2} P^{\prime} C \\
& =-\left[\begin{array}{ll}
I & 0
\end{array}\right] T_{R}\left(R^{\prime} R\right)^{-1 / 2} P^{\prime} C \\
& =-\left[\begin{array}{ll}
I & 0
\end{array}\right] R\left(R^{\prime} R\right)^{-1} P^{\prime} C \\
& =-B\left(\alpha^{2} I-P^{\prime} P\right)^{-1} P^{\prime} C
\end{aligned}
$$

By recalling that $\left\|N\left(Z_{1}^{o}\right)\right\|=\alpha$, point $\left.i i\right)$ follows.

## Appendix G

## Worst Case Analysis

This appendix makes basically reference to the paper by Doyle et al. [17].
Lemma G. 1 Consider the system

$$
\begin{aligned}
\dot{x} & =A x+B q \\
z & =C x
\end{aligned}
$$

with $A$ stable, $x(0)=x_{0}$ and $\|G(s)\|_{\infty}<\gamma$, where $G(s):=C(s I-A)^{-1} B$. Then

$$
\sup _{q \in R H_{2}}\left[\|z\|_{2}^{2}-\gamma^{2}\|q\|_{2}^{2}\right]=x_{0}^{\prime} Q_{\infty} x_{0}
$$

where $Q_{\infty}$ is the symmetric, positive semidefinite and stabilizing solution of the Riccati equation (in the unknown $Q$ )

$$
0=Q A+A^{\prime} Q+\gamma^{-2} Q B B^{\prime} Q+C^{\prime} C
$$

that is such that the matrix $A+\gamma^{-2} B B^{\prime} Q_{\infty}$ is stable.
Proof The existence of the solution of the Riccati equation is guaranteed by the assumption $\|G(s)\|_{\infty}<\gamma$ (Theorem 2.13). Further, it is

$$
\begin{aligned}
\frac{d}{d t} x^{\prime} Q_{\infty} x & =x^{\prime} A^{\prime} Q_{\infty} x+q^{\prime} B^{\prime} Q_{\infty} x+x^{\prime} Q_{\infty} A x+x^{\prime} Q_{\infty} B q \\
& =x^{\prime}\left(Q_{\infty} A+A^{\prime} Q_{\infty}\right) x+q^{\prime} B^{\prime} Q_{\infty} x+x^{\prime} Q_{\infty} B q \\
& =-x^{\prime}\left(C^{\prime} C+\gamma^{-2} Q_{\infty} B B^{\prime} Q_{\infty}\right) x+q^{\prime} B^{\prime} Q_{\infty} x+x^{\prime} Q_{\infty} B q \\
& =-x^{\prime} C^{\prime} C x-\left(\gamma q-\gamma^{-1} B^{\prime} Q_{\infty} x\right)^{\prime}\left(\gamma q-\gamma^{-1} B^{\prime} Q_{\infty} x\right)+\gamma^{2} q^{\prime} q \\
& =-\|z\|^{2}+\gamma^{2}\|q\|^{2}-\gamma^{2}\left\|q-\gamma^{-2} B^{\prime} Q_{\infty} x\right\|^{2}
\end{aligned}
$$

Observe that $x \in R H_{2}$ and $z \in R H_{2}$ whenever $q \in R H_{2}$, since $A$ is stable. Therefore, by integrating over the interval $[0, \infty)$ the first and last terms of the this series of equalities one obtains, for each $q \in R H_{2}$,

$$
\|z\|_{2}^{2}-\gamma^{2}\|q\|_{2}^{2}=x_{0}^{\prime} Q_{\infty} x_{0}-\gamma^{2}\left\|q-\gamma^{-2} B^{\prime} Q_{\infty} x\right\|_{2}^{2} \leq x_{0}^{\prime} Q_{\infty} x_{0}
$$

The choice $q=\gamma^{-2} B^{\prime} Q_{\infty} x$ is consistent with the above derivation since $\dot{x}=(A+$ $\left.\gamma^{-2} B B^{\prime} Q_{\infty}\right) x$ is a stable system by assumption. With this choice it is $\|z\|_{2}^{2}-\gamma^{2}\|q\|_{2}^{2}=$ $x_{0}^{\prime} Q_{\infty} x_{0}$, the maximum possible value.

In the following lemma reference is made to the stable $\left(\Pi_{s}\right)$ and antistable $\left(\Pi_{a}\right)$ orthogonal projections which have been introduced in Definition 2.33.

Lemma G. 2 Consider the system

$$
\begin{align*}
\dot{x} & =A x+B_{1} q_{1}+B_{2} q_{2}  \tag{G.1}\\
z & =C x \tag{G.2}
\end{align*}
$$

with $A$ stable. Further, consider the subspace of $R L_{2}$

$$
\mathcal{Q}:=\left\{\left[\begin{array}{ll}
q_{1}^{\prime} & q_{2}^{\prime}
\end{array}\right]^{\prime}:=q, \quad q_{1} \in R H_{2}^{\perp}, \quad q_{2} \in R L_{2}\right\}
$$

and the operator $\Xi: \mathcal{Q} \rightarrow R H_{2}$ defined by

$$
\Xi: q \mapsto \Xi q:=\Pi_{s}\left[G_{1}(s) q_{1}+G_{2}(s) q_{2}\right]
$$

where, for $i=1,2, G_{i}(s):=C(s I-A)^{-1} B_{i}$. Then, denoting with $\gamma$ a positive scalar,

$$
\begin{equation*}
\sup _{\substack{q \in \mathcal{Q} \\\|q\|_{2}=1}}\|\Xi q\|_{2}<\gamma \tag{G.3}
\end{equation*}
$$

if and only if the following two conditions hold:
i) There exists the symmetric, positive semidefinite and stabilizing solution $Q_{\infty}$ of the Riccati equation (in the unknown $Q$ )

$$
0=Q A+A^{\prime} Q+\gamma^{-2} Q B_{2} B_{2}^{\prime} Q+C^{\prime} C
$$

that is such that the matrix $A+\gamma^{-2} B_{2} B_{2}^{\prime} Q_{\infty}$ is stable.
ii) $r_{s}\left(Q_{\infty} L_{r}\right)<\gamma^{2}$, where $L_{r}$ is the solution of the Lyapunov equation (in the unknown $L$ )

$$
0=L A^{\prime}+A L+B_{1} B_{1}^{\prime}+B_{2} B_{2}^{\prime}
$$

Proof Observe that

$$
\begin{aligned}
\sup _{\substack{q \in \mathcal{Q} \\
\|q\|_{2}=1}}\|\Xi q\|_{2} & \geq \sup _{\substack{q \in \mathcal{Q} \\
\|q\|_{2}=1, q_{1}=0}}\|\Xi q\|_{2} \\
& \geq \sup _{\substack{q_{2} \in R L_{2} \\
\left\|q_{2}\right\|_{2}=1}}\left\|\Pi_{s} G_{2}(s) q_{2}\right\|_{2} \\
& \geq \sup _{\substack{q_{2} \in R H_{2} \\
\left\|q_{2}\right\|_{2}=1}}\left\|G_{2}(s) q_{2}\right\|_{2}=\left\|G_{2}(s)\right\|_{\infty}
\end{aligned}
$$

where the last equality sign follows from Theorem 2.12. If condition (G.3) holds, then Theorem 2.13 proves the necessity of condition (i). Therefore, the proof is carried out by showing that, under condition (i), condition (ii) is necessary and sufficient for condition (G.3) to hold.

Notice that, for any $q \in \mathcal{Q}$, it is

$$
q=\Pi_{s} q+\Pi_{a} q=\left[\begin{array}{c}
0 \\
\Pi_{s} q_{2}
\end{array}\right]+\Pi_{a} q
$$

Having decomposed $q$ into two orthogonal terms implies

$$
\begin{equation*}
\|\Xi q\|_{2}^{2}-\gamma^{2}\|q\|_{2}^{2}=\|\Xi q\|_{2}^{2}-\gamma^{2}\left\|\Pi_{s} q_{2}\right\|_{2}^{2}-\gamma^{2}\left\|\Pi_{a} q\right\|_{2}^{2} \tag{G.4}
\end{equation*}
$$



Figure G.1: The block-scheme representation of $\Xi$

Thanks to the linearity of $\Xi$, condition (G.3) is equivalent to (also recall Remark 2.13)

$$
\sup _{0 \neq q \in \mathcal{Q}} \frac{\|\Xi q\|_{2}^{2}}{\|q\|_{2}^{2}}<\gamma^{2}
$$

which implies that, for each $q \in \mathcal{Q}, q \neq 0,\|\Xi q\|_{2}^{2}-\gamma^{2}\|q\|_{2}^{2}<0$. Therefore, condition (G.3) is violated if and only if

$$
\sup _{0 \neq q \in \mathcal{Q}}\left[\|\Xi q\|_{2}^{2}-\gamma^{2}\|q\|_{2}^{2}\right] \geq 0
$$

By taking into account eq. (G.4), condition (G.3) is violated if and only if

$$
\begin{equation*}
\sup _{0 \neq q \in \mathcal{Q}}\left[\|\Xi q\|_{2}^{2}-\gamma^{2}\left\|\Pi_{s} q\right\|_{2}^{2}-\gamma^{2}\left\|\Pi_{a} q\right\|_{2}^{2}\right] \geq 0 \tag{G.5}
\end{equation*}
$$

With reference to fig. G. 1 (which supplies a representation of the operator $\Xi$ where the input $q_{2}$ has been decomposed into its stable and antistable component), recall that, in the time domain, a signal belonging to $R H_{2}^{\perp}$ is zero for $t>0$, whereas a signal belonging to $R H_{2}$ is zero for $t<0$. Therefore, the value of the state of system (G.1),(G.2) at time $t=0$ only depends on the inputs which belong to $R H_{2}^{\perp}$ so that, given $x(0):=x_{0}$, the maximization of the first two terms in eq. (G.5) can be performed by ignoring the antistable part of the inputs to system (G.1),(G.2), so that

$$
\left.\sup _{0 \neq q \in \mathcal{Q}}\left[\|\Xi q\|_{2}^{2}-\gamma^{2}\left\|\Pi_{s} q\right\|_{2}^{2}\right]\right|_{x(0)=x_{0}}=\left.\sup _{0 \neq q_{2} \in R H_{2}}\left[\|z\|_{2}^{2}-\gamma^{2}\left\|q_{2}\right\|_{2}^{2}\right]\right|_{x(0)=x_{0}}
$$

Therefore,

$$
\begin{align*}
& \left.\sup _{0 \neq q \in \mathcal{Q}}\left[\|\Xi q\|_{2}^{2}-\gamma^{2}\left\|\Pi_{s} q\right\|_{2}^{2}-\gamma^{2}\left\|\Pi_{a} q\right\|_{2}^{2}\right]\right|_{x(0)=x_{0}}= \\
& =\left.\sup _{0 \neq q_{2} \in R H_{2}}\left[\|z\|_{2}^{2}-\gamma^{2}\left\|q_{2}\right\|_{2}^{2}\right]\right|_{x(0)=x_{0}} \quad-\left.\inf _{q \in R H_{2}^{\perp}} \gamma^{2}\|q\|_{2}^{2}\right|_{x(0)=x_{0}} \tag{G.6}
\end{align*}
$$

Observe that $\sup [M-N]=\sup [M]-\inf [N]$ whenever the choice of $M$ is independent of the choice of $N$. This fact has been exploited in writing down eq. (G.6).

Assume, for the moment being, that the pair $\left(A,\left[B_{1} B_{2}\right]\right)$ is reachable. Thanks to Lemma G.1, the first term on the right hand side of eq. (G.6) equals $x_{0}^{\prime} Q_{\infty} x_{0}$, while the second term equals, as it will be shown later on, $-\gamma^{2} x_{0}^{\prime} L_{r}^{-1} x_{0}$ (observe that
$L_{r}>0$ since the pair $\left(A,\left[\begin{array}{ll}B_{1} & B_{2}\end{array}\right]\right)$ is reachable, see Lemma C.1). Thus, eq. (G.6) becomes

$$
\begin{equation*}
\left.\sup _{0 \neq q \in \mathcal{Q}}\left[\|\Xi q\|_{2}^{2}-\gamma^{2}\left\|\Pi_{s} q\right\|_{2}^{2}-\gamma^{2}\left\|\Pi_{a} q\right\|_{2}^{2}\right]\right|_{x(0)=x_{0}}=x_{0}^{\prime}\left(Q_{\infty}-\gamma^{2} L_{r}^{-1}\right) x_{0} \tag{G.7}
\end{equation*}
$$

Therefore, condition (G.5) is verified if and only if there exists $x_{0}$ such that the right hand side of eq. (G.7) is nonnegative. Hence, condition (G.3) is violated if and only if there exists $x_{0}$ such that

$$
x_{0}^{\prime}\left(Q_{\infty}-\gamma^{2} L_{r}^{-1}\right) x_{0} \geq 0
$$

that is if and only if $r_{s}\left(Q_{\infty} L_{r}\right) \geq \gamma^{2}$ (recall Lemma B.12).
The equality

$$
\left.\inf _{q \in R H_{2}^{\frac{1}{2}}}\|q\|_{2}^{2}\right|_{x(0)=x_{0}}=x_{0}^{\prime} L_{r}^{-1} x_{0}
$$

can be proved by making reference to the (classical) optimal control problem relative to the system

$$
\begin{aligned}
& \dot{x}=A x+B q \\
& x(\tau)=0, \quad x(0)=x_{0}, \quad \tau<0
\end{aligned}
$$

and the performance index

$$
J_{\tau}=\int_{\tau}^{0} q^{\prime} q d t
$$

where $\tau$ is a given parameter. The Hamilton-Jacobi theory leads to recognizing

$$
q_{\tau}^{o}(t)=B^{\prime} e^{-A^{\prime} t}\left[\int_{\tau}^{0} e^{-A v} B B^{\prime} e^{-A^{\prime} v} d v\right]^{-1} x_{0}, \quad t \leq 0
$$

as an optimal control, so that, as $\tau \rightarrow-\infty$ (recall eq. (2.21) and what has been said just after) it is

$$
\begin{aligned}
q^{o}(t) & =\lim _{\tau \rightarrow-\infty} q_{\tau}^{o}(t) \\
& =B^{\prime} e^{-A^{\prime} t}\left[\int_{-\infty}^{0} e^{-A v} B B^{\prime} e^{-A^{\prime} v} d v\right]^{-1} x_{0} \\
& =B^{\prime} e^{-A^{\prime} t} L_{r}^{-1} x_{0}, \quad t \leq 0
\end{aligned}
$$

Notice that the so obtained control belongs to $R H_{2}^{\perp}$ since $A$ is stable. Finally,

$$
\begin{aligned}
\left.\inf _{q \in R H_{2}^{+}}\|q\|_{2}^{2}\right|_{x(0)=x_{0}} & =\int_{-\infty}^{0} q^{o}(t)^{\prime} q^{o}(t) d t \\
& =x_{0}^{\prime} L_{r}^{-1}\left[\int_{-\infty}^{0} e^{-A t} B B^{\prime} e^{-A^{\prime} t} d t\right] L_{r}^{-1} x_{0} \\
& =x_{0}^{\prime} L_{r}^{-1} x_{0}
\end{aligned}
$$

The reachability assumption of the pair $\left(A,\left[B_{1} B_{2}\right]\right)$ is now relaxed. Assume to have already decomposed the system into the reachable and unreachable parts, so that $x:=\left[\begin{array}{ll}x_{r}^{\prime} & x_{u}^{\prime}\end{array}\right]^{\prime}$ and

$$
A:=\left[\begin{array}{cc}
A_{r} & A_{x} \\
0 & A_{u}
\end{array}\right], \quad B:=\left[\begin{array}{cc}
B_{1 r} & B_{2 r} \\
0 & 0
\end{array}\right], C:=\left[\begin{array}{ll}
C_{r} & C_{u}
\end{array}\right]
$$

where the pair ( $A_{r},\left[B_{1 r} B_{2 r}\right]$ ) is reachable. Obviously, the states $x_{0}$ to be considered in eq. (G.6) are only those belonging to the set $\mathcal{S}_{0}:=\left\{x \mid x=\left[x_{r 0}^{\prime} 0\right]^{\prime}\right\}$, so that condition (G.5) is satisfied (hence condition (G.3) is violated) if and only if there exists $x_{0} \in \mathcal{S}_{0}$ corresponding to which it is

$$
\begin{equation*}
\left.\sup _{0 \neq q_{2} \in R H_{2}}\left[\|z\|_{2}^{2}-\gamma^{2}\left\|q_{2}\right\|_{2}^{2}\right]\right|_{x(0)=x_{0}}-\left.\inf _{q \in R H_{2}^{\frac{1}{2}}} \gamma^{2}\|q\|_{2}^{2}\right|_{x(0)=x_{0}} \geq 0 \tag{G.8}
\end{equation*}
$$

Observe that the stabilizing solution $Q_{\infty}$ of the Riccati equation is given the form

$$
Q_{\infty}=\left[\begin{array}{ll}
Q_{r} & Q_{x} \\
Q_{x}^{\prime} & Q_{u}
\end{array}\right]
$$

where $Q_{r}$ is the symmetric, positive definite and stabilizing solution of the Riccati equation (in the unknown $V$ )

$$
0=V A_{r}+A_{T}^{\prime} V+\gamma^{-2} V B_{2 r} B_{2 r}^{\prime} V+C_{r}^{\prime} C_{r}
$$

As already found, the first term on the left hand side of eq. (G.8) equals $x_{r 0}^{\prime} Q_{r} x_{r 0}$, while, by again exploiting the previous discussion about the optimal control problem (relative to the reachable part only), it is easy to find that the second term on the left hand side of eq. (G.8) is given by $-\gamma^{-2} x_{r 0}^{\prime} L_{r r}^{-1} x_{r 0}$, where $L_{r r}$ is the solution of the Lyapunov equation (in the unknown $W$ )

$$
0=W A_{r}^{\prime}+A_{r} W+B_{1 r} B_{1 r}^{\prime}+B_{2 r} B_{2 r}^{\prime}
$$

Therefore, condition (G.3) is violated if and only if $r_{s}\left(Q_{r} L_{r r}\right) \geq \gamma^{2}$. It is straightforward to check that, thanks to the peculiar structure of matrices $A$ and $\left[B_{1} B_{2}\right]$, the unique solution of the Lyapunov equation referred to in the statement of the theorem is

$$
L_{r}=\left[\begin{array}{cc}
L_{r r} & 0 \\
0 & 0
\end{array}\right]
$$

so that $r_{s}\left(Q_{r} L_{r r}\right) \geq \gamma^{2}$ is equivalent to the condition $r_{s}\left(Q_{\infty} L_{r}\right) \geq \gamma^{2}$.

## Appendix H

## Convex Functions and Sets

The material included in this appendix is standard in convex optimization. It is mainly based on books [37], [41] and [53]. However, to our present needs, it is important to keep in mind that here all variables and functions are defined in some subset of $R^{n \times m}$. In other terms, instead of vectors we work with $n \times m$ dimensional matrices. All sets to be handled are assumed closed.

First, let us introduce the algebraic structure to be dealt with. The inner product of two matrices $X$ and $Y$ belonging to $R^{n \times m}$ is defined as

$$
<X, Y>:=\operatorname{trace}\left[X^{\prime} Y\right]
$$

which induces the so called Frobenius norm, that is

$$
\|X\|_{F}:=\sqrt{<X, X>}=\sqrt{\operatorname{trace}\left[X^{\prime} X\right]}
$$

The geometric interpretation is clear. Matrix $X \in R^{n \times m}$ may be thought as a point in $R^{n \times m}$ whose distance to the origin is $\|X\|_{F}$ just defined.

Definition H. 1 (Convex sets) $A$ set $\Omega$ in $R^{n \times m}$ is convex if $\forall X_{1}, X_{2} \in \Omega$ the point $X=\alpha X_{1}+(1-\alpha) X_{2} \in \Omega$ for every $\alpha \in[0,1]$.

In other words the line segment between any two points in a convex set is entirely contained in the set. The empty set is, by definition, a convex set. Furthermore, the above definition can be alternatively stated in terms of a number, say $N \geq 2$, of feasible points. Actually, for $X_{i} \in \Omega, i=1,2, \cdots, N$ the convexity of $\Omega$ assures that

$$
\begin{equation*}
X=\sum_{i=1}^{N} \xi_{i} X_{i} \in \Omega \tag{H.1}
\end{equation*}
$$

for all scalars $\xi_{i}, i=1,2, \cdots, N$ such that

$$
\begin{equation*}
\xi_{i} \geq 0, \quad \sum_{i=1}^{N} \xi_{i}=1 \tag{H.2}
\end{equation*}
$$

This fact puts in evidence a very interesting property of convex sets. Given a bounded convex set $\Omega$ it is always possible to choose $N$ (possibly infinite) points $X_{i} \in \Omega, i=1,2, \cdots, N$ such that all $X \in \Omega$ can be written as (H.1). The points $X_{i}$ with this property are called extreme points of $\Omega$. A precise definition of this concept is as follows.

Definition H. 2 (Extreme point) Let $X$ be a point belonging to a convex set $\Omega$. It is an extreme point of $\Omega$ if there are no points $X_{i} \neq X$ in $\Omega, i=1,2, \cdots, N$ such that, (H.1) holds for some $\xi_{1}, \xi_{2}, \cdots, \xi_{N}$ satisfying (H.2).

Two convex sets of particular interest are the convex polyhedron and the convex cone. The first is defined as a convex set with a finite number of extreme points. The latter is a convex set such that if $X \in \Omega$ then $\lambda X \in \Omega$ for all $\lambda>0$. With a slight abuse of notation, we also call a convex cone any set for which the above property holds for all $\lambda \geq \lambda_{0}$ with $\lambda_{0}>0$.

Definition H. 3 (Convex hull) Let $\Gamma$ be a subset of $R^{n \times m}$. The convex hull of $\Gamma$ denoted $\operatorname{co}(\Gamma)$ is the smallest convex set containing $\Gamma$.

It is clear, from the above definition that $\Omega=\operatorname{co}(\Omega)$ whenever $\Omega$ is a convex set. Furthermore, if $\Omega$ is a convex polyhedron then

$$
\Omega=\operatorname{co}\left\{X_{i}, \quad i=1,2, \cdots, N\right\}
$$

where $X_{1}, X_{2}, \cdots, X_{N}$ are all extreme points of $\Omega$. The set of all symmetric positive definite matrices

$$
\Omega=\left\{X \in R^{n \times n}: X>0\right\}
$$

is a convex cone. Obviously, the same is true for the set of all symmetric negative definite matrices.

Let $\Lambda$ be a nonzero matrix in $R^{n \times m}$ and let $X_{0}$ be an arbitrary matrix (point) in $R^{n \times m}$. An hyperplane is the set of all points $X \in R^{n \times m}$ such that $X-X_{0}$ is orthogonal to $\Lambda$, or in more precise terms

$$
\left.0=<\Lambda, X-X_{0}\right\rangle=\operatorname{trace}\left[\Lambda^{\prime}\left(X-X_{0}\right)\right]
$$

Defining the scalar $c:=$ trace $\left[\Lambda^{\prime} X_{0}\right]$ an hyperplane can be characterized by all matrices $X \in R^{n \times m}$ such that

$$
\operatorname{trace}\left[\Lambda^{\prime} X\right]=c
$$

which puts in evidence that an hyperplane in $R^{n \times m}$ is nothing more than a linear variety of dimension $n m-1$. This concept is useful to get the following results valid for closed convex sets, that is, for convex sets which contain all boundary points.

Lemma H. 1 (Separating hyperplane) Let $\Omega$ be a closed convex set and consider $X_{0} \notin \Omega$. There exists a matrix $\Lambda_{0}$ and a scalar $c_{0}$ such that the hyperplane defined by $\left.<\Lambda_{0}, X\right\rangle=c_{0}$ separates $X_{0}$ from $\Omega$. That is, the following two conditions hold simultaneously :

$$
\begin{aligned}
& \text { i) }<\Lambda_{0}, X_{0}><c_{0} \\
& \text { ii) }<\Lambda_{0}, X>\geq c_{0}, \quad \forall X \in \Omega .
\end{aligned}
$$

Proof The proof comes from the definition of the new matrix

$$
Y_{0}:=\operatorname{argmin}\left\{\left\|X-X_{0}\right\|_{F}: X \in \Omega\right\}
$$

which enables us to get

$$
\Lambda_{0}:=Y_{0}-X_{0}, \quad c_{0}:=<Y_{0}-X_{0}, Y_{0}>
$$

such that both conditions are simultaneously verified.
Point i) Using the above matrix we have

$$
\begin{aligned}
<\Lambda_{0}, X_{0}>-c_{0} & =<Y_{0}-X_{0}, X_{0}>-<Y_{0}-X_{0}, Y_{0}> \\
& =-\left\|Y_{0}-X_{0}\right\|_{F}^{2}<0
\end{aligned}
$$

where the last inequality follows from the fact that $Y_{0} \neq X_{0}$.
Point ii) From the definition of matrix $Y_{0}$ we have

$$
\left\|X-X_{0}\right\|_{F}^{2} \geq\left\|Y_{0}-X_{0}\right\|_{F}^{2}, \quad \forall X \in \Omega
$$

However, for any $X \in \Omega$ given, the point $\lambda X+(1-\lambda) Y_{0} \in \Omega$ for all $\lambda \in[0,1]$ by a consequence of the convexity of $\Omega$. In this case, the above inequality reveals that

$$
<Y_{0}-X_{0}, X-Y_{0}>\geq-\left(\frac{\lambda}{2}\right)\left\|X-Y_{0}\right\|_{F}^{2}
$$

holds. Using this inequality for $\lambda \rightarrow 0$ we finally obtain

$$
<Y_{0}-X_{0}, X-Y_{0}>\geq 0, \quad \forall X \in \Omega
$$

This inequality, together with

$$
\begin{aligned}
<\Lambda_{0}, X>-c_{0} & =<Y_{0}-X_{0}, X>-<Y_{0}-X_{0}, Y_{0}> \\
& =<Y_{0}-X_{0}, X-Y_{0}>
\end{aligned}
$$

conclude the proof of the theorem proposed.
Lemma H. 2 (Supporting hyperplane) Let $\Omega$ be a closed convex set and consider $X_{0}$ in the boundary of $\Omega$. There exists a matrix $\Lambda_{0}$ and a scalar $c_{0}$ such that the hyperplane $<\Lambda_{0}, X>=c_{0}$ supports $\Omega$ at $X_{0}$. That is, the following two conditions hold simultaneously :

$$
\begin{aligned}
& \text { i) }<\Lambda_{0}, X_{0}>=c_{0} \\
& \text { ii) }<\Lambda_{0}, X>\geq c_{0}, \quad \forall X \in \Omega .
\end{aligned}
$$

Lemma H. 1 says that given a convex set and an exterior point then it is always possible to determine an hyperplane which contains the point in one of its half spaces while the convex set lays entirely in the other half space. Lemma H. 2 generalizes this result to cope with a point in the boundary of a closed convex set.

Let us now turns our attention to real valued functions defined in $R^{n \times m}$. A function $f(X): R^{n \times m} \rightarrow R$ is continuous if for each $X_{0} \in R^{n \times m}$

$$
\lim _{X \rightarrow X_{0}} f(X)=f\left(X_{0}\right)
$$

All functions handled here are continuous. The gradient of a continuous function $f(\cdot)$ at $X=X_{0}$ is a $n \times m$ matrix, denoted $\nabla f\left(X_{0}\right)$ and defined as

$$
\nabla f\left(X_{0}\right):=\left\{\frac{\partial f}{\partial X_{i j}}\left(X_{0}\right), \quad i=1,2, \cdots, n ; j=1,2, \cdots m\right\}
$$

For functions of matrices, a simple way to calculate gradients is from the concept of directional derivative. The directional derivative of a differentiable function $f(\cdot)$ at $X \in R^{n \times m}$ in the direction $Y \in R^{n \times m}$ is

$$
\begin{aligned}
D f(X, Y) & :=\lim _{\epsilon \rightarrow 0} \frac{f(X+\epsilon Y)-f(X)}{\epsilon} \\
& =<\nabla f(X), Y> \\
& =\operatorname{trace}\left[\nabla f(X)^{\prime} Y\right]
\end{aligned}
$$

In other words, given a function $f(\cdot)$, if for any matrix $Y \in R^{n \times m}$, it is possible to write

$$
\begin{equation*}
f(X+\epsilon Y)=f(X)+\epsilon \operatorname{trace}\left[G(X)^{\prime} Y\right]+\mathcal{O}\left(\epsilon^{2}\right) \tag{H.3}
\end{equation*}
$$

such that $\lim _{\epsilon \rightarrow 0} \mathcal{O}\left(\epsilon^{2}\right) / \epsilon=0$ then $f(\cdot)$ is differentiable and an adequate choice of directions leads to

$$
\begin{equation*}
\nabla f(X)=G(X) \tag{H.4}
\end{equation*}
$$

Example H. 1 Using the above result, it is simple to calculate the gradient of the function $f(X):=\operatorname{det}[X]$ where $X \in R^{n \times n}$ and nonsingular. Indeed, elementary matrices properties yield, for $\epsilon \in R$ arbitrarily small

$$
\begin{aligned}
\operatorname{det}[X+\epsilon Y] & =\operatorname{det}[X] \operatorname{det}\left[I+\epsilon X^{-1} Y\right] \\
& =\operatorname{det}[X] \prod_{i=1}^{n}\left[1+\epsilon \lambda_{i}\left(X^{-1} Y\right)\right] \\
& =\operatorname{det}[X]\left(1+\epsilon \operatorname{trace}\left[X^{-1} Y\right]\right)+\mathcal{O}\left(\epsilon^{2}\right)
\end{aligned}
$$

showing that $\nabla f(X)=\operatorname{det}[X]\left(X^{\prime}\right)^{-1}$. The same steps can be adopted to evaluate the gradient of the function $f(X):=\log \operatorname{det}[X]$ defined for all $X \in R^{n \times n}$ such that $\operatorname{det}[X]>0$. Indeed, from the above we get

$$
\begin{aligned}
f(X+\epsilon Y) & =f(X)+\sum_{i=1}^{n} \log \left[1+\epsilon \lambda_{i}\left(X^{-1} Y\right)\right] \\
& =f(X)+\epsilon \operatorname{trace}\left[X^{-1} Y\right]+\mathcal{O}\left(\epsilon^{2}\right)
\end{aligned}
$$

implying that $\nabla f(X)=\left(X^{\prime}\right)^{-1}$.
Remark H. 1 Equation (H.4) must be applied with care. In some important cases, it does not apply directly. To clarify its correct use, consider a function $f(\cdot): \mathcal{S} \rightarrow R$, where $\mathcal{S}$ denotes the set of all $n \times n$ symmetric matrices. It is clear that this is a function of $n(n+1) / 2$ independent variables. Assuming that (H.3) holds for $Y=Y^{t}$ we observe that

$$
\operatorname{trace}\left[G(X)^{\prime} Y\right]=\sum_{i=1}^{n} G(X)_{i i} Y_{i i}+2 \sum_{i=1}^{n} \sum_{j>i}^{n} G(X)_{i j} Y_{i j}
$$

which provides

$$
\nabla f(X)=G(X)+G(X)^{\prime}-\operatorname{diag}[G(X)]
$$

For instance, the gradient matrix of the linear function $f(X):=\operatorname{trace}\left[R^{\prime} X\right]$ defined for all square matrix $X \in R^{n \times n}$ is $\nabla f(X)=R$. On the other hand, the gradient of the same function defined in $\mathcal{S}$ is $\nabla f(X)=R+R^{\prime}-\operatorname{diag}[R]$.

Definition H. 4 (Convex functions) A function $f(\cdot) \in R$ defined in a convex set $\Omega$ is convex if $\forall X_{1}, X_{2} \in \Omega$ and $X=\alpha X_{1}+(1-\alpha) X_{2}$ there holds $f(X) \leq \alpha f\left(X_{1}\right)+$ $(1-\alpha) f\left(X_{2}\right)$ for every $\alpha \in[0,1]$.

There are several different but equivalent ways to test a function for convexity. Of course each one depends on the previous informations we have for the function we are work with. The next lemma provides some equivalent tests for twice differentiable functions.

Lemma H. 3 (Convexity) Let $f(\cdot)$ be a twice differentiable function, defined in a convex set $\Omega \subset R^{n \times m}$. Assume that $\Omega$ contains an interior point. The following are equivalent :
i) $f(\cdot)$ is convex over $\Omega$.
ii) $f(Y) \geq f(X)+<\nabla f(X), Y-X>, \quad \forall X, Y \in \Omega$.
iii) The second order variation

$$
V f(X, Y):=\left.\frac{d^{2} f}{d \epsilon^{2}}(X+\epsilon Y)\right|_{\epsilon=0}
$$

is such that $V f(X, Y) \geq 0$ for all $X \in \Omega$ and all $Y \in R^{n \times m}$.
Dealing with optimal control problems, it is frequently necessary to handle nondifferentiable functions. The next lemma characterizes convexity in this important case.

Lemma H. 4 (Convexity) Let $f(\cdot)$ be a function, defined in a convex set $\Omega \subset R^{n \times m}$. Assume that $\Omega$ contains an interior point. The following are equivalent:
i) $f(\cdot)$ is convex over $\Omega$.
ii) For each $X \in \Omega$, there exits a matrix $\Lambda(X)$ with finite norm such that

$$
f(Y) \geq f(X)+<\Lambda(X), Y-X>, \quad \forall Y \in \Omega
$$

iii) The epigraph of function $f$, namely

$$
\text { epi } f:=\{(X, \gamma): X \in \Omega, \quad f(X) \leq \gamma\}
$$

is a convex set.
Proof The equivalence between points $i$ ) and $i i i$ ) follows immediately from the observation that $f$ is convex over $\Omega$ if and only if function $g(X, \gamma):=f(X)-\gamma$ is convex over $\Omega \times R$.
ii) $\Longrightarrow i)$ Consider $X_{1}$ and $X_{2}$ two arbitrary matrices in $\Omega$. Define $X=\alpha X_{1}+$ $(1-\alpha) X_{2}$ for $\alpha \in[0,1]$. Setting $Y=X_{1}$ and $Y=X_{2}$ we have

$$
\begin{aligned}
& f\left(X_{1}\right) \geq f(X)+<\Lambda(X), X_{1}-X> \\
& f\left(X_{2}\right) \geq f(X)+<\Lambda(X), X_{2}-X>
\end{aligned}
$$

Multiplying the first inequality by $\alpha$, the second by $1-\alpha$ and adding terms we get

$$
\alpha f\left(X_{1}\right)+(1-\alpha) f\left(X_{2}\right) \geq f(X)
$$

which is the desired result.
$i) \Longrightarrow i i)$ Under the assumption that $f$ is convex, from part $i i i$ ) the set epi $f$ is convex as well. Using Lemma H.2, there exists a supporting hyperplane to epi $f$ in
the boundary point $\left(X_{0}, \gamma_{0}\right)$ where $\gamma_{0}=f\left(X_{0}\right)$ and $X_{0} \in \Omega$. That is, there exist a matrix $\Lambda_{0}$ and a scalar $\lambda_{0}$ such that all $(X, \gamma) \in$ epi $f$ satisfy

$$
<\Lambda_{0}, X-X_{0}>+\lambda_{0}\left(\gamma-\gamma_{0}\right) \geq 0
$$

We notice that the above inequality imposes $\lambda_{0}>0$. In fact, it is simple to check that $\left(X_{0}, \gamma\right) \notin$ epi $f$ provided $\gamma<\gamma_{0}$ leading to $\lambda_{0}\left(\gamma-\gamma_{0}\right)<0$. As a consequence we have

$$
\gamma \geq f\left(X_{0}\right)+<-\lambda_{0}^{-1} \Lambda_{0}, X-X_{0}>
$$

which turns out to be true for all $(X, \gamma) \in$ epi $f$. Then, it is verified for the particular point $(X, \gamma)$ with $\gamma=f(X)$ of epi $f$. Defining $\Lambda\left(X_{0}\right):=-\Lambda_{0} / \lambda_{0}$, the proof is concluded.

The level set of a function $f(\cdot)$ defined in $\Omega$ is the set of all $X \in \Omega$ such that $f(X)$ is not greater than a fixed value. More precisely, let $\alpha$ be a fixed scalar, the level set of $f$ is

$$
L_{\alpha} f:=\{X: X \in \Omega, \quad f(X) \leq \alpha\}
$$

It must be clear that the level set of a function is a subset of $\Omega$ for each value of $\alpha$ given. On the contrary, the epi $f$ is a subset of $\Omega \times R$. Then, it is not surprising that while, convexity of $f$ is equivalent to the convexity of epi $f$, the same is not true for $L_{\alpha} f$. Indeed, if $f$ is convex in $\Omega$ the level set $L_{\alpha} f$ is convex for all $\alpha \in R$ but the latter statement is not generally sufficient to assure convexity of $f$.

Definition H. 5 (Subgradient) Let $f(\cdot)$ be a convex function defined in a convex set $\Omega \subset R^{n \times m}$. Matrix $\Lambda \in R^{n \times m}$ is said to be a subgradient of $f$ at a point $X$ if

$$
f(Y) \geq f(X)+<\Lambda, Y-X>
$$

holds for all $Y \in \Omega$. The set of all subgradients of $f$ at $X$ is denoted $\partial f(X)$.
Lemma H. 5 The set $\partial f(X)$ is convex.
Proof Consider $\Lambda_{1}, \Lambda_{2} \in \partial f(X)$. From Definition H. 5 we have, for all $Y \in \Omega$

$$
\begin{aligned}
& f(Y) \geq f(X)+<\Lambda_{1}, Y-X> \\
& f(Y) \geq f(X)+<\Lambda_{2}, Y-X>
\end{aligned}
$$

Multiplying the first inequality by $\alpha$, the second one by $1-\alpha$ and adding terms, yields

$$
f(Y) \geq f(X)+<\alpha \Lambda_{1}+(1-\alpha) \Lambda_{2}, Y-X>
$$

implying that $\Lambda:=\alpha \Lambda_{1}+(1-\alpha) \Lambda_{2} \in \partial f(X)$, which is the stated property.
It is important to mention that the set of subgradients is in addition a closed set. An important concept to deal with nondifferentiable convex functions is the one sided directional derivative. The right directional derivative of function $f$ at $X \in R^{n \times m}$ in the direction $Y \in R^{n \times m}$ is

$$
D_{+} f(X, Y):=\lim _{\epsilon \rightarrow 0^{+}} \frac{f(X+\epsilon Y)-f(X)}{\epsilon}
$$

as well as, the left directional derivative in the same direction is

$$
D_{\ldots} f(X, Y):=\lim _{\epsilon \rightarrow 0^{-}} \frac{f(X+\epsilon Y)-f(X)}{\epsilon}
$$

Simple calculations show that each one sided directional derivative is one related to the other by the equation

$$
\begin{equation*}
D_{-} f(X, Y)=-D_{+} f(X,-Y) \tag{H.5}
\end{equation*}
$$

indicating that, generally, they are different. Indeed, if they coincide then the function under consideration is differentiable at $X$. This important feature is now discussed from the result of the next lemma.
Lemma H. 6 Let $f(\cdot)$ be a convex function defined in the convex set $\Omega$. At any interior point of $\Omega$, it follows that

$$
D_{+} f(X, Y)=\max \left\{\operatorname{trace}\left[\Lambda^{\prime} Y\right]: \Lambda \in \partial f(X)\right\}
$$

Proof For any $\Lambda \in \partial f(X)$, the definition of subgradient yields

$$
\begin{aligned}
D_{+} f(X, Y) & =\lim _{\epsilon \rightarrow 0^{+}} \frac{f(X+\epsilon Y)-f(X)}{\epsilon} \\
& \geq \operatorname{trace}\left[\Lambda^{\prime} Y\right]
\end{aligned}
$$

the proof is then concluded from the fact that at any interior point of $\Omega$ the set $\partial f(X)$ is closed and bounded and both sides of the above inequality must be equal for some $\Lambda \in \partial f(X)$.

The above result used together (H.5) enables us to get

$$
D_{-} f(X, Y)=\min \left\{\operatorname{trace}\left[\Lambda^{\prime} Y\right]: \Lambda \in \partial f(X)\right\}
$$

which implies that both, the right and the left directional derivatives are equal in any direction $Y \in R^{n \times m}$ whenever the set of subgradients contains only one element. In this case the function is differentiable, $\partial f(X)=\{\nabla f(X)\}$ and the directional derivative is a linear function of the direction $Y \in R^{n \times m}$.
Example H. 2 Let us calculate the set of subgradients $\partial g(X)$ of the function

$$
g(X):=\max \left\{f_{i}(X): i=1,2, \cdots, N\right\}
$$

where all functions $f_{i}(\cdot), i=1,2, \cdots, N$ are supposed to be convex and differentiable in all points of $R^{n \times m}$. It is a simple matter to verify that function $g$ is convex as well. Defining the set

$$
J(X):=\left\{j: f_{j}(X)=g(X)\right\}
$$

and taking $Y \in R^{n \times m}$ we get

$$
\begin{aligned}
g(Y) & =\max \left\{f_{i}(Y): i=1,2, \cdots, N\right\} \\
& \geq f_{j}(Y) \\
& \geq f_{j}(X)+<\nabla f_{j}(X), Y-X> \\
& \geq g(X)+<\nabla f_{j}(X), Y-X>
\end{aligned}
$$

and so $\nabla f_{j}(X) \in \partial g(X)$ for all $j \in J(X)$. Even more, $\nabla f_{j}(X)$ is an extreme point of $\partial g(X)$. Taking into account the result of Lemma H. 5 we then conclude that

$$
\partial g(X)=\operatorname{co}\left\{\nabla f_{j}(X): j \in J(X)\right\}
$$

Accordingly, if there is only one index, say $k \in J(X)$ then $g$ is differentiable and $\nabla g(X)=$ $\nabla f_{k}(X)$.


Figure H.1: Level sets of function $g(X)$

When dealing with nondifferentiable convex functions, it is important to keep in mind that many properties valid under differentiability may fail. For instance, if $f$ is differentiable at $X \in R^{n \times m}$ then $-Y=\nabla f(X) \neq 0$ is a direction such that the directional derivative is negative, as a consequence, there exists a step length $\epsilon>0$ such that $f(X+\epsilon Y)<f(X)$. On the other hand, if $f$ is not differentiable at $X \in R^{n \times m}$ this property may be lost. In other words, it does not hold if we just consider an arbitrary direction $-Y \in \partial f(X)$. This common situation is illustrated in fig. H.1. The function $g(X)$ is the one treated in Example H. 2 where $N=3$ and $f_{i}(X), i=1,2,3$ are linear. At point $X$ the function is not differentiable because $f_{1}(X)$ and $f_{2}(X)$ coincide. Both indicated directions are in the set $\partial g(X)$. Even though, the right directional derivative in the direction $Y=-Y_{2}$ is clearly negative but the right directional derivative in the direction $Y=-Y_{1}$ is positive. Hence, for $-Y=Y_{1} \in \partial g(X)$ the previously mentioned descent property does not hold. In the direction $Y=-Y_{1}$ the function increases.

For nondifferentiable functions, there is a way to determine the steepest descent direction. It is given by (notice the constraint on the maximum length of the direction)

$$
Y_{s t}:=\operatorname{argmin}\left\{D_{+} f(X, Y):\|Y\|_{F}^{2} \leq 1\right\}
$$

which, from Lemma H. 6 can be determined by solving

$$
\min _{\|Y\|_{F}^{2} \leq 1} \max _{\Lambda \in \partial f(X)} \operatorname{trace}\left[\Lambda^{\prime} Y\right]
$$

Recalling that all constraints are convex, closed and bounded, this problem is equivalent to

$$
\max _{\Lambda \in \partial f(X)} \min _{\|Y\|_{F}^{2} \leq 1} \operatorname{trace}\left[\Lambda^{\prime} Y\right]
$$

The minimization is readily solved by keeping in mind that the objective function is an inner product of two matrices whose minimum provides

$$
Y_{s t}:=-\frac{\Lambda_{s t}}{\left\|\Lambda_{s t}\right\|_{F}}
$$

where $\Lambda_{s t}$ is given by

$$
\Lambda_{s t}:=\operatorname{argmin}\left\{\|\Lambda\|_{F}: \Lambda \in \partial f(X)\right\}
$$

The steepest descent direction is obtained from the selection among all subgradients, the one of minimum Frobenius norm. This task is, in general, very difficult to be accomplished numerically since the set $\partial f(X)$ may not be explicitly known. In Example H. 2 (recall fig. H.1) $Y_{s t}=-Y_{2}$ is the steepest direction at point $X$. Finally, it is also interesting to observe that under differentiability we get $Y_{s t}=-\nabla f(X) /\|\nabla f(X)\|_{F}$ as expected. Furthermore, in the steepest descent direction the right directional derivative provides, for any other direction $Y$ the lower bound

$$
\begin{equation*}
D_{+} f(X, Y) \geq D_{+} f\left(X, Y_{s t}\right)=-\left\|\Lambda_{s t}\right\|_{F} \tag{H.6}
\end{equation*}
$$

It is obvious that if $0 \in \partial f(X)$, the calculation needed to determine the steepest descent direction $Y_{s t}$ can not be performed because $\Lambda_{s t}=0$. The possible occurrence of this fact is however of particular importance as indicated in the next lemma.

Lemma H. 7 Let $f(\cdot)$ be a convex function defined in $R^{n \times m}$. Matrix $X^{\star}$ minimizes $f$ if and only if $0 \in \partial f\left(X^{\star}\right)$.

Proof Suppose $0 \in \partial f\left(X^{\star}\right)$, from the convexity of $f$ we have, for all $X \in R^{n \times m}$

$$
f(X) \geq f\left(X^{\star}\right)+<0, X-X^{\star}>\geq f\left(X^{\star}\right)
$$

implying that $X^{\star}$ minimizes $f$ indeed. Conversely, if $X^{\star}$ minimizes $f$ then in any direction $Y \in R^{n \times m}$, the right directional derivative must be nonnegative. In view of (H.6) this occurs provided $\Lambda_{s t}=0$ implying that $0 \in \partial f\left(X^{\star}\right)$.

We are now in position to solve the following optimization problem

$$
\begin{equation*}
\min \{f(X): X \in \mathcal{X}\} \tag{H.7}
\end{equation*}
$$

which is called a convex programming problem provided $i$ ) the objective function $f(\cdot): \Omega \rightarrow R$ is convex in $\Omega$ which by its turn is a convex subset of $R^{n \times m}$ and $i i$ ) the constraint set $\mathcal{X}$ is a closed convex subset of $\Omega$. To avoid pathological cases, it is also assumed that $\mathcal{X}$ is bounded and contains an interior point.

Lemma H. 8 Consider the convex problem (H.7). The following hold :
i) If $X^{\star}$ is a minimum then it is a global minimum.
ii) Matrix $X^{\star}$ minimizes $f$ over $\mathcal{X}$ if and only if there exists $\Lambda^{\star} \in \partial f\left(X^{\star}\right)$ such that

$$
<\Lambda^{\star}, X-X^{\star}>\geq 0, \quad \forall X \in \mathcal{X}
$$

Proof Only point $i i$ ) will be proved since the first point restates an important but very known property of convex problems.

Point ii) Suppose there exists $\Lambda^{\star} \in \partial f\left(X^{\star}\right)$ with the above property. From the convexity of $f$ we have, for all $X \in \mathcal{X}$,

$$
f(X) \geq f\left(X^{\star}\right)+<\Lambda^{\star}, X-X^{\star}>\geq f\left(X^{\star}\right)
$$

implying that $X^{\star}$ minimizes $f$ indeed. Conversely, if $X^{\star}$ minimizes $f$ over $\mathcal{X}$ then, we first notice that due to the convexity of $\mathcal{X}$, for any $X \neq X^{\star} \in \mathcal{X}$, the direction $Y=X-X^{*}$ is always a feasible direction in the sense that

$$
X^{\star}+\epsilon Y \in \mathcal{X}
$$

for all $\epsilon \in[0,1]$. So, due to the fact that $X^{\star}$ is optimal we must have

$$
D_{+} f\left(X^{\star}, Y\right) \geq 0
$$

Then, from Lemma H.6, it is true that for all $X \in \mathcal{X}$

$$
\max _{\Lambda \in \partial f\left(X^{\star}\right)}<\Lambda, X-X^{\star}>\geq 0
$$

that is

$$
\min _{X \in \mathcal{X}} \max _{\Lambda \in \partial f\left(X^{\star}\right)}<\Lambda, X-X^{\star}>\geq 0
$$

Using the fact that all constraints in the above problem are convex, bounded and closed, it can be written in the equivalent form

$$
\max _{\Lambda \in \partial f\left(X^{\star}\right)} \min _{X \in \mathcal{X}}<\Lambda, X-X^{\star}>\geq 0
$$

which allows us to conclude that there exists $\Lambda^{\star} \in \partial f\left(X^{\star}\right)$ such that

$$
\min _{X \in \mathcal{X}}<\Lambda^{\star}, X-X^{\star}>\geq 0
$$

which is the same to say that $<\Lambda^{\star}, X-X^{\star}>\geq 0$ for all $X \in \mathcal{X}$ and the proof is concluded.

Once again, it is to be noticed that under differentiability, Lemma H. 7 reduces to the classical property $\nabla f\left(X^{\star}\right)=0$ to characterize a minimum of a convex function. As well as, under the same assumption, the second part of Lemma H. 8 imposes that $<\nabla f\left(X^{\star}\right), X-X^{\star}>\geq 0, \forall X \in \mathcal{X}$ for global optimality.

Problem (H.7) is frequently called Primal. With the purpose to solve it more efficiently, convexity allows us to determine another equivalent but in many instances easier problem to be solved called Dual. For convex problems the Primal and the Dual versions are equivalent in the sense that both provide the same optimal solution. Duality gap does not exist. The determination of the dual problem is based on the following result.
Lemma H. 9 (Minimax) Consider $\mathcal{X}$ and $\mathcal{Y}$ two closed and bounded convex sets. Let the continuous function $L(X, Y)$ be convex with respect to $X$ and concave with respect to $Y$. The following equality holds

$$
\begin{equation*}
\max _{Y \in \mathcal{Y}} \min _{X \in \mathcal{X}} L(X, Y)=\min _{X \in \mathcal{X}} \max _{Y \in \mathcal{Y}} L(X, Y) \tag{H.8}
\end{equation*}
$$

Lemma H. 9 states that under the given conditions, function $L(\cdot)$ admits a saddle point which is in fact the optimal solution of the min/max problem. Adding to (H.7) a new convex constraint, that is

$$
\min _{X \in \mathcal{X}}\{f(X): g(X) \leq 0\}
$$

and defining the associated Lagrangian function

$$
L(X, Y):=f(X)+<Y, g(X)>
$$

in the domain

$$
\mathcal{Y}:=\left\{Y \geq 0: \min _{X \in \mathcal{X}} L(X, Y)>-\infty\right\}
$$

the above lemma can be applied. Equality (H.8) provides a pair of equivalent problems. One of then is the primal while the other is the dual.

## Appendix I

## Convex Programming Numerical Tools

In this appendix some convex programming numerical tools are discussed. The material is mainly based on the natural generalization, to the nondifferentiable case, of two classical convex programming algorithms described in the classical books [37] and [41]. The method of centers is based on the book [10]. More efficient methods (but also much more involved from a theoretical point of view) can be found in reference [45].

The general form of a convex programming problem is

$$
\min \{f(X): X \in \mathcal{X}\}
$$

where $f(\cdot)$ is a convex function and $\mathcal{X}$ is a convex set. The concept of epigraph allows us to rewrite it in the equivalent form

$$
\min \{\gamma:(X, \gamma) \in \operatorname{epi} f\}
$$

which implies that, with no loss of generality, the objective function of a general convex programming problem can be considered to be a linear function. In addition, all problems we have manipulated in the previous chapters, are such that the convex set epi $f$ can be alternatively written (possibly by adding new variables and constraints) as a LMI (Linear Matrix Inequality). Consequently, the general problem we have to solve numerically is

$$
\begin{equation*}
\min \left\{c^{\prime} x: \mathcal{A}(x) \geq 0\right\} \tag{I.1}
\end{equation*}
$$

where $c$ and $x$ are $n$ dimensional real vectors and

$$
\mathcal{A}(x):=A_{0}+\sum_{i=1}^{n} x_{i} A_{i}
$$

with $A_{0}, A_{1}, \cdots, A_{n}$ being $m \times m$ dimensional symmetric matrices and $x_{i} \in R$ denotes the $i$-th entry of the $n$ dimensional vector $x$. It is assumed that the feasible set defined by all vectors $x \in R^{n}$ such that $\mathcal{A}(x) \geq 0$ is bounded. This assumption is verified in all problems of our interest already discussed in the previous chapters. Notice however that this set is convex but not polyhedral then, in general, it has an infinite number of extreme points.


Figure I.1: The set $\mathcal{A}(x) \geq 0$

There is no difficulty to consider more general problems with the structure (I.1) and with an additional equality constraint $B x=b$, where $B \in R^{r \times n}$ is a matrix of full row rank. If necessary, a preliminary change of variables, allows the partitioning $B=$ [ $B_{1} B_{2}$ ] where $B_{1} \in R^{r \times r}$ is nonsingular. Doing the same with vector $x^{\prime}=\left[x_{1}^{\prime} x_{2}^{\prime}\right]$, the linear constraint provides

$$
x_{1}=B_{1}^{-1} b-B_{1}^{-1} B_{2} x_{2}
$$

and adopting again the above partitioning to $c^{\prime}=\left[c_{1}^{\prime} c_{2}^{\prime}\right]$, the feasibility with respect to the same equality constraint yields

$$
c^{\prime} x=c_{1}^{\prime} B_{1}^{-1} b+\bar{c}_{2}^{\prime} x_{2}
$$

where $\bar{c}_{2}^{\prime}:=c_{2}^{\prime}-c_{1}^{\prime} B_{1}^{-1} B_{2}$ and

$$
\begin{aligned}
\mathcal{A}(x) & =\mathcal{A}\left(\left[\begin{array}{c}
B_{1}^{-1} b-B_{1}^{-1} B_{2} x_{2} \\
x_{2}
\end{array}\right]\right) \\
& :=\overline{\mathcal{A}}_{2}\left(x_{2}\right)
\end{aligned}
$$

The consequence is that the equality constraint can easily be handled by eliminating itself and formulating a new problem as (I.1) which depends exclusively on the variable $x_{2} \in R^{n-r}$. From the numerical point of view the remaining problem to be solved is simpler since it presents a smaller number of free variables.

Example I. 1 The geometry of the feasible set of problem (I.1) is now illustrated by means of a simple example. Consider the LMI $\mathcal{A}(x) \geq 0$ be given by $n=2$ and matrices

$$
A_{0}=\left[\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 / 2 & 0 \\
0 & 0 & 0 & 1
\end{array}\right], A_{1}=\left[\begin{array}{llll}
0 & 1 & 0 & 0 \\
1 & 0 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & 0
\end{array}\right], A_{2}=\left[\begin{array}{llll}
1 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 \\
0 & 0 & 0 & 1 \\
0 & 0 & 1 & 0
\end{array}\right]
$$

Figure I. 1 shows the convex set defined by this LMI. It is obtained by the intersection of the convex regions $x_{2} \geq x_{1}^{2}-1$ and $x_{1} \geq x_{2}^{2}-1 / 2$.

Algorithm I. 1 (Separating hyperplane algorithm) This is one of the simplest algorithms that can be applied to solve the stated convex problem (I.1). In the following it is discussed its convergence as well as its limitations to handle large scale problems.

1. Determine a convex polyhedral set $\mathcal{P}_{0}$ containing the overall feasible set of Problem I.1, that is

$$
\{x: \mathcal{A}(x) \geq 0\} \subset \mathcal{P}_{0}
$$

and set the iteration index $k=0$.
2. Solve the linear programming problem

$$
\min \left\{c^{\prime} x: x \in \mathcal{P}_{k}\right\}
$$

If it does not admit a feasible solution, the same is true for Problem I. 1 - stop. Otherwise, let $x_{k}$ be its optimal solution.
3. If $\mathcal{A}\left(x_{k}\right) \geq 0$ then $x_{k}$ solves Problem I. 1 - stop.
4. Determine a separating hyperplane $a_{k}^{\prime} x=c_{k}$ which separates $x_{k}$ from the feasible convex set $\{x: \mathcal{A}(x) \geq 0\}$. Define

$$
\mathcal{P}_{k+1}:=\mathcal{P}_{k} \cap\left\{x: a_{k}^{\prime} x \geq c_{k}\right\}
$$

set the iteration index $k \leftarrow k+1$ and go back to step 2.
When the algorithm stops the global optimal solution (if any) of Problem I. 1 is provided or it is answered that it is infeasible.

Several points have to be analyzed in details. The first one concerns the determination of the polyhedral set $\mathcal{P}_{0}$. It always exists from the boundedness assumption on the feasible set of Problem I. 1 introduced before. For instance, it can be taken as

$$
\mathcal{P}_{0}:=\left\{x:\left|x_{i}\right| \leq \rho, \quad 1=1,2, \cdots, n\right\}
$$

where $\rho>0$ is a scalar sufficiently large.
In step 2, it is more efficient to solve the stated linear programming problem by means of a dual method. Doing this, it is possible to take advantage from the fact that the sets $\mathcal{P}_{k+1}$ and $\mathcal{P}_{k}$ differ one to the other by only one new linear constraint added in step 4 . Indeed, in iteration $k+1$ a feasible dual solution is readily given by $\lambda_{k+1}^{\prime}:=\left[\lambda_{k}^{\prime} 0\right]$, where $\lambda_{k}$ is the optimal dual solution already obtained in the previous iteration.

In practice, the stopping condition in step 3 has to be changed to $\mathcal{A}\left(x_{k}\right) \geq-\epsilon$ where $\epsilon>0$ is an arbitrarily small parameter. This is necessary because,

$$
\begin{equation*}
\{x: \mathcal{A}(x) \geq 0\} \subset \cdots \subset \mathcal{P}_{k+1} \subset \mathcal{P}_{k}, \quad \forall k=0,1, \cdots \tag{I.2}
\end{equation*}
$$

implies that the algorithm evolves from the outside of the feasible set. Of course, the number of iterations to reach a certain precision increases as $\epsilon$ decreases.

Finally, it must be clear how to calculate the separating hyperplane needed in step 4. We claim that the vector $a_{k}$ and the scalar $c_{k}$ can be taken as

$$
\begin{aligned}
a_{k} & :=\left[\begin{array}{lll}
z_{k}^{\prime} A_{1} z_{k} & \cdots & z_{k}^{\prime} A_{n} z_{k}
\end{array}\right]^{\prime} \\
c_{k} & :=-z_{k}^{\prime} A_{0} z_{k}
\end{aligned}
$$

where $z_{k} \in R^{m}$ is any unitary norm eigenvector associated to the minimum eigenvalue of $\mathcal{A}\left(x_{k}\right)$. Actually, this important property follows immediately from the equality

$$
a_{k}^{\prime} x-c_{k}=z_{k}^{\prime} \mathcal{A}(x) z_{k}
$$

valid for all $x \in R^{n}$. Indeed, for all feasible vectors that is for those vectors such that $\mathcal{A}(x) \geq 0$ obviously $a_{k}^{\prime} x-c_{k} \geq 0$. On the other hand, for $x=x_{k}$ we get

$$
a_{k}^{\prime} x_{k}-c_{k}=\lambda_{\min }\left[\mathcal{A}\left(x_{k}\right)\right]<0
$$

since by construction the point $x_{k}$ is not feasible as it has been previously tested in step 3. It is interesting to observe that if the feasible set of Problem I. 1 is not empty, then $\left\|a_{k}\right\| \neq 0$. This is true because the function $g(x):=\lambda_{\text {min }}[\mathcal{A}(x)]$ is concave and so the inequality

$$
\begin{aligned}
g(x) & \leq z_{k}^{\prime} \mathcal{A}(x) z_{k} \\
& \leq a_{k}^{\prime} x-c_{k} \\
& \leq g\left(x_{k}\right)+a_{k}^{\prime}\left(x-x_{k}\right), \forall x \in R^{n}
\end{aligned}
$$

implies that $a_{k} \in \partial g\left(x_{k}\right)$ and $a_{k}=0$ is possible if and only if $x_{k}$ maximizes $g(\cdot)$. Since $g\left(x_{k}\right)<0$ then it is impossible to have any other $x \in R^{n}$ such that $g(x) \geq 0$. So, $\left\|a_{k}\right\|=0$ may occur if and only if the feasible set of Problem I. 1 is empty.

It remains to prove the global convergence of the Separating hyperplane algorithm which is done in the next lemma.

Lemma I.1 Suppose the feasible set of Problem (I.1) is not empty and the Separating hyperplane algorithm generates the sequence of points $\left\{x_{k}\right\}$. Any limit point of this sequence is a global optimal solution to Problem (I.1).

Proof In the $k$-th iteration of the Separating hyperplane algorithm, the linear constraint (cut) $a_{k}^{\prime} x \geq c_{k}$ has been added to the linear programming problem to be solved in step 2. Consequently in any subsequent iteration, say $l>k$, it must be verified for $x=x_{l}$, that is $a_{k}^{\prime} x_{l} \geq c_{k}$ which yields

$$
\begin{aligned}
0 & \leq a_{k}^{\prime} x_{l}-c_{k} \\
& \leq\left(a_{k}^{\prime} x_{k}-c_{k}\right)+a_{k}^{\prime}\left(x_{l}-x_{k}\right) \\
& \leq \lambda_{\min }\left[\mathcal{A}\left(x_{k}\right)\right]+a_{k}^{\prime}\left(x_{l}-x_{k}\right)
\end{aligned}
$$

or equivalently

$$
\lambda_{\min }\left[\mathcal{A}\left(x_{k}\right)\right] \geq a_{k}^{\prime}\left(x_{k}-x_{l}\right) \geq-\left\|a_{k}\right\|\left\|x_{k}-x_{l}\right\|
$$

Now, since $\left\|a_{k}\right\|<\infty$ in all iterations, the conclusion is that as $k$ and $l$ go to infinity, the algorithm provides a vector $x$ such that $\mathcal{A}(x) \geq 0$, that is $x$ is a feasible solution to Problem (I.1). On the other hand, due to (I.2), in any iteration we have

$$
c^{\prime} x_{k} \leq \min \left\{c^{\prime} x: \mathcal{A}(x) \geq 0\right\}
$$

which shows, by continuity, that the limit point $x$ is optimal indeed.
The Separating hyperplane algorithm has some important features. First, unlike many nonlinear programming methods, it does not require any line search. The practical implementation is simple and depends basically on the development of a powerful dual-simplex routine to solve the linear programming problem on step 2. The use of a dual method also provides useful informations to drop non binding constraints, which of course, keeps reduced the number of constraints to be handled and so increases numerical efficiency. Unfortunately, the Separating hyperplane algorithm as presented
before, may have a poor rate of convergence. Indeed, it can be estimated that it converges arithmetically or at most geometrically with a ratio that goes to unity as the dimension of the problem increases. More specifically

$$
\left\|x_{k}-x_{o p t}\right\| \leq \alpha \beta^{k}
$$

for some constants $\alpha>0$ and $0<\beta<1$ where $\beta=\beta(n)$ goes to one as $n$ increases. Based on this, it is predicted (and practically verified) that it does not perform well for solving large scale problems.

It is also possible to recognize that the determination of the separating hyperplane plays a central role in the numerical efficiency of the method. The deepness of the cut added in each iteration, appears to be of great importance. This leads to the introduction of the so called Supporting hyperplane algorithm which works with the deepest cut that can be calculated in each iteration. For that we assume the feasible set of Problem (I.1) is not empty and we have previously calculated an interior point $y \in R^{n}$ such that $\mathcal{A}(y)>0$. At any infeasible point $x_{k}$, it is possible to calculate a step size $0<\alpha_{k}<1$ which defines a point in the line $d_{k}:=y-x_{k}$ and, in the same time, is on the boundary of the constraint $\mathcal{A}(x) \geq 0$. To this end, we have to find the scalar $\alpha_{k}$ solution to the nonlinear equation (in the unknown $\alpha$ )

$$
\begin{aligned}
0 & =\lambda_{\min }\left[\mathcal{A}\left(y-\alpha d_{k}\right)\right] \\
& =\lambda_{\min }\left[\mathcal{A}(y)-\alpha\left(\mathcal{A}\left(d_{k}\right)-\mathcal{A}(0)\right)\right]
\end{aligned}
$$

which after simple algebraic manipulations provides

$$
\alpha_{k}:=\frac{1}{\mu_{k}}
$$

where

$$
\mu_{k}:=\lambda_{\max }\left[\mathcal{A}(y)^{-1 / 2}\left(\mathcal{A}\left(d_{k}\right)-\mathcal{A}(0)\right) \mathcal{A}(y)^{-1 / 2}\right]
$$

All operations indicated above are well defined since $\mathcal{A}(y)>0$ implies that $\mathcal{A}(y)$ is nonsingular. Furthermore, from the above, $\mu_{k} \leq 0$ implies that

$$
\begin{aligned}
0 & \geq \mathcal{A}\left(d_{k}\right)-\mathcal{A}(0) \\
& \geq \mathcal{A}(y)-\mathcal{A}\left(x_{k}\right)
\end{aligned}
$$

that is $\mathcal{A}\left(x_{k}\right) \geq \mathcal{A}(y)>0$ which puts in evidence that for any infeasible vector $x_{k}$ we necessarily have $\mu_{k}>0$. The point

$$
\bar{x}_{k}:=y-\alpha_{k} d_{k}
$$

is, by construction, on the boundary of the feasible set of Problem (I.1) and for its determination only an additional eigenvalue calculation is needed. The Supporting hyperplane algorithm is stated in the sequel. The main idea is to use the above information to define iteratively a supporting hyperplane at the boundary point $\bar{x}_{k}$.
Algorithm I. 2 (Supporting hyperplane algorithm) It is completely based on the previous algorithm where the separating hyperplane is replaced by a supporting hyperplane. Although more efficient, the limitation to handle large scale problems remains.

1. Determine an interior feasible point $y \in R^{n}$, a convex polyhedral set $\mathcal{P}_{0}$ containing the overall feasible set of Problem I.1, that is

$$
\{x: \mathcal{A}(x) \geq 0\} \subset \mathcal{P}_{0}
$$

and set the iteration index $k=0$.
2. Solve the linear programming problem

$$
\min \left\{c^{\prime} x: x \in \mathcal{P}_{k}\right\}
$$

and let $x_{k}$ be its optimal solution.
3. If $\mathcal{A}\left(x_{k}\right) \geq 0$ then $x_{k}$ solves Problem I. 1 - stop. Otherwise, calculate the point $\bar{x}_{k}$ on the boundary of the feasible set and go to the next step.
4. Determine a supporting hyperplane $a_{k}^{\prime} x=c_{k}$ at the boundary point $\bar{x}_{k}$. Define

$$
\mathcal{P}_{k+1}:=\mathcal{P}_{k} \cap\left\{x: a_{k}^{\prime} x \geq c_{k}\right\}
$$

set the iteration index $k \leftarrow k+1$ and go back to step 2.
Once again, when the algorithm stops the global optimal solution of Problem I. 1 is provided.

Let us now discuss an important class of convex programming methods, called Interior point methods. These methods apply to the solution of an approximate version of Problem (I.1) given in the form

$$
\begin{equation*}
\inf \left\{c^{\prime} x: \mathcal{A}(x)>0\right\} \tag{I.3}
\end{equation*}
$$

Clearly, this problem is equivalent to Problem (I.1) provided the LMI $\mathcal{A}(x) \geq 0$ admits an interior point, that is a vector $x \in R^{n}$ such that $\mathcal{A}(x)>0$. In this case, the equivalence between problems (I.1) and (I.3) holds in the sense that their optimal solutions are arbitrarily close one to the other. In the developments that follow, we work with Problem (I.3).

The main idea comes from the definition of the analytic center of a LMI. The analytic center of the LMI

$$
\mathcal{A}(x)=A_{0}+\sum_{i=1}^{n} x_{i} A_{i}>0
$$

is the vector $x_{a c} \in R^{n}$ such that

$$
\begin{equation*}
x_{a c}:=\operatorname{argmin}\{-\log \operatorname{det}[\mathcal{A}(x)]: \mathcal{A}(x)>0\} \tag{I.4}
\end{equation*}
$$

The objective function of the above problem can be interpreted as a barrier function for the LMI under consideration. Indeed, as $x$ goes to the boundary of the feasible set, at least one eigenvalue of $\mathcal{A}(x)$ goes to zero and enforces the objective function to be arbitrarily large. Moreover, the following properties are of great importance in the calculations that follow.

Lemma I. 2 The function $p(x):=-\log \operatorname{det}[\mathcal{A}(x)]$, defined in the open convex set $\mathcal{A}(x)>0$ is such that :
i) Function $p(x)$ is convex.
ii) At any $x \mid \mathcal{A}(x)>0$, the gradient of $p(x)$ is

$$
\nabla p(x)_{i}:=-\operatorname{trace}\left[\mathcal{A}(x)^{-1} A_{i}\right], i=1,2, \cdots, n
$$

iii) At any $x \mid \mathcal{A}(x)>0$, the Hessian matrix of $p(x)$ is

$$
H(x)_{i j}:=\operatorname{trace}\left[\mathcal{A}(x)^{-1} A_{i} \mathcal{A}(x)^{-1} A_{j}\right], i, j=1,2, \cdots, n
$$

Proof The proof is based on the concavity of the scalar function $\log (z)$ in the interval $z>0$ which implies that $\log (z) \leq z-1$ for all $z>0$. Using this and any two vectors such that $\mathcal{A}(x)>0$ and $\mathcal{A}(y)>0$, we get

$$
\begin{aligned}
p(y)-p(x) & =-\log \operatorname{det}\left[\mathcal{A}(x)^{-1} \mathcal{A}(y)\right] \\
& =-\sum_{l=1}^{m} \log \lambda_{l}\left[\mathcal{A}(x)^{-1} \mathcal{A}(y)\right] \\
& \geq-\sum_{l=1}^{m}\left\{\lambda_{l}\left[\mathcal{A}(x)^{-1} \mathcal{A}(y)\right]-1\right\}
\end{aligned}
$$

and consequently

$$
\begin{aligned}
p(y) & \geq p(x)-\operatorname{trace}\left[\mathcal{A}(x)^{-1} \mathcal{A}(y)-I\right] \\
& \geq p(x)-\operatorname{trace}\left[\mathcal{A}(x)^{-1} \mathcal{A}(y)-\mathcal{A}(x)^{-1} \mathcal{A}(x)\right] \\
& \geq p(x)-\sum_{i=1}^{n} \operatorname{trace}\left[\mathcal{A}(x)^{-1} A_{i}\right]\left(y_{i}-x_{i}\right) \\
& \geq p(x)+\nabla p(x)^{\prime}(y-x)
\end{aligned}
$$

which proves the first two points of the lemma proposed. The last one is proved by simple partial differentiation of $\nabla p(x)_{i}$ with respect to the variable $x_{j}$. The proof is concluded.

Example I. 2 For the same LMI of example I.1, figure I. 2 shows the level set of $\operatorname{det}[\mathcal{A}(x)]=$ $\alpha>0$ which for $\mathcal{A}(x)>0$ and $\beta=-\log (\alpha)$ coincides with that of $p(x)=\beta$. It is clearly seen that in the interior of the LMI the level set for each value of $\alpha>0$ defines a convex set. Moreover, the closed region approaches to the boundary of the LMI as $\alpha$ goes to zero. Outside this region, there exist points for which the determinant of the affine function $\mathcal{A}(x)$ attains the same level but with an even number of negative eigenvalues. Finally, using part, ii) of Lemma I. 2 we solve $\nabla p(x)=0$ to get the analytic center $x_{c a}=[0.59020 .4236]$.

Further inspection reveals that function $p(x)$ is in fact strictly convex which means that the Hessian matrix $H(x)$ is positive definite whenever the vector $x \in R^{n}$ is such that $\mathcal{A}(x)>0$. Hence, the analytic center of the LMI can very efficiently be calculated by the following well known Modified Newton's method.

Algorithm I. 3 (Modified Newton's method) Assume an initial point $x_{0}$ such that $\mathcal{A}\left(x_{0}\right)>0$ is given. Then, perform the following iterations until convergence.

1. Determine the gradient vector $\nabla p\left(x_{k}\right)$ and the Hessian matrix $H\left(x_{k}\right)$.
2. Determine the descent direction $d_{k}:=H\left(x_{k}\right)^{-1} \nabla p\left(x_{k}\right)$. If within some prespecified precision $\left\|d_{k}\right\|=0$ - stop. Otherwise, go to the next step.
3. Determine the optimal step length $\alpha_{k}$ given by

$$
\alpha_{k}:=\operatorname{argmin} p\left(x_{k}-\alpha d_{k}\right)
$$



Figure I.2: The level set of $\operatorname{det}[\mathcal{A}(x)]$
3. Update $x_{k+1}=x_{k}-\alpha_{k} d_{k}$, set the iteration index $k \leftarrow k+1$ and go back to step 1.

When the algorithm stops the analytic center is approximately given by $x_{a c}=x_{k}$.
For the complete implementation of this algorithm, it remains to calculate the optimal step length $\alpha_{k}$ defined in step 3 . This is accomplished with no great difficulty. Indeed, simple calculations put in evidence that

$$
\begin{equation*}
p\left(x_{k}-\alpha d_{k}\right)=p\left(x_{k}\right)-\sum_{l=1}^{m} \log \left[1-\alpha e_{k l}\right] \tag{I.5}
\end{equation*}
$$

where

$$
e_{k l}:=\lambda_{l}\left[\mathcal{A}\left(x_{k}\right)^{-1 / 2}\left(\mathcal{A}\left(d_{k}\right)-\mathcal{A}(0)\right) \mathcal{A}\left(x_{k}\right)^{-1 / 2}\right]
$$

is the $l$-th component of the $m$ dimensional real vector $e_{k}:=\left[\begin{array}{llll}e_{k 1} & e_{k 2} & \cdots & e_{k m}\end{array}\right]^{\prime}$. Hence, the derivative of $p\left(x_{k}-\alpha d_{k}\right)$ with respect to $\alpha$ provides

$$
\begin{aligned}
\frac{d p}{d \alpha}\left(x_{k}-\alpha d_{k}\right) & =\sum_{l=1}^{m} \frac{e_{k l}}{1-\alpha e_{k l}} \\
& =\sum_{l=1}^{m}\left[\frac{e_{k l}^{2}}{\beta-e_{k l}}+e_{k l}\right]
\end{aligned}
$$

where $\beta:=1 / \alpha$. Setting the right hand side of the above equation to zero and taking into account that

$$
\delta_{k}^{2}:=-\sum_{l=1}^{m} e_{k l}=-\left.\frac{d p}{d \alpha}\left(x_{k}-\alpha d_{k}\right)\right|_{\alpha=0}>0
$$

since, $d_{k}$ is a nonzero descent direction, the optimal step size $\alpha=1 / \beta$ solves the nonlinear equation

$$
\sum_{l=1}^{m}\left[\frac{e_{k l}^{2}}{\beta-e_{k l}}\right]-\delta_{k}^{2}=0
$$

The solution of this nonlinear equation is not simple to be determined unless we realize it can be rewritten in terms of the following determinantal equation

$$
\operatorname{det}\left[\delta_{k}^{2}-e_{k}^{\prime}\left(\beta I-E_{k}\right)^{-1} e_{k}\right]=0
$$

where matrix $E_{k} \in R^{m \times m}$ is defined as $E_{k}:=\operatorname{diag}\left[e_{k 1}, e_{k 2}, \cdots, e_{k m}\right]$. Finally, the last equation together with some elementary determinant manipulations provides

$$
\operatorname{det}\left[\beta I-\left(E_{k}+\delta_{k}^{-2} e_{k} e_{k}^{\prime}\right)\right]=0
$$

which makes clear that the optimal step size $\alpha_{k}$ is given by

$$
\begin{equation*}
\alpha_{k}=\lambda_{\max }^{-1}\left[E_{k}+\left(\frac{e_{k}}{\delta_{k}}\right)\left(\frac{e_{k}}{\delta_{k}}\right)^{\prime}\right] \tag{I.6}
\end{equation*}
$$

This shows that to determine the optimal steep size, we have to calculate all eigenvalues of a symmetric matrix in order to define the vector $e_{k}$ and finally to calculate the maximum eigenvalue of the symmetric matrix indicated in (I.6). To reduce this amount of calculations, in some cases, we have a great advantage to get a suboptimal step size as indicated in the sequel. It comes to light from the observation that any positive step size less than the optimal one may also be used to assure that function $p(x)$ is reduced in the direction $-d_{k}$. To get such a suboptimal step size, we proceed by establishing the following equality

$$
\begin{aligned}
\sum_{l=1}^{m} e_{k l}^{2} & =\left.\frac{d^{2} p}{d \alpha^{2}}\left(x_{k}-\alpha d_{k}\right)\right|_{\alpha=0} \\
& =d_{k}^{\prime} H\left(x_{k}\right) d_{k} \\
& =\nabla p\left(x_{k}\right)^{\prime} d_{k} \\
& =-\left.\frac{d p}{d \alpha}\left(x_{k}-\alpha d_{k}\right)\right|_{\alpha=0} \\
& =\delta_{k}^{2}
\end{aligned}
$$

which together with (I.6) implies that

$$
\begin{aligned}
\alpha_{k}^{-1} & =\lambda_{\max }\left[E_{k}+\left(\frac{e_{k}}{\delta_{k}}\right)\left(\frac{e_{k}}{\delta_{k}}\right)^{\prime}\right] \\
& \leq \lambda_{\max }\left[E_{k}\right]+\delta_{k}^{-2} \sum_{l=1}^{m} e_{k l}^{2} \\
& \leq 1+\lambda_{\max }\left[E_{k}\right]
\end{aligned}
$$

consequently, a suboptimal step size, denoted as $\alpha_{k}^{+}$is given by

$$
\alpha_{k}^{+}:=\frac{1}{1+\mu_{k}}
$$

where

$$
\mu_{k}:=\lambda_{\max }\left[\mathcal{A}\left(x_{k}\right)^{-1 / 2}\left(\mathcal{A}\left(d_{k}\right)-\mathcal{A}(0)\right) \mathcal{A}\left(x_{k}\right)^{-1 / 2}\right]
$$

It is to be noticed that $\mu_{k}>0$ since otherwise all $e_{k l} \leq 0$ which implies from (I.5) that the feasible set is unbounded, a situation avoided by our previous assumption.

It is also interesting to see that the above formula for the suboptimal step size is very similar to the one introduced before for the calculation of a point on the boundary of the feasible set of the LMI under consideration.

Let us now use the concept of analytic center to calculate the optimal solution of Problem (I.3). Obviously it can be equivalently stated in the form

$$
\inf \left\{\gamma: \mathcal{A}(x)>0, \quad \gamma-c^{\prime} x>0\right\}
$$

or, in terms of only one augmented LMI

$$
\begin{equation*}
\inf \{\gamma: \mathcal{B}(x, \gamma)>0\} \tag{I.7}
\end{equation*}
$$

where

$$
\mathcal{B}(x, \gamma):=\left[\begin{array}{cc}
\mathcal{A}(x) & 0 \\
0 & \gamma-c^{\prime} x
\end{array}\right]
$$

This LMI depends on both variables namely the vector $x \in R^{n}$ and the scalar $\gamma$. However, for $\gamma$ fixed, let us define as before the analytic center

$$
x_{a c}(\gamma):=\operatorname{argmin}\{-\log \operatorname{det}[\mathcal{B}(x, \gamma)]: \mathcal{B}(x, \gamma)>0\}
$$

where it is indicated the dependence of the analytic center with respect to the scalar $\gamma$ and that the minimization must be done with respect to $x \in R^{n}$ only. The curve $x_{a c}(\gamma)$ obtained for all possible values of $\gamma$ is called the Path of centers and plays a central role to the numerical solution of Problem (I.3) as indicated in the next algorithm.

Algorithm I. 4 (Method of centers) Assume an initial pair ( $x_{0}, \gamma_{0}$ ) is given, such that simultaneously $\mathcal{A}\left(x_{0}\right)>0$ and $\gamma_{0}>c^{\prime} x_{0}$. Choose $0<\theta<1$ and $\epsilon>0$ sufficiently small and perform the following iterations until convergence.

1. $\gamma_{k+1}=(1-\theta) c^{\prime} x_{k}+\theta \gamma_{k}$
2. $x_{k+1}=x_{a c}\left(\gamma_{k+1}\right)$
3. If $\gamma_{k+1}-c^{\prime} x_{k+1}<\epsilon / m$ - stop. Otherwise set the iteration index $k \leftarrow k+1$ and go back to step 1 .

When the algorithm stops the optimal solution to Problem (I.7) is found within $\epsilon$.
It is important to recognize that the rule in step 1 , never produces infeasibility on the analytic center determination in step 2. Actually, assume that in a generic iteration $k \geq 0$ we have $\mathcal{B}\left(x_{k}, \gamma_{k}\right)>0$. With the formula stated in step 1 , we get

$$
\gamma_{k+1}-c^{\prime} x_{k}=\theta\left(\gamma_{k}-c^{\prime} x_{k}\right)>0
$$

which implies that $\mathcal{B}\left(x_{k}, \gamma_{k+1}\right)>0$ and consequently the vector $x_{k}$ can be used to initialize the Modified Newton's method for the determination of the analytic center $x_{a c}\left(\gamma_{k+1}\right)$. In practice, it is verified that this simple initialization procedure is very effective as far as numerical efficiency is concerned.

Lemma I. 3 The Method of centers converges geometrically to the optimal solution of Problem (I.3).

Proof Denote ( $x_{o p t}, \gamma_{o p t}$ ) the optimal solution of Problem (I.3). For $\gamma=\gamma_{k+1}$ fixed, Lemma I. 2 enables us to write the optimality conditions to characterize the analytic center $x_{a c}\left(\gamma_{k+1}\right)$. So, due to step 2 we must have

$$
\frac{c_{i}}{\gamma_{k+1}-c^{\prime} x_{k+1}}=\operatorname{trace}\left[\mathcal{A}\left(x_{k+1}\right)^{-1} A_{i}\right], \quad i=1,2, \cdots, n
$$

which gives (recall that $\gamma_{o p t}=c^{\prime} x_{o p t}$ )

$$
\frac{c^{\prime} x_{k+1}-\gamma_{o p t}}{\gamma_{k+1}-c^{\prime} x_{k+1}}=\operatorname{trace}\left[\mathcal{A}\left(x_{k+1}\right)^{-1}\left(\mathcal{A}\left(x_{k+1}\right)-\mathcal{A}\left(x_{o p t}\right)\right)\right]
$$

Now, define the scalar $\phi$ as being

$$
\phi:=\sup _{\mathcal{A}(x)>0} \operatorname{trace}\left[\mathcal{A}(x)^{-1}\left(\mathcal{A}(x)-\mathcal{A}\left(x_{o p t}\right)\right)\right]
$$

and observe that $0 \leq \phi \leq m$. Actually, the lower bound is obtained from the simple observation that $x=x_{\text {opt }}$ is feasible and the upper bound is a consequence of the fact that $\mathcal{A}\left(x_{\text {opt }}\right)>0$. Then, the inequality

$$
\phi \gamma_{k+1}+\gamma_{o p t} \geq(1+\phi) c^{\prime} x_{k+1}
$$

holds in all iterations. Using the update of step 1, namely

$$
c^{\prime} x_{k+1}=\frac{\gamma_{k+2}-\theta \gamma_{k+1}}{1-\theta}
$$

simple algebraic manipulations put in evidence that

$$
\gamma_{k+2}-\gamma_{o p t} \leq \frac{\phi+\theta}{\phi+1}\left(\gamma_{k+1}-\gamma_{o p t}\right)
$$

which proves that the Method of centers converges geometrically. This concludes the proof of the Lemma proposed.

This proof is of great practical importance for two main reasons. First, if the stopping criterion in step 3 is verified then

$$
\begin{aligned}
c^{\prime} x_{k+1}-\gamma_{o p t} & \leq \phi\left(\gamma_{k+1}-c^{\prime} x_{k+1}\right) \\
& \leq m(\epsilon / m)=\epsilon
\end{aligned}
$$

and the optimal solution is found within the prespecified precision level $\epsilon>0$ imposed by the designer. Second, the ratio of geometric convergence, such that

$$
0 \leq c^{\prime} x_{k}-c^{\prime} x_{o p t} \leq \alpha \beta(\phi)^{k}
$$

for some $\alpha>0$, is estimated as being

$$
\beta(\phi):=\frac{\phi+\theta}{\phi+1}
$$

which is an increasing function of $\phi$. The worst estimation is then obtained for $\phi=m$ providing thus $\beta(m)$. It is important to realize that, doing this, the conclusion is that the Method of centers converges geometrically but with a ratio that goes to unity


Figure I.3: Convergence behavior
as the dimension of the problem to be solved increases. In other words, it performs, under this worst case analysis, as the Separating hyperplane algorithm. However, it is possible to introduce in the Method of centers a simple modification in order to get much better convergence behavior. Indeed, if in the determination of the analytic center $x_{a c}(\gamma)$ the objective function is changed to

$$
-\log \operatorname{det}[\mathcal{A}(x)]-m \log \left(\gamma-c^{\prime} x\right)
$$

which is nothing more than to redefine the augmented LMI by replacing the scalar $\gamma-c^{\prime} x$ by the $m \times m$ diagonal matrix $\left(\gamma-c^{\prime} x\right) I$, then the same reasoning used in the proof of Lemma I. 3 yields the new estimate for the ratio of geometric convergence

$$
\begin{equation*}
\beta(\phi):=\frac{\phi+\theta m}{\phi+m} \leq \frac{1+\theta}{2}<1 \tag{I.8}
\end{equation*}
$$

The worst case for the ratio of convergence is now independent of the problem dimension. It depends only on the parameter $0<\theta<1$ to be fixed by the designer.
Example I. 3 Consider the LMI of example I. 1 and $c^{\prime}=\left[\begin{array}{ll}0 & 1\end{array}\right]$. The optimal solution of problem (I.1) is found to be $x_{o p t}=[0.3132-0.9018]^{\prime}$ and $\gamma_{o p t}=c^{\prime} x_{o p t}=-0.9018$. Figure I. 3 shows the objective function per iteration calculated by Algorithm I. 2 (dashed line) and by Algorithm I. 4 (solid line). As expected, the first evolves through infeasible points and the objective function is increasing. On the contrary, the second is always feasible and the objective function is decreasing.

Finally it is also to be noticed that the determination of a feasible starting point required for many algorithms can be done with no great difficulty. Given the LMI $\mathcal{A}(x) \geq 0$, the problem is to find (if one exists) a vector $x$ such that $\mathcal{A}(x)>0$. This is accomplished by solving the auxiliary convex problem

$$
\min \{\lambda: \mathcal{A}(x)+\lambda I \geq 0\}
$$

which presents two interesting properties. First, the pair $(x, \lambda):=\left(0,1-\lambda_{\min }\left[A_{0}\right]\right)$ satisfies the LMI constraint strictly and second at the optimal solution $\mathcal{A}\left(x_{o}\right)>0$ if and only if $\lambda_{o}<0$. Of course, the optimal solution of the above problem does not need to be exactly calculated. If in some iteration the current value of the auxiliary variable $\lambda$ becomes negative, the current value of the vector $x$ is strictly feasible and the search process may be stopped.

## Bibliography

[1] J. Ackermann. Robust Control. Springer Verlag, 1993.
[2] T. Basar and P. Bernhard. $H_{\infty}$ Optimal Control and Related Minimax Design Problems. Birkhäuser, 1991.
[3] B.R. Barmish. Necessary and Sufficient Conditions for Quadratic Stabilizability of an Uncertain System. J. Opt. Th. Appl. 46, pp.399-408, 1985.
[4] B.R. Barmish. Stabilization of Uncertain Systems Via Linear Control. IEEE Trans. Auto. Contr. AC-28, pp.848-850, 1983.
[5] D.S. Bernstein and W.M. Haddad. LQG Control with an $H_{\infty}$ Performance Bound: A Riccati Equation Approach. IEEE Trans. Auto. Contr. AC-34, pp.293305, 1989.
[6] J. Bernussou, P.L.D. Peres and J.C. Geromel. A Linear Oriented Procedure for Quadratic Stabilization of Uncertain Systems. Syst. Contr. Lett. 13, pp.65-72, 1989.
[7] S. Bittanti, A.J. Laub and J.C. Willems (Eds). The Riccati Equations. Springer Verlag, 1991.
[8] P. Bolzern, P. Colaneri and G. De Nicolao. On the Computation of Upper Covariance Bounds for Perturbed Linear Systems. IEEE Trans. Auto. Contr. AC-39, pp.623-626, 1994.
[9] S.P. Boyd, V. Balakrishnan and P.T. Kabamba. A Bisection Method for Computing the $H_{\infty}$ Norm of a Transfer Matrix and Related Problems. Mathematics of Control, Signals, and Syst. 2, pp.207-219, 1989.
[10] S.P. Boyd, L. El Ghaoui, E. Feron and V. Balakrishnan. Linear Matrix Inequalities in System and Control Theory. Siam Studies in Applied Mathematics, 1994.
[11] S.S.L. Chang and T.K.C. Peng. Adaptive Guaranteed Cost Control of Systems with Uncertain Parameters. IEEE Trans. Auto. Contr. AC-17, pp.474-483, 1972.
[12] C.E. de Souza, U. Shaked and N. Fu. Robust $H_{\infty}$ Filtering with Parametric Uncertainty and Deterministic Input Signal. Proc. $31^{\text {st }}$ Conf. Dec. Contr. Tucson, USA, pp.2305-2310,1992.
[13] J.C. Doyle. Robustness of Multiloop Linear Feedback Systems. Proc. $17^{\text {th }}$ Conf. on Dec. and Contr. USA, pp.12-18, 1978.
[14] J.C. Doyle. Analysis of Feedback Systems with Structured Uncertainties. IEE Proc. Pt.D. 129, pp.242-250, 1982.
[15] J.C. Doyle. Structured Uncertainty in Control System Design. Proc. $24^{\text {th }}$ Conf. on Dec. and Contr. Ft. Lauderdale, USA, pp.260-265, 1985.
[16] J.C. Doyle, B.A. Francis, A.M. Tannenbaum. Feedback Control Theory. Maxwell MacMillan Int. Ed., 1992.
[17] J.C. Doyle, K. Glover, P.P. Khargonekar and B.A. Francis. State Space Solutions to Standard $H_{2}$ and $H_{\infty}$ Control Problems. IEEE Trans. Auto. Contr. AC-34, pp.831-847, 1989.
[18] J.C. Doyle, J.E. Wall and G. Stein. Performance and Robustness Analysis for Structured Uncertainty. Proc. $21^{\text {st }}$ Conf. on Dec. and Contr. pp.629-636, 1982.
[19] B.A. Francis. A Course in $H_{\infty}$ Control Theory. Lectures Notes on Control and Information Sciences. Springer Verlag, 1987.
[20] N. Fu, C.E. de Souza and L. Xie. $H_{\infty}$ Estimation for Uncertain Systems. Int. J. Rob. and Nonlin. Contr. 2, pp.82-105, 1992.
[21] T.T. Georgiou and M.C. Smith. Optimal Robustness in the Gap Metric: Controller Design for Distributed Plants. IEEE Trans. Auto. Contr. AC-37, pp.11331143, 1992.
[22] J.C. Geromel, P.L.D. Peres and J. Bernussou. On a Convex Parameter Space Method for Linear Control Design of Uncertain Systems. SIAM J. Contr. and Opt. 29, pp.381-402, 1991.
[23] J.C. Geromel, J. Bernussou and P.L.D. Peres. Decentralized Control Through Parameter Space Optimization. Automatica 30, pp.1565-1578, 1994.
[24] J.C. Geromel, P.L.D. Peres and S.R. Souza. $\mathcal{H}_{2}$ Guaranteed Cost Control for Uncertain Continuous-time Linear Systems. Syst. Contr. Lett. 19, pp.23-27, 1992.
[25] J.C. Geromel. On the Determination of a Diagonal Solution of the Lyapunov Equation. IEEE Trans. Auto. Contr. AC-30, pp.404-406, 1985.
[26] J.C. Geromel and A.O.E. Santo. On the Robustness of Linear Continuous-time Dynamic Systems. IEEE Trans. Auto. Contr. AC-31, pp.1136-1138, 1986.
[27] W.M. Haddad and D.S. Bernstein. Robust Stabilization with Positive Real Uncertainty: Beyond the Small Gain Theorem. Syst. Contr. Lett. 17, pp.191-208, 1991.
[28] W.M. Haddad and D.S. Bernstein. Controller Design with Regional Pole Constraints. IEEE Trans. Auto. Contr. AC-37, pp.54-69, 1992.
[29] T. Kailath. Linear Systems. Prentice Hall, 1980.
[30] E. Kaszkurewicz and A. Bhaya. Robust Stability and Diagonal Liapunov Functions. SIAM J. Matrix Anal. Appl. 14, pp.508-520, 1993.
[31] P.P. Khargonekar, I.R. Petersen and K. Zhou. Robust Stabilization of Uncertain Linear Systems: Quadratic Stabilizability and $H_{\infty}$ Control Theory. IEEE Trans. Auto. Contr. AC-35, pp.356-361, 1990.
[32] P.P. Khargonekar and M.A. Rotea. Mixed $H_{2} / H_{\infty}$ Control: A Convex Optimization Approach. IEEE Trans. Auto. Contr. AC-36, pp.824-837, 1991.
[33] B. Kouvarikatis and A.G.J. MacFarlane. Geometric Approach to Analysis and Synthesis of System Zeros: Part I - Square Systems. Int. J. Control. 23, pp.149166, 1976.
[34] B. Kouvarikatis and A.G.J. MacFarlane. Geometric Approach to Analysis and Synthesis of System Zeros: Part II - Nonsquare Systems. Int. J. Control. 23, pp.167-181, 1976.
[35] V. Kucera. The LQG and $H_{2}$ Designs: Two Different Problems? Proc. 2 ${ }^{\text {nd }}$ Europ. Contr. Conf.. Groningen, The Netherlands, pp.334-337, 1988.
[36] H. Kwakernaak. The Polynomial Approach to the $H_{\infty}$ Optimal Regulation. $H_{\infty}$ Control. E. Mosca and L. Pandolfi Eds. Lecture Notes in Mathematics, 1496. Springer-Verlag, 1990.
[37] L.S. Lasdon. Optimization Theory for Large Systems. Macmillan Publishing Co., 1970.
[38] C.L. Lawson and R.J. Hanson. Solving Least-Square Problems. Prentice Hall, 1974.
[39] W.S. Levine and M. Athans. On the Determination of the Optimal Constant Output Feedback Gains for Linear Multivariable Systems. IEEE Trans. Auto. Contr. AC-15, pp.44-48, 1970.
[40] D.J.N. Limebeer, B.D.O. Anderson and B. Hendel. A Nash Game Approach to Mixed $H_{2} / H_{\infty}$ Control. IEEE Trans. Auto. Contr. AC-39, pp.69-82, 1994.
[41] D.G. Luenberger. Linear and Nonlinear Programming. Addison-Wesley Publishing Co., 1984.
[42] A.G.J. MacFarlane and N. Karcanias. Poles and Zeros of Linear Multivariable Systems: a Survey of the Algebraic, Geometric and Complex Variable Theory. Int. J. Control. 24, pp.33-74, 1976.
[43] J.M. Maciejowski. Multivariable Feedback Design. Addison Wesley, 1989.
[44] T. Mita, K. Kuriyamna and K.Z. Liu. Revision of FI Result in $H_{\infty}$ Control and Parametrization of State Feedback Controllers. Trans. Soc. Instrum. Contr. Eng. 28, pp.801-809, 1992.
[45] Y. Nesterov and A. Nemirovsky. Interior-point Polynomial Methods in Convex Programming. Siam Studies in Applied Mathematics, 1994.
[46] S. Parrott. On a Quotient Norm and the Sz.-Nagy-Foias Lifting Theorem. J. Functional Analysis. 30, pp.311-328, 1978.
[47] P.L.D. Peres and J.C. Geromel. An Alternate Numerical Solution to the Linear Quadratic Problem. IEEE Trans. Auto. Contr. AC-39, pp.198-202, 1994.
[48] P.L.D. Peres, J.C. Geromel and S.R. Souza. Optimal $\mathcal{H}_{\infty}$ - State Feedback Control for Continuous-Time Linear Systems. J. of Opt. Theory and App. 82, pp.343359, 1994.
[49] P.L.D. Peres, J.C. Geromel and S.R. Souza. $\mathcal{H}_{\infty}$ Guaranteed Cost Control for Uncertain Continuous-time Linear Systems. Syst. Contr. Lett. 20, pp.413-418, 1993.
[50] P.L.D. Peres, J.C. Geromel and J. Bernussou. Quadratic Stabilizability of Linear Uncertain Systems in Convex-bounded Domains. Automatica 29, pp.491-493, 1993.
[51] I.R. Petersen and D.C. McFarlane. Robust State Estimation for Uncertain Systems. Proc. $30^{\text {th }}$ Conf. Dec. Contr. Brighton, England, pp.2630-2631, 1991.
[52] S.C. Power. Hankel Operator in Hilbert Space. Pitman, 1982.
[53] R.T. Rockafellar. Convex Analysis. Princeton Univ. Press, 1970.
[54] M.A. Rotea. The Generalized $H_{2}$ Control Problem. Automatica 29, pp373-386, 1993.
[55] W. Rudin. Functional Analysis. McGraw Hill, 1973.
[56] M.G. Safonov and D.J.N. Limebeer. Simplifying the $H_{\infty}$ Theory via Loop Shifting. Proc. $27^{\text {th }}$ Conf. on Dec. and Contr.. Austin, USA, pp.1399-1404, 1988.
[57] G. Stein and J.C. Doyle. Beyond Singular Values and Loop Shapes. AIAA J. of Guid. Contr. and Dynam. 14, pp.5-16, 1991.
[58] A.A. Stoorvogel. The singular $H_{\infty}$ Control Problem with Dynamic Measurement Feedback. SIAM J. Contr. and Opt. 29, pp.160-184, 1991.
[59] A.A. Stoorvogel. The $H_{\infty}$ Control Problem. Prentice Hall, 1992.
[60] M. Vidyasagar. Control Systems Synthesis. MIT Press, 1985.
[61] M. Vidyasagar. Control Systems Analysis. Prentice Hall, 1993.
[62] D.A. Wilson. Convolution and Hankel Operator Norms for Linear Systems. IEEE Trans. Auto. Contr. AC-34, pp.94-98, 1989.
[63] D.C. Youla, H.A. Jabr and J.J. Bongiorno. Modern Wiener-Hopf Design of Optimal Controllers: Part II. IEEE Trans. Auto. Contr. AC-21, pp.319-338, 1976.
[64] G. Zames. Feedback and Optimal Sensitivity: Model Reference Transformations, Multiplicative Seminorms and Approximate Inverses. IEEE Trans. Auto. Contr. AC-26, pp.301-320, 1981.
[65] K. Zhou and P.P. Khargonekar. Robust Stabilization of Linear Systems with Norm-bounded Time-varying Uncertainty. Syst. Contr. Lett. 10, pp.17-20, 1988.

## Index

## A

Actuators failure, 286
Actuators uncertainty, 124
Additive perturbations, 130, 294
Adjoint Hankel operator, 63
Adjoint Laurent operator, 63
Adjoint operator, 44
Adjoint system, 5
Admissible controller
in $\mathrm{RH}_{2}, 88$
Affine function, 196
Algorithm
interior point, 364
method of centers, 368
separating hyperplane, 360
supporting hyperplane, 363
Augmented system, 126
Auxiliary problem, 211
Auxiliary system, 234, 249

## B

Block-diagonal structure, 251
Blocking property, 32

## C

Computing norm
$R H_{2}, 51$
$R H_{2}^{\perp}, 51$
$R H_{\infty}, 54,56,57,59$
$R L_{2}, 52$
$R L_{\infty}, 52$
Hankel operator, 67
Continuity, 351
Control design
in $R H_{\infty}, 214$
Controllability, 325
Controller
central, $83,218,272$
dynamic, 201
observer-based, 201
stabilizing, 198
structure, 250
Controllers
stabilizing, 80
Convergence, 260
Convex cone, 350
Convex function, 195, 352
Convex hull, 350
Convex optimization, 349
Convex polyhedron, 350
Convex problem, 196
Convex programming
numerical tools, 359
problem, 357
Convex set, 195, 349
bounded, 206
Convexity, 196, 245
Coprime, 10
left, 10, 18
right, 10,18
Covariance bound, 128

## D

D - stability, 252
Damping ratio, 236
Decentralized control exact, 254
Decoupling zeros, 35 input, 36
input-output, 36
output, 35
Descent property, 356
Description
input-output, 27
state-space, 27
transfer function, 27
Detectability, 325
Differentiable function, 196
Directional derivative, 355
Disturbance feedforward, 100
Disturbance feedforward problem, 154
Double coprime factorization, 74

## E

Eigenvalue, 40, 44
Eigenvector, 44
Epigraph, 353
Extended Lyapunov lemma, 315
Extended matrix, 212
Extreme matrix, 265
Extreme point, 350

## F

Factorization inner outer, 47
Feasible starting point, 370
Feedback
decentralized, 250
static output, 250
Feedback connection, 84, 127
First order approximation, 123
Frobenius norm, 349
Full control, 110
Full information, 138
Full information problem, 137, 209, 278
in $\mathrm{RH}_{2}, 90$
Functional space, 45
$R H_{2}$ and $\mathrm{RH}_{2}^{\perp}, 48$
$R H_{\infty}, 46$
$R L_{2}, 47$
$R L_{\infty}, 45$
norm, 45,48

## G

Givens transformation, 304
Global minimization, 242
Gradient, 351
Greatest common left divisor, 7, 16
Greatest common right divisor, 6, 16
Guaranteed cost, 267, 287
$R \mathrm{H}_{2}, 273$
in $R H_{\infty}, 279$
problem, 290, 298

## H

Hankel operator, 61
Householder transformation, 303

## I

Inner function, 46
Inner product, 48
Invariant zeros, 30
changes of basis, 32
rank property, 31
Inverse system, 38

## J

Joint convexity, 225
Joint sensitivity performance, 126

## K

Kalman filter, 109
Kalman test, 325
Kalman-Yacubovitch lemma, 296

## L

Laurent operator, 61
Linear matrix inequalities
analytic center, 364
Linear matrix inequality, 208, 221
Linear operator, 43, 182
boundedness, 43
continuity, 43
norm, 43
rank, 44
Linear quadratic differential game, 151
Loop shifting, 173, 337
LQ problem, 98, 118
LQS problem, 99, 116
Lur'e design, 294
Lyapunov equation, 50
Lyapunov function, 265

## M

Man/nature game, 268
Matrix
spectral radius, 5
Min/max problem, 268
Minimum phase, 46
Mixed control problem, 219, 281
Model matching problem, 186
Modified Newton's method, 365
Multiplicative perturbations, 130, 289

## N

Nash game approach, 229
Nehari theorem, 68
Nominal design, 122, 126
Nominal performance, 132
Nominal system, 266

## O

O - stability, 256
Observability, 325

Optimal filtering, 107
Optimal regulator problem, 104
Optimal state filtering, 159
Optimization
in $R H_{2}, 195$
in $R H_{\infty}, 195$
Orthogonal projection, 189
antistable, 61
stable, 61
Orthogonality
assumption, 104, 156, 165
Outer function, 46
Output estimation, 213
Output estimation problem, 105, 137, 157
Output feedback, 200
stabilizing gains, 256
static, 255
Output feedback design, 232, 240, 248
Output feedback problem, 212, 222
in $\mathrm{RH}_{2}, 210$

## P

Parametric perturbations, 179, 288
Parametrization
$\mathcal{F}_{2 \gamma}, 94,107,116$
$\mathcal{F}_{\infty \gamma}, 149,159,179$
Parrott's theorem, 341
Partial information, 214
Partial information problem, 113, 137, 167, 182
Passivity theorem, 295
PBH test, 325
Performance deterioration, 279
Persidiskii design, 289
Persidiskii theorem, 291
Pole location requirements, 236
Pole placement, 86, 236
regional, 281
Poles, 26
closed-loop, 272
Polynomial, 5, 301
coprime, 302
degree, 6
relative degree, 15
Polynomial matrix, 5
elementary operations, 7
left divisor, 6
rank, 6
right divisor, 6
unimodular, 6
Pseudoinverse, 41

## Q

Quadratic stability, 265

## R

Rational matrix, 5, 15
analytic, 5
greatest common divisors, 24
left divisor, 16
proper, 15
rank, 15
right divisor, 16
strictly proper, 15
unimodular, 15
Reachability, 326
Reconstructability, 325
Region
circular, 236
vertical strip, 239
Riccati equation, 54
stabilizing solution, 316,318
Riccati inequality, 57
Robust design, 125
stability, 130
Robust parametric filtering, 162
Robust performance, 125
complementary sensitivity, 125
control sensitivity, 125
sensitivity, 125,134
Robust stability, $125,132,265$
Robust stabilizability, 152
Root clustering, 283
Root locus, 240

## S

Schur complement, 198, 312
Self-adjoint operator, 45
Sensitivity function, 122
complementary, 122
input, 122
Sensitivity performance, 126
complementary, 123, 126
input, 123, 126
Sensors uncertainty, 124
Separating hyperplane, 350
Separation property, 118, 178, 211
Shaping function, 123
Simultaneous matching, 125

Singular problem, 104
Singular value, 40
decomposition, 40
maximum, 122
Singular values, 303
Small gain theorem, 127
Smith form, 14
Smith-McMillan form, 21
Stability, 5, 50, 195
external, 77
input-output, 5
internal, 5, 71, 73
uncertain matrix, 127
Stabilizability, 69,326
quadratic, 264
Stabilizing controller
quadratic, 268
Standard problem
in $\mathrm{RH}_{2}, 89$
in $R H_{\infty}, 125,137$
State estimation error, 159
State feedback design, 237
State feedback problem, 205, 273
in $\mathrm{RH}_{2}, 202$
in $R H_{\infty}, 217$
State space, 4
Steepest descent direction, 357
Step size
optimal, 367
suboptimal, 367
Strictly positive real, 296
extended, 296
Structural constraints, 281
Structured perturbations, 230
Subgradient, 354
Supporting hyperplane, 351
System zeros, 38

## T

Time domain specifications, 281
Time response decay, 236
Time-domain specifications, 242
Transfer function, 4
proper, 5
realizations, 4
strictly proper, 5
Transmission poles, 27
characterization, 29
Transmission zeros, 27
characterization, 29
rank property, 28
square system, 28
Transpose system, 5
Two-block scheme, 327

## U

Uncertain system, 264
stability, 264
Uncertainty domain
comparison, 271
convex bounded, 265
norm bounded, 266, 288
Unstructured perturbations, 123

## W

Wonham test, 325
Worst case analysis, 343
Worst disturbance, 179

## Y

Youla parametrization, 81

## Z

Zeros, 26


[^0]:    ${ }^{1}$ Single-input single-output
    ${ }^{2}$ Multi-input multi-output
    ${ }^{3}$ Linear quadratic gaussian

