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# The extended symplectic pencil and the finite-horizon LQ problem with two-sided boundary conditions 

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#### Abstract

This note introduces a new approach to the solution of a very general class of finite-horizon optimal control problems for discrete-time systems. This approach provides a parametric expression for the optimal control sequences, as well as the corresponding optimal state trajectories, by exploiting a new decomposition of the so-called extended symplectic pencil. This decomposition provides an original strategy for a more direct solution of the problem with no need of the system-theoretic hypotheses (including regularity of the symplectic pencil) that have always been assumed in the literature so far.


## 1 Introduction

This paper focuses on a very general class of finite-horizon linear-quadratic (LQ) problems with affine constraints at the end-points. These problems are not just important per se. In fairly recent literature it has been shown that LQ problems are becoming increasingly useful as building blocks to solve complex optimisation problems, broken down into two or more LQ subproblems, each one with constraints at the end-points. In particular, finite-horizon LQ problems with constraints at the end-points $[13,14]$ are intermediate steps in the solution of

[^0]$H_{2}$ receding-horizon problems and the minimisation of regulation transients in switching linear plants.

The aim of this paper is to present a method to solve the most general class of finite-horizon LQ optimal control problems in the discrete time with positive semidefinite cost index and affine constraints at the end-points. The proposed solution is based on a procedure for the parameterisation of the set of trajectories generated by the so-called extended symplectic difference equation (ESDE). The idea of solving finite-horizon LQ problems by exploiting expressions of the trajectories generated by the ESDE originated in the papers [3] and [4]. In the past literature, however, the problem solution was always essentially based on two "opposite" solutions of the associated discrete algebraic Riccati equation (with some extra tricks to deal with the case when the closed-loop matrix is singular and hence no pairs of completely opposite solutions exist). This point of view always requires some controllability-type assumption and the extended symplectic pencil [17] to be regular and devoid of generalised eigenvalues on the unit circle. The goal of this paper is to propose a new point of view aimed at a more direct and simple solution to this problem, without requiring system-theoretic assumptions. The technique presented here only requires a solution of the so-called generalised discrete-time algebraic Riccati equation, which may exist even when the symplectic pencil is not regular (in which case the standard discrete algebraic Riccati equation does not admit solutions, let alone pairs of "opposite" solutions). Such solution is used to derive a decomposition of the extended symplectic pencil that yields a natural parameterisation of the solutions of the symplectic difference equation. Thus, while for practical purposes our paper simply provides a generalisation (yet in three different directions) with respect to the existing literature, its different point of view casts a new light on the theoretical comprehension of this problem and on its connections to the classical cornerstones of linear systems theory.

For a better description of the features and the generality of our framework, we illustrate all our results in a running example in which the underlying system is not modulus controllable, and the extended symplectic pencil is not regular (so that the methods in previous literature cannot be used).

## 2 Statement of the problem

Consider the linear time-invariant discrete-time system governed by the difference equation

$$
\begin{equation*}
x(t+1)=A x(t)+B u(t), \tag{1}
\end{equation*}
$$

where, for all $t \in \mathbb{N}, x(t) \in \mathbb{R}^{n}$ is the state, $u(t) \in \mathbb{R}^{m}$ is the control input, $A \in \mathbb{R}^{n \times n}$ and $B \in \mathbb{R}^{n \times m}$. Let $N \in \mathbb{N} \backslash\{0\}$ be the length of the time horizon. Let $V_{0}, V_{N} \in \mathbb{R}^{q \times n}$
and $v \in \mathbb{R}^{q}$; consider

$$
\begin{equation*}
V_{0} x(0)+V_{N} x(N)=v, \tag{2}
\end{equation*}
$$

which represents a two-point boundary-value affine constraint on the states at the end-points. With no loss of generality, we can consider $V \triangleq\left[V_{0} V_{N}\right]$ to be of full row-rank. In the case where $q=0$, the matrices $V_{0}, V_{N}, V$ and the vector $v$ are considered to be void.
Let $\Pi=\left[\begin{array}{cc}Q & S \\ S^{\top} & R\end{array}\right]=\Pi^{\top} \geq 0$ be a square $(n+m)$-dimensional matrix with $Q \in \mathbb{R}^{n \times n}$, $S \in \mathbb{R}^{n \times m}$ and $R \in \mathbb{R}^{m \times m}$ (note that we do not assume the non-singularity of $R$ ). We denote by $\Sigma$ the Popov triple $(A, B, \Pi)$. Finally, let $H=\left[\begin{array}{cc}H_{1} & H_{2} \\ H_{2}^{\top} & H_{3}\end{array}\right]=H^{\top} \geq 0$ with $H_{1}, H_{2}, H_{3} \in \mathbb{R}^{n \times n}$ and $h_{0}, h_{N} \in \mathbb{R}^{n}$.

Problem 1 Find $u(t), t \in\{0, \ldots, N-1\}$ and $x(t), t \in\{0, \ldots, N\}$, minimising

$$
\begin{align*}
J(x, u) & \triangleq \sum_{t=0}^{N-1}\left[\begin{array}{ll}
x^{\top}(t) & u^{\top}(t)
\end{array}\right] \Pi\left[\begin{array}{l}
x(t) \\
u(t)
\end{array}\right] \\
& +\left[\begin{array}{ll}
x^{\top}(0)-h_{0}^{\top} & \left.x^{\top}(N)-h_{N}^{\top}\right] H\left[\begin{array}{c}
x(0)-h_{0} \\
x(N)-h_{N}
\end{array}\right]
\end{array}, .\right. \tag{3}
\end{align*}
$$

under the constraints (1-2).
As discussed in [4], the formulation of Problem 1 is very general, since the cost index in (3) involves the most general type of positive semidefinite quadratic penalisation on the extreme states, and (2) represents the most general affine constraint on these states. As particular cases of Problem 1 we have 1) the standard case where $x(0)$ is assigned and $x(N)$ is weighted in (3); 2) the fixed end-point case, where the states at the end-points are sharply assigned; $\mathbf{3}$ ) the point-to-point case, where the extreme values of an output $y(t)=C x(t)$ are constrained to be equal to two assigned vectors $y_{0}$ and $y_{N}$, respectively. Further non-standard LQ problems that can be useful in practice are particular cases of Problem 1: consider for example an LQ problem in which the states at the end-points $x(0)$ and $x(N)$ are not assigned, but they are constrained to be equal, i.e., $x(0)=x(N)$. This case can be obtained by Problem 1 by setting $V_{0}=I_{n}, V_{N}=-I_{n}$ and $v=0$.

Lemma 1 [4, Lemma 1] If $u(t)$ and $x(t)$ are optimal for Problem 1 , then $\lambda(t) \in \mathbb{R}^{n}$, $t \in\{0, \ldots, N\}$ and $\eta \in \mathbb{R}^{s}$ exist such that $x(t), \lambda(t), u(t)$ and $\eta$ satisfy the set of
equations

$$
\begin{align*}
& x(t+1)=A x(t)+B u(t) \quad t \in\{0, \ldots, N-1\},  \tag{4}\\
& \quad V\left[\begin{array}{c}
x(0) \\
x(N)
\end{array}\right]=v,  \tag{5}\\
& \lambda(t)=Q x(t)+A^{\top} \lambda(t+1)+S u(t) \quad t \in\{0, \ldots, N-1\},  \tag{6}\\
& {\left[\begin{array}{c}
-\lambda(0) \\
\lambda(N)
\end{array}\right]=H\left[\begin{array}{c}
x(0)-h_{0} \\
x(N)-h_{N}
\end{array}\right]+V^{\top} \eta,}  \tag{7}\\
& 0=S^{\top} x(t)+B^{\top} \lambda(t+1)+R u(t) \quad t \in\{0, \ldots, N-1\} . \tag{8}
\end{align*}
$$

Conversely, if equations (4-8) admit solutions $x(t), u(t), \lambda(t), \eta$, then $x(t), u(t)$ minimise $J(x, u)$ subject to the constraints (1-2).

## 3 The generalised Riccati equation and the extended symplectic system

Since in the present setting we are not assuming that $R$ is positive definite, (8) cannot be solved in $u(t)$ to obtain a set of $2 n$ equations in $x(t)$ and $\lambda(t)$. A convenient form in which (4), (6) and (8) can be written, that does not require inversion of $R$, is the descriptor form

$$
\begin{equation*}
F p(t+1)=G p(t) \quad t \in\{0, \ldots, N-1\} \tag{9}
\end{equation*}
$$

where

$$
F \triangleq\left[\begin{array}{ccc}
I_{n} & O & O \\
O & -A^{\top} & O \\
O & -B^{\top} & O
\end{array}\right], \quad G \triangleq\left[\begin{array}{ccc}
A & O & B \\
Q & -I_{n} & S \\
S^{\top} & O & R
\end{array}\right], \quad p(t) \triangleq\left[\begin{array}{c}
x(t) \\
\lambda(t) \\
u(t)
\end{array}\right] .
$$

Notice that there is a small issue in the equivalence between equations (4), (6) and (8) and equation (9). In fact, $u(N)$ does not appear in (4), (6) and (8). Notice, however, that when $u(N)$ appears in (9) it is multiplied by 0 , hence its value is irrelevant. Therefore, we can say that equations (4), (6) and (8) and equation (9) are equivalent, modulo the (arbitrary) value of $u(N)$. The matrix pencil $G-z F$ is known as the extended symplectic pencil, [11, 9], herein denoted concisely by $\operatorname{ESP}(\Sigma)$. In this paper we do not make the assumption of regularity of this pencil.

We will show how to obtain a decomposition of $\operatorname{ESP}(\Sigma)$ that can be used to solve Problem 1 by exploiting the solutions of the following constrained matrix equation

$$
\begin{gather*}
X=A^{\top} X A-\left(A^{\top} X B+S\right)\left(R+B^{\top} X B\right)^{\dagger}\left(B^{\top} X A+S^{\top}\right)+Q  \tag{10}\\
\operatorname{ker}\left(R+B^{\top} X B\right) \subseteq \operatorname{ker}\left(A^{\top} X B+S\right), \tag{11}
\end{gather*}
$$

where (10) has been obtained from the standard discrete algebraic Riccati equation (DARE) by replacing the inverse with the Moore-Penrose pseudo-inverse. Eq. (10) is known in the literature as the generalised discrete-time algebraic Riccati equation $\operatorname{GDARE}(\Sigma),[15,8] \operatorname{GDARE}(\Sigma)$ with the additional constraint given by (11) is sometimes referred to as constrained generalised discrete-time algebraic Riccati equation $\operatorname{CGDARE}(\Sigma)$. Clearly (10) constitutes a generalisation of the classic $\operatorname{DARE}(\Sigma)$, in the sense that any solution of $\operatorname{DARE}(\Sigma)$ is also a solution of $\operatorname{GDARE}(\Sigma)$ - and therefore also of $\operatorname{CGDARE}(\Sigma)$ - but the vice-versa is not true in general. Results on the existence of solutions of $\operatorname{GDARE}(\Sigma)$ in terms of deflating subspaces of the extended symplectic pencil are given in [8] and [9]. We now introduce a standing assumption.

Assumption 3.1 Assume that $\operatorname{CGDARE}(\Sigma)$ has solutions.
Notice that Assumption 3.1 is generically satisfied. The situations in which $\operatorname{CGDARE}(\Sigma)$ does not admit solutions happen to be extremely pathological. Indeed, to the best of the authors' knowledge, no necessary and sufficient existence conditions expressed in terms of the problem data are available for $\operatorname{CGDARE}(\Sigma)$. There are, however, very weak sufficient conditions (see e.g. modulus controllability, [4]) that guarantee existence of solutions of $\operatorname{DARE}(\Sigma)$ - and therefore also of $\operatorname{CGDARE}(\Sigma)$. On the other hand, $\operatorname{CGDARE}(\Sigma)$ generalises $\operatorname{DARE}(\Sigma)$, and may admit solutions even when $\operatorname{DARE}(\Sigma)$ does not. Thus, even in cases in which the aforementioned weak system-theoretic conditions are not satisfied, $\operatorname{CGDARE}(\Sigma)$ may still have solutions. Such solutions can be computed via a reduction to a reduced-order $\operatorname{DARE}(\Sigma)$, see the MATLAB ${ }^{\circledR}$ routine rdare. m in [2], see also [7]. We now introduce some notation that will be used throughout the paper. First, to any matrix $X=X^{\top} \in \mathbb{R}^{n \times n}$ we associate the following matrices:

$$
\begin{align*}
& S_{X} \triangleq A^{\top} X B+S, \quad R_{X} \triangleq R+B^{\top} X B, \quad G_{X} \triangleq I_{m}-R_{X}^{\dagger} R_{X},  \tag{12}\\
& K_{X} \triangleq R_{X}^{\dagger} S_{X}^{\top}, \quad A_{X} \triangleq A-B K_{X} . \tag{13}
\end{align*}
$$

The term $R_{X}^{\dagger} R_{X}$ is the orthogonal projector that projects onto range $R_{X}^{\dagger}=\operatorname{range} R_{X}$ so that $G_{X}$ is the orthogonal projector that projects onto $\operatorname{ker} R_{X}$. Hence, $\operatorname{ker} R_{X}=$ range $G_{X}$.

Example 3.1 The following Popov triple is used as a running example throughout the paper:

$$
A=\left[\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right], \quad B=\left[\begin{array}{ll}
2 & 0 \\
1 & 1
\end{array}\right], \quad Q=\left[\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right], \quad S=R=\left[\begin{array}{ll}
0 & 0 \\
0 & 0
\end{array}\right] .
$$

The extended symplectic pencil in this case is not regular. As such, $\operatorname{DARE}(\Sigma)$ does not admit solutions. On the other hand, in this case $\operatorname{CGDARE}(\Sigma)$ admits
the solution $X=\operatorname{diag}\{0,1\}$, that can be computed by resorting to the algorithm proposed in [2]. In this case, $R_{X}=\left[\begin{array}{ll}1 & 1 \\ 1 & 1\end{array}\right]$, and $A_{X}=\operatorname{diag}\{1,0\}$. Observe that the spectrum of $A_{X}$ is not unmixed, see e.g. [4].

The following result adapts [6, Lemma 2.5] to the case when $G-z F$ may be singular.

Lemma 2 Let $X=X^{\top}$ be a solution of $\operatorname{CGDARE}(\Sigma)$. Then, two invertible matrices $U_{X}, V_{X} \in \mathbb{R}^{2 n+m}$ exist such that

$$
U_{X}(G-z F) V_{X}=\left[\begin{array}{ccc}
A_{X}-z I_{n} & O & B  \tag{14}\\
O & I_{n}-z A_{X}^{\top} & O \\
O & -z B^{\top} & R_{X}
\end{array}\right]
$$

Proof: By direct computation we find

$$
U_{X}(G-z F) V_{X}=\left[\begin{array}{ccc}
A_{X}-z I_{n} & O & B \\
\Xi_{21} & I_{n}-z A_{X}^{\top} & \Xi_{23} \\
\Xi_{23}^{\top} & -z B^{\top} & R_{X}
\end{array}\right]
$$

with

$$
U_{X} \triangleq\left[\begin{array}{ccc}
I_{n} & O & O \\
A_{X}^{\top} X & I_{n} & -K_{X}^{\top} \\
B^{\top} X & O & I_{m}
\end{array}\right] \quad \text { and } \quad V_{X} \triangleq\left[\begin{array}{ccc}
I_{n} & O & O \\
X & -I_{n} & O \\
-K_{X} & O & I_{m}
\end{array}\right] .
$$

The term $\Xi_{21}$ is given by

$$
\begin{aligned}
\Xi_{21}= & A_{X}^{\top} X A-A_{X}^{\top} X B K_{X}+Q-X-S K_{X}-K_{X}^{\top} S^{\top}+K_{X}^{\top} R K_{X} \\
& -z\left(A^{\top} X-A_{X}^{\top} X+K_{X}^{\top} B^{\top} X\right) .
\end{aligned}
$$

The term multiplying $z$ is zero since $A_{X}=A-B K_{X}$. Moreover, since $\operatorname{GDARE}(\Sigma)$ can be written as $X=A^{\top} X A-S_{X} K_{X}+Q$ we find $\Xi_{21}=K_{X}^{\top}\left(R_{X} K_{X}-S_{X}^{\top}\right)=$ $S_{X} R_{X}^{\dagger} R_{X} R_{X}^{\dagger} S_{X}^{\top}-S_{X} R_{X}^{\dagger} S_{X}^{\top}=0$. Finally, $\Xi_{23}=A^{\top} X B-z X B-K_{X}^{\top} B^{\top} X B+S+$ $z X B-K_{X}^{\top} R=S_{X} G_{X}$. In view of (11), we have $S_{X} G_{X}=0$, so that (14) holds.

Remark 1 It is known that the dynamics associated with a matrix pencil is governed by its generalised eigenvalues. ${ }^{1}$ If $X$ is a solution of $\operatorname{CGDARE}(\Sigma)$, from (14) we have

$$
\begin{equation*}
\operatorname{det}(G-z F)=(-1)^{n} \cdot \operatorname{det}\left(A_{X}-z I_{n}\right) \cdot \operatorname{det}\left(I_{n}-z A_{X}^{\top}\right) \cdot \operatorname{det} R_{X} . \tag{15}
\end{equation*}
$$

[^1]When $R_{X}$ is non-singular (i.e. $X$ is also a solution of $\operatorname{DARE}(\Sigma)$ ), the generalised eigenvalues of $G-z F$ are immediately seen to be given by the eigenvalues of $A_{X}$, the reciprocal of the non-zero eigenvalues of $A_{X}$, and a generalised eigenvalue at infinity whose algebraic multiplicity is equal to $m$ plus the algebraic multiplicity of the eigenvalue of $A_{X}$ at the origin. When the matrix $R_{X}$ is singular, the computation of the generalised eigenvalues of $G-z F$ is much more complex. Indeed, in such case (15) still holds but provides no information since $\operatorname{det} R_{X}=0$. We show this fact with a simple example.

Example 3.2 Consider Example 3.1. Matrix $X=\operatorname{diag}\{0,1\}$ is a solution of $\operatorname{CGDARE}(\Sigma)$, and the corresponding closed-loop matrix is $A_{X}=\operatorname{diag}\{1,0\}$. From Lemma 2 we find

$$
U_{X}(G-z F) V_{X}=\left[\begin{array}{cc|cc|cc}
1-z & 0 & 0 & 0 & 2 & 0 \\
0 & -z & 0 & 0 & 1 & 1 \\
\hline 0 & 0 & 1-z & 0 & 0 & 0 \\
0 & 0 & 0 & 1 & 0 & 0 \\
\hline 0 & 0 & -2 z & -z & 1 & 1 \\
0 & 0 & 0 & -z & 1 & 1
\end{array}\right]
$$

whose normal rank (which coincides with that of $G-z F$ ) is easily seen to be equal to 5 . The eigenvalues of $A_{X}$ are 0 and 1 . However, it is not true that $z=1$ is a generalised eigenvalue of $G-z F$. In fact, a direct check shows that the rank of $G-F$ is equal to $5 .{ }^{2}$

From these considerations, it turns out that when $R_{X}$ is singular, the computation of the eigenstructure of the pencil $G-z F$ is more difficult, and requires a different machinery. This machinery hinges on a decomposition of the matrix pencil $G-z F$ for which we need to introduce the following notation. Consider a change of coordinates in the input space $\mathbb{R}^{m}$ induced by the $m \times m$ orthogonal matrix $T=\left[\begin{array}{ll}T_{1} & T_{2}\end{array}\right]$ where range $T_{1}=\operatorname{range} R_{X}$ and range $T_{2}=\operatorname{range} G_{X}=\operatorname{ker} R_{X}$. From [5, Theorems 4.3-4.4], $T$ is independent of the solution $X$ of $\operatorname{CGDARE}(\Sigma)$. Thus $T^{\top} R_{X} T=\operatorname{diag}\left\{R_{X, 0}, O\right\}$, where $R_{X, 0}$ is invertible. Its dimension is denoted by $m_{1}$. Consider the block matrix $\hat{T} \triangleq \operatorname{diag}\left(I_{n}, I_{n}, T\right)$. Defining the matrices $B_{1} \triangleq B T_{1}$ and $B_{2} \triangleq B T_{2}$ we get

$$
\hat{T}^{\top}\left(U_{X}(G-z F) V_{X}\right) \hat{T}=\left[\begin{array}{cccc}
A_{X}-z I_{n} & O & B_{1} & B_{2}  \tag{16}\\
O & I_{n}-z A_{X}^{\top} & O & O \\
O & -z B_{1}^{\top} & R_{X, 0} & O \\
O & -z B_{2}^{\top} & O & O
\end{array}\right] .
$$

[^2]From ker $R_{X}=\operatorname{range} G_{X}$, we obtain range $B_{2}=\operatorname{range}\left(B G_{X}\right)$. Matrix $B_{1}$ has $m_{1}$ columns. Let $m_{2} \triangleq m-m_{1}$ be the number of columns of $B_{2}$. Let us take $U=$ [ $U_{1} U_{2}$ ] such that $U_{1}$ spans the reachable subspace associated with the pair $\left(A_{X}, B_{2}\right)$, denoted by $\mathscr{R}_{X}$, and $U_{2}$ is such that $U$ is invertible. We have

$$
\begin{align*}
U^{-1} A_{X} U & =\left[\begin{array}{cc}
A_{X, 11} & A_{X, 12} \\
O & A_{X, 22}
\end{array}\right], \quad U^{-1} B_{2}=\left[\begin{array}{c}
B_{21} \\
O
\end{array}\right] \\
U^{-1} B_{1} & =\left[\begin{array}{l}
B_{11} \\
B_{12}
\end{array}\right] \tag{17}
\end{align*}
$$

Now, we are ready to state the main result of the paper.
Theorem 1 Let Assumption 3.1 hold. Two invertible matrices $\hat{U}_{X}$ and $\hat{V}_{X}$ exist such that

$$
\begin{align*}
& \hat{U}_{X}(G-z F) \hat{V}_{X}= \\
& {\left[\begin{array}{cc|c|ccc}
A_{X, 11}-z I_{r} B_{21} & O & A_{X, 12} & O & B_{11} \\
\hline O & O & I_{r}-z A_{X, 11}^{\top} & O & O & O \\
O & O & -z B_{21}^{\top} & O & O & O \\
\hline O & O & O & A_{X, 22}-z I_{n-r} & O & B_{12} \\
O & O & -z A_{X, 12}^{\top} & O & I_{n-r}-z A_{X, 22}^{\top} O \\
O & O & -z B_{11}^{\top} & O & -z B_{12}^{\top} & R_{X, 0}
\end{array}\right],} \tag{18}
\end{align*}
$$

where the pair $\left(A_{X, 11}, B_{21}\right)$ is reachable and $R_{X, 0}$ is invertible. Moreover, the matrix pencil $P_{1}(z) \triangleq\left[\begin{array}{ccc}A_{X, 22}-z I_{n-r} & O & B_{12} \\ 0 & I_{n-r}-z A_{X, 22}^{\top} & o \\ 0 & -z B_{12}^{\top} & R_{X, 0}\end{array}\right]$ in (18) is regular, and the generalised eigenvalues of the pencil $G-z F$ are the generalised eigenvalues of $P_{1}(z)$.

The proof of Theorem 1 can be found in the Appendix. The decomposition introduced in Theorem 1 essentially isolates the regular part $P_{1}(z)$ of the pencil $G-z F$. A consequence of this fact is that, unlike the regular case, not all the eigenvalues of $A_{X}$ appear as generalised eigenvalues of $\operatorname{ESP}(\Sigma)$. Indeed, from (18) we have the following

Corollary 1 The finite generalised eigenvalues of $G-z F$ are the uncontrollable eigenvalues of the pair $\left(A_{X}, B_{2}\right)$ plus the reciprocals of those eigenvalues that are not zero.

Example 3.3 Consider Example 3.1. Using the solution $X=\operatorname{diag}\{0,1\}$ of $\operatorname{CGDARE}(\Sigma)$, the null-space and image of $R_{X}$ are respectively spanned by the vectors $\left[\begin{array}{c}-1 \\ 1\end{array}\right]$ and
$\left[\begin{array}{l}1 \\ 1\end{array}\right]$. By taking $T=\left[\begin{array}{cc}1 & -1 \\ 1 & 1\end{array}\right]$ we obtain $T^{\top} R_{X} T=\operatorname{diag}\{4,0\}$. Hence, in this case $m_{1}=m_{2}=1$. We partition $B T$ as $B T=\left[\begin{array}{cc}2 & -2 \\ 2 & 0\end{array}\right]$, so that $B_{1}=\left[\begin{array}{l}2 \\ 2\end{array}\right]$ and $B_{2}=\left[\begin{array}{c}-2 \\ 0\end{array}\right]$. The normal rank of $\operatorname{ESP}(\Sigma)$ is equal to $2 n+m_{1}=5$. The generalised eigenvalues of $G-z F$ are given by the uncontrollable eigenvalues of the pair $\left(A_{X}, B_{2}\right)=\left(\left[\begin{array}{ll}1 & 0 \\ 0 & 0\end{array}\right],\left[\begin{array}{c}-2 \\ 0\end{array}\right]\right)$ plus their reciprocals. Therefore, $\operatorname{ESP}(\Sigma)$ has a generalised eigenvalue at the origin. Since $A_{X, 22}=0$ and $B_{12}=2$, it also has an eigenvalue at infinity with multiplicities equal to the multiplicities of the zero eigenvalue of $\left[\begin{array}{ll}0 & 0 \\ 2 & 0\end{array}\right]$. By writing this pencil in the form given by (18), we get

$$
P(z)=\left[\begin{array}{cc|c|ccc}
1-z & -2 & 0 & 0 & 0 & 2 \\
\hline 0 & 0 & 1-z & 0 & 0 & 0 \\
0 & 0 & 2 z & 0 & 0 & 0 \\
\hline 0 & 0 & 0 & -z & 0 & 2 \\
0 & 0 & 0 & 0 & 1 & 0 \\
0 & 0 & -2 z & 0 & -2 z & 4
\end{array}\right]
$$

from which we see that zero is indeed the only finite generalised eigenvalue of $\operatorname{ESP}(\Sigma)$.

## 4 Solution of the LQ problem

We now consider the problem in the basis constructed in the previous section. Let $\left[\begin{array}{l}x_{1}(t) \\ x_{2}(t)\end{array}\right]=U^{-1} x(t)$ be the coordinates of the state in this basis, partitioned conformably with $U$. Similarly, let $\left[\begin{array}{l}\lambda_{1}(t) \\ \lambda_{2}(t)\end{array}\right]=U^{\top} \lambda(t)$ and $T^{\top} u(t)=\left[\begin{array}{l}u_{1}(t) \\ u_{2}(t)\end{array}\right]$. In this section, we show that in this basis the problem can be easily solved in closed form. More precisely, we first parameterise the solutions of (9) in terms of $x_{1}(0)$, $x_{2}(0), x_{2}(N)$ and $\lambda_{2}(N)$. Then we parameterise the optimal values of $x_{1}(0), x_{2}(0)$, $x_{2}(N)$ and $\lambda_{2}(N)$ by imposing the boundary conditions.

In the new bases, equations (9) can be written for $t \in\{0, \ldots, N-1\}$ as

$$
\begin{align*}
x_{1}(t+1) & =A_{X, 11} x_{1}(t)+B_{21} u_{1}(t)+A_{X, 12} x_{2}(t)+B_{11} u_{2}(t),  \tag{19}\\
\lambda_{1}(t) & =A_{X, 11}^{\top} \lambda_{1}(t+1),  \tag{20}\\
0 & =-B_{21}^{\top} \lambda_{1}(t+1),  \tag{21}\\
x_{2}(t+1) & =A_{X, 22} x_{2}(t)+B_{12} u_{2}(t),  \tag{22}\\
\lambda_{2}(t) & =A_{X, 22}^{\top} \lambda_{2}(t+1)+A_{X, 12}^{\top} \lambda_{1}(t+1),  \tag{23}\\
u_{2}(t) & =R_{X, 0}^{-1} B_{12}^{\top} \lambda_{2}(t+1)+R_{X, 0}^{-1} B_{11}^{\top} \lambda_{1}(t+1) . \tag{24}
\end{align*}
$$

Since by construction the pair $\left(A_{X, 11}, B_{21}\right)$ is reachable, $\operatorname{ker}\left[\begin{array}{c}A_{X, 11}^{\top} \\ B_{21}^{\top}\end{array}\right]=\{0\}$, which means (20-21) yield $\lambda_{1}(t)=0$ for all $t \in\{0, \ldots, N-1\}$. Thus, (23-24) can be simplified as

$$
\begin{equation*}
\lambda_{2}(t)=A_{X, 22}^{\top} \lambda_{2}(t+1), \quad \text { and } \quad u_{2}(t)=R_{X, 0}^{-1} B_{12}^{\top} \lambda_{2}(t+1) . \tag{25}
\end{equation*}
$$

It is clear at this point that we can parameterise all the trajectories generated by the difference equations (22) and (25) in terms of $x_{2}(0)$ and $\lambda_{2}(N)$. Indeed, the first of (25) leads to

$$
\begin{equation*}
\lambda_{2}(t)=\left(A_{X, 22}^{\top}\right)^{N-t} \lambda_{2}(N) \quad \forall t \in\{0, \ldots, N\} . \tag{26}
\end{equation*}
$$

This expression can be plugged into $u_{2}(t)$, and gives

$$
\begin{equation*}
u_{2}(t)=R_{X, 0}^{-1} B_{12}^{\top}\left(A_{X, 22}^{\top}\right)^{N-t-1} \lambda_{2}(N) . \tag{27}
\end{equation*}
$$

Plugging (26) and (27) into (22) gives

$$
\begin{align*}
x_{2}(t)= & A_{X, 22}^{t} x_{2}(0) \\
& +\sum_{j=0}^{t-1} A_{X, 22}^{t-j-1} B_{12} R_{X, 0}^{-1} B_{12}^{\top}\left(A_{X, 22}^{\top}\right)^{N-j-1} \lambda_{2}(N) . \tag{28}
\end{align*}
$$

It is worth observing that

$$
\begin{equation*}
x_{2}(N)=A_{X, 22}^{N} x_{2}(0)+P \lambda_{2}(N), \tag{29}
\end{equation*}
$$

where

$$
P \triangleq \sum_{j=0}^{N-1} A_{X, 22}^{N-j-1} B_{12} R_{X, 0}^{-1} B_{12}^{\top}\left(A_{X, 22}^{\top}\right)^{N-j-1} .
$$

It is easy to see that matrix $P$ can be re-written as $P=\sum_{j=0}^{N-1} A_{X, 22}^{j} B_{12} R_{X, 0}^{-1} B_{12}^{\top}\left(A_{X, 22}^{\top}\right)^{j}$. Therefore, $P$ satisfies the discrete Lyapunov equation

$$
P=A_{X, 22} P A_{X, 22}^{\top}-A_{X, 22}^{N} B_{12} R_{X, 0}^{-1} B_{12}^{\top}\left(A_{X, 22}^{\top}\right)^{N}+B_{12} R_{X, 0}^{-1} B_{12}^{\top} .
$$

If $A_{X, 22}$ has unmixed spectrum, this equation can be used to determine $P$ instead of computing the sum in (29). At this point we can solve (19), which can be written as

$$
\begin{equation*}
x_{1}(t+1)=A_{X, 11} x_{1}(t)+B_{21} u_{1}(t)+\xi(t) \tag{30}
\end{equation*}
$$

where $\xi(t)=A_{X, 12} x_{2}(t)+B_{11} u_{2}(t)$. Using (28) and (27) we find

$$
\begin{aligned}
& \xi(t)=A_{X, 12} A_{X, 22}^{t} x_{2}(0)+\left(B_{11} R_{X, 0}^{-1} B_{12}^{\top}\left(A_{X, 22}^{\top}\right)^{N-t-1}\right. \\
& \left.\quad+A_{X, 12} \sum_{j=0}^{t-1} A_{X, 22}^{t-j-1} B_{12} R_{X, 0}^{-1} B_{12}^{\top}\left(A_{X, 22}^{\top}\right)^{N-j-1}\right) \lambda_{2}(N)
\end{aligned}
$$

Let $R_{1}=\left[B_{21}\left|A_{X, 11} B_{21}\right| A_{X, 11}^{2} B_{21}|\cdots| A_{X, 11}^{N-1} B_{21}\right]$ and $R_{2}=\left[I\left|A_{X, 11}\right| A_{X, 11}^{2} \mid\right.$ $\left.\cdots \mid A_{X, 11}^{N-1}\right]$. Then, we can write $x_{1}(N)=A_{X, 11}^{N} x_{1}(0)+R_{2} \Xi+R_{1} U_{1}$ where $\Xi \triangleq$ $\left[\begin{array}{c}\xi(N-1) \\ \vdots \\ \xi(0)\end{array}\right]$ and $U_{1} \triangleq\left[\begin{array}{c}u_{1}(N-1) \\ \vdots \\ u_{1}(0)\end{array}\right]$. We assume that $N$ is greater than the controllability index of the pair $\left(A_{X, 11}, B_{21}\right)$. All the solutions of this equation are parameterised by

$$
\begin{equation*}
U_{1}=R_{1}^{\dagger}\left(x_{1}(N)-A_{X, 11}^{N} x_{1}(0)-R_{2} \Xi\right)+\left(I-R_{1}^{\dagger} R_{1}\right) v_{1} \tag{31}
\end{equation*}
$$

where $v_{1}$ is arbitrary.

### 4.1 Boundary conditions

In the new basis, the state, co-state and transversality equations can be written again as in (4), (6) and (8), where $A, B, Q, S, V, H, h_{0}$ and $h_{N}$ are replaced by $\tilde{A}=U^{-1} A U, \tilde{B}=U^{-1} B, \tilde{Q}=U^{\top} Q U, \tilde{S}=U^{\top} S, \tilde{V}=V\left[\begin{array}{cc}U & O \\ O & U\end{array}\right], \tilde{H}=$ $\left[\begin{array}{cc}U & O \\ O & U\end{array}\right]^{\top} H\left[\begin{array}{cc}U & O \\ O & U\end{array}\right], \tilde{h}_{0}=U^{-1} h_{0}$ and $\tilde{h}_{N}=U^{-1} h_{N}$, respectively. Now, let us consider the boundary conditions. In this basis, if we partition $\tilde{V}_{0}$ and $\tilde{V}_{N}$ conformably with the state vector, i.e, $\tilde{V}_{0}=\left[\begin{array}{ll}\tilde{V}_{0,1} & \tilde{V}_{0,2}\end{array}\right]$ and $\tilde{V}_{N}=\left[\begin{array}{ll}\tilde{V}_{N, 1} & \tilde{V}_{N, 2}\end{array}\right]$, (5) can be re-written in this basis as

$$
\left[\begin{array}{cccc}
\tilde{V}_{0,1} & \tilde{V}_{N, 1} & \tilde{V}_{0,2}+\tilde{V}_{N, 2} A_{X, 22}^{N} & \tilde{V}_{N, 2} \tag{32}
\end{array}\right] x=v,
$$

where $x=\left[x_{1}^{\top}(0) x_{1}^{\top}(N) x_{2}^{\top}(0) \lambda_{2}^{\top}(N)\right]^{\top}$. Let us now consider (7), and let $\tilde{H}=\left[\begin{array}{cc}\tilde{H}_{1} & \tilde{H}_{2} \\ \tilde{H}_{3} & \tilde{H}_{4}\end{array}\right]$ be partitioned conformably with $\left[\begin{array}{l}x_{1}(t) \\ x_{2}(t)\end{array}\right]$, where $\tilde{H}_{3}=\tilde{H}_{2}^{\top}$. Finally, let $K_{0}$ and $K_{N}$ be basis matrices for $\operatorname{ker} \tilde{V}_{0}$ and $\operatorname{ker} \tilde{V}_{N}$, respectively, to be used to eliminate the multiplier $\eta$. Thus, (7) can be re-written as

$$
\begin{gather*}
-K_{0}^{\top}\left[\begin{array}{l}
\lambda_{1}(0) \\
\lambda_{2}(0)
\end{array}\right]=K_{0}^{\top} \tilde{H}_{1}\left(\left[\begin{array}{l}
x_{1}(0) \\
x_{2}(0)
\end{array}\right]-\tilde{h}_{0}\right)+K_{0}^{\top} \tilde{H}_{2}\left(\left[\begin{array}{l}
x_{1}(N) \\
x_{2}(N)
\end{array}\right]-\tilde{h}_{N}\right),  \tag{33}\\
K_{N}^{\top}\left[\begin{array}{l}
\lambda_{1}(N) \\
\lambda_{2}(N)
\end{array}\right]=K_{N}^{\top} \tilde{H}_{3}\left(\left[\begin{array}{l}
x_{1}(0) \\
x_{2}(0)
\end{array}\right]-\tilde{h}_{0}\right)+K_{N}^{\top} \tilde{H}_{4}\left(\left[\begin{array}{l}
x_{1}(N) \\
x_{2}(N)
\end{array}\right]-\tilde{h}_{N}\right) . \tag{34}
\end{gather*}
$$

Now, defining $\tilde{H}_{i}=\left[\begin{array}{cc}\tilde{H}_{i, 11} & \tilde{H}_{i, 12} \\ \tilde{H}_{i, 21} & \tilde{H}_{i, 22}\end{array}\right]$ and also $\tilde{H}_{i}^{1}=\left[\begin{array}{c}\tilde{H}_{i, 11} \\ \tilde{H}_{i, 21}\end{array}\right]$ and $\tilde{H}_{i}^{2}=\left[\begin{array}{c}\tilde{H}_{i, 12} \\ \tilde{H}_{i, 22}\end{array}\right]$ for $i \in$ $\{1,2,3,4\}$, in this basis (5) and (7) can be expressed as the single linear equation

$$
\begin{equation*}
\left[\frac{F_{1}}{F_{2}}\right] x=g, \tag{35}
\end{equation*}
$$

where

$$
\begin{aligned}
& F_{1}=\left[\begin{array}{l|l|l|l}
\tilde{V}_{0,1} & \tilde{V}_{N, 1} & \tilde{V}_{0,2}+\tilde{V}_{N, 2} A_{X, 22}^{N} & \tilde{V}_{N, 2}
\end{array}\right], \\
& F_{2}=\operatorname{diag}\left\{K_{0}^{\top}, K_{N}^{\top}\right\}\left[\begin{array}{c|c|c|c}
\tilde{H}_{1} & \tilde{H}_{1}^{2}+\tilde{H}_{2}^{2} A_{X, 22}^{N} & \tilde{H}_{2}^{1} & {\left[\begin{array}{c}
-\tilde{H}_{2,12} P \\
\left(A_{X, 22}^{\top}\right)^{N}-\tilde{H}_{2,22} P
\end{array}\right]} \\
\tilde{H}_{3}^{1} & \tilde{H}_{3}^{2}+\tilde{H}_{4}^{2} A_{X, 22}^{N} & \tilde{H}_{4}^{1} & {\left[\begin{array}{c}
\tilde{H}_{4,12} P \\
\tilde{H}_{4,22} P-I
\end{array}\right]}
\end{array}\right], \\
& g=\left[\frac{v}{\operatorname{diag}\left\{K_{0}^{\top}, K_{N}^{\top}\right\} \tilde{H}\left[\begin{array}{l}
h_{0} \\
\tilde{h}_{N}
\end{array}\right]}\right] .
\end{aligned}
$$

We have just proved the following result.
Theorem 2 Under Assumption 3.1, Problem 1 admits solutions if and only if (35) does. For any solution $x=\left[\begin{array}{llll}x_{1}^{\top}(0) & x_{1}^{\top}(N) & x_{2}^{\top}(0) & \lambda_{2}^{\top}(N)\end{array}\right]^{\top}$ we get an optimal initial state $x(0)=\left[\begin{array}{l}x_{1}(0) \\ x_{2}(0)\end{array}\right]$ and a class of optimal controls parameterised by (27) and (31). The solutions obtained in this way are all the solutions of Problem 1.

Example 4.1 Consider a finite-horizon LQ problem in the time interval $\{0, \ldots, N\}$, involving the matrices given in Example 3.1. The initial and final states are constrained to be equal, i.e., $x(0)=x(N)$. Let $H=I_{2 n}, h_{0}=\left[\begin{array}{l}h_{1} \\ h_{2}\end{array}\right]$ and $h_{N}=0$. As aforementioned, $X=\operatorname{diag}\{0,1\}$ is a solution of $\operatorname{CGDARE}(\Sigma)$, leading to $A_{X}=$ $\operatorname{diag}\{1,0\}$. By taking $T=\left[\begin{array}{cc}1 & -1 \\ 1 & 1\end{array}\right]$, we obtained $T^{\top} R_{X} T=\operatorname{diag}\{4,0\}$, so that $R_{0, X}=4, B_{1}=\left[\begin{array}{l}2 \\ 2\end{array}\right]$ and $B_{2}=\left[\begin{array}{c}-2 \\ 0\end{array}\right]$. Therefore, the reachable subspace of the pair $\left(A_{X}, B_{2}\right)$ is spanned by $\left[\begin{array}{l}1 \\ 0\end{array}\right]$, which means this system is already in the desired basis. Thus, $A_{X, 11}=1, A_{X, 12}=A_{X, 22}=0, B_{11}=B_{12}=2$ and $B_{21}=-2$. In this case, (22) and (25) yield $x_{2}(t+1)=B_{12} u_{2}(t), \lambda_{2}(t)=0 \cdot \lambda_{2}(t+1)$, and $u_{2}(t)=R_{X, 0}^{-1} B_{12}^{\top} \lambda_{2}(t+1)$. This implies that $\lambda_{2}(t)=0$ for all $t \in\{0, \ldots, N-1\}$ and is equal to $\lambda_{2}(N)$ for $t=N$, so that $u_{2}(t)=0$ for all $t \in\{0, \ldots, N-2\}$ and $u_{2}(N-1)=R_{X, 0}^{-1} B_{12}^{\top} \lambda_{2}(N)$. Thus $x_{2}(t)$ is equal to $x_{2}(0)$ at $t=0$, is equal to zero for $t \in\{1, \ldots, N-1\}$ and is equal to $B_{12}^{\top} R_{X, 0}^{-1} B_{12}^{\top} \lambda_{2}(N)=\lambda_{2}(N)$ for $t=N$. In this basis, (5) gives rise to $x_{1}(0)=x_{1}(N)$ and $x_{2}(0)=x_{2}(N)=\lambda_{2}(N)$, which are linear in $x_{1}(N)$ and $\lambda_{2}(N)$, while (33-34) can be written as $x_{1}(0)+x_{1}(N)=h_{1}$ and
$x_{2}(0)+x_{2}(N)+\lambda_{2}(0)-\lambda_{2}(N)=h_{2}$. Since $\lambda_{2}(0)=0$ and $x_{2}(N)=\lambda_{2}(N)$, the latter can be written as $x_{2}(0)=h_{2}$. Therefore, the boundary conditions can be written in the form (35). This linear equation admits only the solution $x_{1}(0)=x_{1}(N)=$ $h_{1} / 2$ and $x_{2}(0)=\lambda_{2}(N)=h_{2}$. Now we can compute the optimal control law. First, $u_{2}(t)$ is zero for all $t \in\{0, \ldots, N-2\}$ and $u_{2}(N-1)=R_{X, 0}^{-1} B_{12}^{\top} \lambda_{2}(N)=h_{2} / 2$. To compute $u_{1}$, we write (19) as $x_{1}(t+1)=1 \cdot x_{1}(t)-2 u_{1}(t)+\xi(t)$. The term $\xi(t)$ is zero for all $t \in\{0, \ldots, N-2\}$ and $\xi(N-1)=B_{11} R_{0, X}^{-1} B_{12}^{\top} \lambda_{2}(N)=\lambda_{2}(N)=h_{2}$. We write (31) explicitly as

$$
\begin{aligned}
x_{1}(N)= & x_{1}(0)+\left[\begin{array}{lllll}
I & A_{X, 11} & A_{X, 11}^{2} & \ldots & A_{X, 11}^{N-1}
\end{array}\right]\left[\begin{array}{c}
h_{2} \\
0 \\
\vdots \\
0
\end{array}\right] \\
& +\underbrace{\left[\begin{array}{llll}
-2 & -2 & \ldots & -2
\end{array}\right]}_{N}\left[\begin{array}{c}
u_{1}(N-1) \\
u_{1}(N-2) \\
\vdots \\
u_{1}(0)
\end{array}\right]
\end{aligned}
$$

which gives

$$
\left[\begin{array}{c}
u_{1}(N-1) \\
u_{1}(N-2) \\
\vdots \\
u_{1}(0)
\end{array}\right]=\frac{h_{2}}{2 N}\left[\begin{array}{c}
1 \\
1 \\
\vdots \\
1
\end{array}\right]+\left[\begin{array}{cccc}
1-N & 0 & \ldots & 0 \\
1 & 2-N & \ldots & 0 \\
\vdots & \vdots & \ddots & \vdots \\
1 & 1 & \ldots & -1 \\
1 & 1 & \ldots & 1
\end{array}\right] v,
$$

where $v$ is arbitrary and represents the degree of freedom in the control $u_{1}$.
Remark 2 So far, we have not considered the problem of existence of solutions for Problem 1. In general, the existence of a state trajectory $x(t)$ satisfying the constraints (1-2) for some $u(t)$ is not ensured, since we have not assumed reachability on (1). A necessary and sufficient condition for the existence of optimal solutions is that there exist state and input trajectories satisfying (1-2) (feasible solutions). In fact, since the optimal control problem formulated in Section 2 involves a finite number of variables - precisely, $L=m \cdot N$ for the control plus $n$ for the initial state - Problem 1 can be restated as a quadratic static optimisation problem in these $L+n$ variables with linear constraints. Thus, a solution to Problem 1 exists if and only if a feasible solution - i.e., a state and input functions satisfying both (1) and (2) - exists.

## Proof of Theorem 1

Recall that we have defined the $m \times m$ orthogonal matrix $T=\left[\begin{array}{ll}T_{1} & T_{2}\end{array}\right]$ where range $T_{1}=\operatorname{range} R_{X}$ and range $T_{2}=\operatorname{range} G_{X}=\operatorname{ker} R_{X}$, so that $T^{\top} R_{X} T=\operatorname{diag}\left\{R_{X, 0}, O\right\}$, where $R_{X, 0}$ is invertible. Its dimension is denoted by $m_{1}$. We also defined $B_{1} \triangleq$ $B T_{1}$ and $B_{2} \triangleq B T_{2}$. Moreover, we considered $U=\left[U_{1} U_{2}\right]$ such that $U_{1}$ spans $\mathscr{R}_{X}$, and (17) holds. Let $\hat{U}=\operatorname{diag}\left\{U, U, I_{m_{1}}, I_{m_{2}}\right\}$. Let $r$ denote the size of $\mathscr{R}_{X}$. Using (17) into (16), and taking into consideration and defining the two unimodular matrices
along with $\hat{U}_{X} \triangleq \Omega_{1} \hat{U}^{-1} \hat{T}^{\top} U_{X}$ and $\hat{V}_{X} \triangleq V_{X} \hat{T} \hat{U} \Omega_{2}$, we get (18). Let $P(z)=$ $\hat{U}_{X}(G-z F) \hat{V}_{X}$. Since in (18) the pair $\left(A_{X, 11}, B_{21}\right)$ is reachable by construction, all the $r$ rows of the submatrix $\left[A_{X, 11}-z I_{r} B_{21}\right]$ are linearly independent for every $z \in \mathbb{C} \cup\{\infty\}$. This also means that of the $r+m_{2}$ columns of $\left[A_{X, 11}-z I_{r} B_{21}\right]$, only $r$ are linearly independent, and this gives rise to the presence of a null-space of $P(z)$ whose dimension $m_{2}$ is independent of $z \in \mathbb{C} \cup\{\infty\}$. We obtain ${ }^{3}$

$$
\begin{aligned}
& \operatorname{rank} P(z)=r+ \\
& \quad \operatorname{rank}\left[\begin{array}{c|ccc}
I_{r}-z A_{X, 11}^{\top} & O & O & O \\
-z B_{21}^{\top} & O & O & O \\
\hline O & A_{X, 22}-z I_{n-r} & O & B_{12} \\
-z A_{X, 12}^{\top} & O & I_{n-r}-z A_{X, 22}^{\top} & O \\
-z B_{11}^{\top} & O & -z B_{12}^{\top} & R_{X, 0}
\end{array}\right] .
\end{aligned}
$$

Now, consider the rank of $\left[\begin{array}{c}I_{r}-z A_{X, 11}^{\top} \\ -z B_{21}^{\top}\end{array}\right]$. Again, since the pair $\left(A_{X, 11}, B_{21}\right)$ is reachable, this rank is constant and equal to $r$ for every $z \in \mathbb{C} \cup\{\infty\}$. Thus, $\operatorname{rank} P(z)=$ $2 r+\operatorname{rank} P_{1}(z)$. Since $\operatorname{det} P_{1}(z)=\operatorname{det}\left(A_{X, 22}-z I_{n-r}\right) \cdot \operatorname{det}\left(I_{n-r}-z A_{X, 22}^{\top}\right) \cdot \operatorname{det} R_{X, 0}$, a value $z \in \mathbb{C}$ can be found for which $\operatorname{det} P_{1}(z) \neq 0$. Hence, the normal rank of $P_{1}(z)$ is equal to $2(n-r)+m_{1}$, and therefore the normal rank of $P(z)$ is $2 r+2(n-r)+m_{1}=2 n+m_{1}$. The generalised eigenvalues of the pencil $P(z)$ are the values $z \in \mathbb{C} \cup\{\infty\}$ for which the rank of $P_{1}(z)$ is smaller than its normal rank $2(n-r)+m_{1}$.

[^3]
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[^1]:    ${ }^{1}$ Recall that a generalised eigenvalue of a matrix pencil $G-z F$ is a value of $z \in \mathbb{C}$ for which the rank of the matrix pencil $G-z F$ is lower than its normal rank.

[^2]:    ${ }^{2}$ We warn that the routine eig.m of the software MATLAB ${ }^{\circledR}$ (version 7.11.0.584(R2010b)) in this case fails to provide the right answer. It indeed returns 1 as a generalised eigenvalue of the pencil $G-z F$.

[^3]:    ${ }^{3}$ Let $\Xi=\left[\begin{array}{cc}\Xi_{11} & \Xi_{12} \\ 0 & \Xi_{22}\end{array}\right]$. Observe that if either $\Xi_{11}$ is full row-rank or $\Xi_{22}$ is full column-rank, then $\operatorname{rank} \Xi=\operatorname{rank} \Xi_{11}+\operatorname{rank} \Xi_{22}$.

