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Simultaneous State and Input Reachability for Linear Time Invariant Systems

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Abstract

In this paper, we give an explicit solution to the behavioral reachability problem for linear time invariant systems, which amounts to finding an explicit control law that reaches a given final input-state pair (u_1, x_1) in a given finite time t_1 . We first tackle the case of state space realizations, and we then extend the obtained results to the case of implicit realizations. For this, we use the geometric approach and some results of the viability theory. Some complements are given about the existing relationships between reachability and pole placement, as well as some notions of unicity and existence of solution.

Keywords:

Linear systems, implicit systems, reachability, geometric approach.

Notation. Script capitals \mathscr{V} , \mathscr{W} , ..., denote finite dimensional linear spaces with elements v, w, \ldots ; the dimension of a space \mathscr{V} is denoted $\dim(\mathscr{V})$; $\mathscr{V} \approx \mathscr{W}$ stands for $\dim(\mathscr{V}) = \dim(\mathscr{W})$; when $\mathscr{V} \subset \mathscr{W}$, $\frac{\mathscr{W}}{\mathscr{V}}$ or \mathscr{W}/\mathscr{V} stand for the quotient

space \mathscr{W} modulo \mathscr{V} ; the direct sum of independent spaces is written as \oplus . $X^{-1}\mathscr{V}$, stands for the inverse image of the subspace \mathscr{V} by the linear transformation X. Given a linear transformation $X: \mathscr{V} \to \mathscr{W}$, $\operatorname{Im} X = X\mathscr{V}$ denotes its image, and $\operatorname{Ker} X$ denotes its kernel; when $\mathscr{V} \approx \mathscr{W}$, we write $X: \mathscr{V} \leftrightarrow \mathscr{W}$; when $\mathscr{U} \subset \mathscr{V}$, $X|_{\mathscr{U}}$ denotes the restriction of X to \mathscr{U} . Given the space $\mathscr{X} = \mathscr{S} \oplus \mathscr{T}$, the natural projection, $P: \mathscr{X} \to \mathscr{S}$, on \mathscr{S} along \mathscr{T} , is also written as $P: \mathscr{X} \to \mathscr{S}//\mathscr{T}$. The special subspaces $\operatorname{Im} B$, $\operatorname{Ker} E$ and $\operatorname{Ker} C$, are denoted by \mathscr{B} , \mathscr{K}_E and \mathscr{K}_C , respectively. The zero dimension subspace is denoted $\{0\}$, and the identity operator is denoted I, namely Ix = x. Given the linear transformations $X: \mathscr{V} \to \mathscr{V}$ and $Y: \mathscr{W} \to \mathscr{V}$, $\langle X \mid \operatorname{Im} Y \rangle$ stands for the subspace of \mathscr{V} : $\operatorname{Im} Y + X \operatorname{Im} Y + \cdots + X^{\dim(\mathscr{V})-1} \operatorname{Im} Y$. The notations A_{F_p} and E_{F_d} stand for $(A + BF_p)$ and $(E - BF_d)$, respectively.

BDM $\{X_1, ..., X_k\}$ denotes a block diagonal matrix whose diagonal blocks are the matrices $X_1, ..., X_k$, and DM $\{x_1, ..., x_k\}$ denotes a diagonal matrix whose diagonal elements are $x_1, ..., x_k$. The notation \mathbb{R}^k stands for the Euclidean space of dimension k. $\underline{e}_k^i \in \mathbb{R}^k$ stands for the vector whose i-th entry is equal to 1 and the other ones are equal to 0. $T_u\{v^T\}$ stands for the upper triangular Toeplitz matrix, whose first row is v^T . * stands for some matrix which exact value has no importance.

 \mathbb{R}^+ , \mathbb{R}^{+*} and \mathbb{Z}^+ , stand for the sets of non negative real numbers, positive real numbers and non negative integers, respectively. $\mathcal{C}^{\infty}(\mathbb{R}^+, \mathcal{V})$ and $\mathcal{L}^{\infty}(\mathbb{R}^+, \mathcal{V})$ are the space of infinitely differentiable functions and the space of bounded functions, $v: \mathbb{R}^+ \to \mathcal{V}$, respectively. $\mathcal{L}_{1}^{loc}(\mathbb{R}^+, \mathcal{V})$ stands for the locally integrable functions.

Geometric Algorithms. Given the linear transformations $X : \mathcal{V} \to \mathcal{W}$, $Y : \mathcal{T} \to \mathcal{W}$, and $Z : \mathcal{V} \to \mathcal{W}$, and the subspace $\mathcal{K} \subset \mathcal{V}$, we have the two following popular geometric algorithms (see mainly Verghese, 1981, Özçaldiran, 1986, Malabre, 1987, 1989, Lewis, 1992):

$$\mathscr{V}^{0}_{[\mathscr{K}:X,Z,Y]} = \mathscr{V} , \quad \mathscr{V}^{\mu+1}_{[\mathscr{K}:X,Z,Y]} = \mathscr{K} \cap X^{-1} \left(Z \mathscr{V}^{\mu}_{[\mathscr{K}:X,Z,Y]} + \operatorname{Im} Y \right) \quad (\text{ALG-V})$$

$$\mathscr{S}^{0}_{[Z,X,Y]} = \{0\}, \quad \mathscr{S}^{\mu+1}_{[Z,X,Y]} = Z^{-1} \left(X \mathscr{S}^{\mu}_{[Z,X,Y]} + \operatorname{Im} Y \right)$$
 (ALG-S)

where $\mu \in \mathbb{Z}^+$. The limit of (ALG–V) is the supremal (X, Z, Y) invariant subspace contained in \mathscr{K} , $\mathscr{V}^*_{[\mathscr{K}:X,Z,Y]} := \sup\{\mathscr{S} \subset \mathscr{K} | X\mathscr{S} \subset Z\mathscr{S} + \operatorname{Im} Y\}$, and the limit of (ALG–S) is the infimal (Z,X,Y) invariant subspace related to $\operatorname{Im} Y$, $\mathscr{S}^*_{[Z,X,Y]} := \inf\{\mathscr{S} \subset \mathscr{V} \mid \mathscr{S} = Z^{-1}(X\mathscr{S} + \operatorname{Im} Y)\}$.

We distinguish two cases.

- For the square brackets $[\mathscr{V}:X,Z,0]$ and [Z,X,0], we write: $\mathscr{V}^*_{[X,Z]}, \mathscr{V}^{\mu}_{[X,Z]}, \mathscr{S}^*_{[Z,X]}$ and $\mathscr{S}^{\mu}_{[Z,X]}$, instead of: $\mathscr{V}^*_{[\mathscr{V}:X,Z,0]}, \mathscr{V}^{\mu}_{[\mathscr{V}:X,Z,0]}, \mathscr{S}^*_{[Z,X,0]}$ and $\mathscr{S}^{\mu}_{[Z,X,0]}, \mathscr{S}^{\mu}_{[Z,X,0]}$ respectively, where $\mu \in \mathbb{Z}^+$.
- For the square bracket $\left[\mathcal{K}_{\overline{C}} : \overline{A}, I, Y \right]$, we write: $\overline{\mathcal{V}}_{Y}^{*}$ and $\overline{\mathcal{V}}_{Y}^{\mu}$, instead of: $\mathcal{V}_{\left[\mathcal{K}_{\overline{C}} : \overline{A}, I, Y\right]}^{*}$ and $\mathcal{V}_{\left[\mathcal{K}_{\overline{C}} : \overline{A}, I, Y\right]}^{\mu}$, where $\mu \in \mathbb{Z}^{+}$.

Subspaces. Note that in the particular case $X = A : \mathcal{X} \to \mathcal{X}, Y = B : \mathcal{U} \to \mathcal{X}$, and Z = I, the equalities $\mathscr{V}^*_{[\mathcal{X}:A,I,B]} = \mathcal{X}$ and $\mathscr{S}^*_{[I,A,B]} = \langle A \mid \mathscr{B} \rangle$ hold true.

Given the linear transformations $X = A : \mathscr{X}_d \to \mathscr{X}_{eq}, Y = B : \mathscr{U} \to \mathscr{X}_{eq}$, and $Z = E : \mathscr{X}_d \to \mathscr{X}_{eq}$, it is observed that

• the supremal (A, E, B) invariant subspace contained in \mathscr{X}_d and the infimal (E, A, B) invariant subspace related to \mathscr{B} , $\mathscr{V}^*_{[\mathscr{X}_d: A, E, B]}$ and $\mathscr{S}^*_{[E, A, B]}$, are identified by $\mathscr{V}_{\mathscr{X}_d}^*$ and $\mathscr{S}_{\mathscr{X}_d}^*$, respectively, and the respective subspaces of their algorithms (ALG–V) and (ALG–S) are identified by $\mathscr{V}_{\mathscr{X}_d}^{\mu}$ and $\mathscr{S}_{\mathscr{X}_d}^{\mu}$ ($\mu \in \mathbb{Z}^+$), respectively;

- the supremal (A, E, B) invariant subspace contained in ℋ_C, Ψ^{*}_[ℋ_C: A, E, B], is identified by Ψ^{*}, and the respective subspaces of its algorithm (ALG–V) are identified by Ψ^μ (μ ∈ Z⁺);
- the unobservable space $\mathscr{V}^*_{[\mathscr{K}_C:A,E,0]}$ is identified by \mathscr{N} ; and the closed loop unobservable space $\mathscr{V}^*_{[\mathscr{K}_C:A_{F_p},E_{F_d},0]}$ is identified by $\mathscr{N}_{(F_p,F_d)}$.

Let us note that:

$$(i) \mathcal{V}_{[\mathcal{K}:A,E,B]}^* = \mathcal{V}_{[\mathcal{K}:A_{F_n},E_{F_d},B]}^*,$$

(ii)
$$\mathscr{S}^*_{[\mathscr{K}:E,A,B]} = \mathscr{S}^*_{[\mathscr{K}:E_{F_d},A_{F_p},B]}$$
, and

(iii) for any F_d , there exists F_p such that: $A_{F_p} \mathcal{V}_{[\mathcal{K}:A_{F_p},E_{F_d},B]}^* \subset E_{F_d} \mathcal{V}_{[\mathcal{K}:A_{F_p},E_{F_d},B]}^*$. The set of such pairs (F_p,F_d) is identified by $\mathbf{F}(\mathcal{V}_{[\mathcal{K}:A,E,B]}^*)$.

1. INTRODUCTION

- One of the most studied concepts in System Theory is the one of reacha-
- bility. This concept is normally associated with the set of vectors that can be
- 4 reached from the origin in a finite time, following trajectories solutions of the
- 5 system, generated by the input system. Here, the term input system refers
- 6 to an exogenous signal which is available for controlling the output system.
- 7 1.1. State Space Representations
- For the case of state space representations $\Re^{ss}(A, B)$,

$$dx/dt = Ax + Bu, (1.1)$$

where $u \in \mathcal{U} \approx \mathbb{R}^m$ is the input variable, $x \in \mathcal{X} \approx \mathbb{R}^n$ is the state variable, and with the usual assumption $\operatorname{Ker} B = \{0\}$, Kalman (1960, 1963) introduced his famous reachability matrix: $\mathcal{R}_{[A,B]} = \begin{bmatrix} B & AB & \cdots & A^{n-1}B \end{bmatrix}$. He showed that given any $x_0, x_1 \in \mathcal{X}$, there exists a control $\operatorname{law}^1 \ u(\cdot) \in \mathcal{C}^{\infty}(\mathbb{R}^+, \mathcal{U})$, generating a trajectory $x(\cdot) \in \mathcal{C}^{\infty}(\mathbb{R}^+, \mathcal{X})$ solution of (1.1), starting from the given initial condition $x(0) = x_0 \in \mathcal{X}$, and reaching the desired final state $x(t_1) = x_1 \in \mathcal{X}$, in a finite time $t_1 \in \mathbb{R}^+$, iff, $\operatorname{rank}(\mathcal{R}_{[A,B]}) = n$; in this case, the representation (1.1) is called reachable. This concept is known as state reachability², and when the pair (A, B) satisfies such a rank condition, we identify it as a state reachable pair.

A. State reachability. Brunovsky (1970) showed that for a given reachable state space representation (1.1), there exist a linear map $F_{\mathcal{B}}: \mathbb{R}^n \to \mathbb{R}^m$ and isomorphisms $T_{\mathcal{B}}: \mathbb{R}^n \to \mathbb{R}^n$ and $G_{\mathcal{B}}: \mathbb{R}^m \to \mathbb{R}^m$, such that the pair $(A_{\mathcal{B}}, B_{\mathcal{B}})$, where $A_{\mathcal{B}} = T_{\mathcal{B}}^{-1}(A + BF_{\mathcal{B}})T_{\mathcal{B}}$ and $B_{\mathcal{B}} = T_{\mathcal{B}}^{-1}BG_{\mathcal{B}}$, is expressed in the Brunovsky

¹ In this paper, we restrain our discussion to infinitely differentiable functions. This is not restrictive since $\mathcal{C}^{\infty}(\mathbb{R}^+, \mathscr{V})$ is dense in $\mathcal{L}_1^{\mathrm{loc}}(\mathbb{R}^+, \mathscr{V})$ (see Polderman & Willems, 1998, Corollary 2.4.12).

² In many text books, this property is called *state controllability*, or simply *controllability*. Let us note that *controllability* only characterizes the system's property of reaching the origin $x_1 = 0$, from any state $x_0 \neq 0$, in a finite time t_1 . Since in the continuous time-invariant linear systems case both properties, *reachability* and *controllability*, are mutually implied, they are often treated indistinguishably, but in the general case of the implicit representations, this is no longer the case; for example the implicit representation, $\begin{bmatrix} 1 & 0 \\ 0 & 0 \end{bmatrix} dx/dt = \begin{bmatrix} 0 & 0 \\ 0 & 1 \end{bmatrix} x + \begin{bmatrix} 1 \\ 0 \end{bmatrix} u$, is trivially controllable but not reachable.

23 canonical form, namely:

$$A_{\mathcal{B}} = \operatorname{BDM} \left\{ A_{\mathcal{B}_1}, \dots, A_{\mathcal{B}_m} \right\}, \ B_{\mathcal{B}} = \operatorname{BDM} \left\{ b_{\mathcal{B}_1}, \dots, b_{\mathcal{B}_m} \right\},$$

$$\left[A_{\mathcal{B}_i} \mid b_{\mathcal{B}_i} \right] = \left[\operatorname{T}_u \left\{ (\underline{e}_{\kappa_i}^2)^T \right\} \mid \underline{e}_{\kappa_i}^{\kappa_i} \right], \ i \in \{1, \dots, m\},$$

$$(1.2)$$

where the set

$$S_{\kappa} = \left\{ \{\kappa_1, \kappa_2, \dots, \kappa_m\} \subset \mathbb{Z} \mid \kappa_1 \ge \kappa_2 \ge \dots \ge \kappa_m \ge 1 \& \sum_{i=1}^m \kappa_i = n \right\} \quad (1.3)$$

25 is known as the set of *reachability indices*. They are also geometrically char-

acterized as follows

$$\operatorname{card}\left\{\kappa_{i} \geq 1\right\} = \dim\left(\mathscr{B}\right) \text{ and } \operatorname{card}\left\{\kappa_{i} \geq \mu\right\} = \dim\left(\frac{\sum_{i=0}^{\mu-1} A^{i}\mathscr{B}}{\sum_{i=0}^{\mu-2} A^{i}\mathscr{B}}\right), \quad \forall \, \mu \geq 2.$$

Another important success was the introduction of the reachable space $\langle A | \mathcal{B} \rangle$. In (Wonham, 1985) is showed that a pair (A, B) is reachable iff:

$$\langle A \mid \mathscr{B} \rangle = \mathscr{X} . \tag{1.4}$$

Note that $\langle A | \mathscr{B} \rangle = \mathscr{X}$ iff rank $(\mathcal{R}_{[A,B]}) = n$. Wonham (1985) showed that the

reachability Gramian $W_{t_1} = \int_0^{t_1} \exp(\tau A) BB^T \exp(\tau A^T) d\tau$, with $t \in [0, t_1]$, is non-

singular iff (1.4) is satisfied. Thus, with the control law

$$u(t) = B^{T} \exp((t_{1} - t)A^{T}) W_{t_{1}}^{-1}(x_{1} - \exp(t_{1}A) x_{0}), \qquad (1.5)$$

we get a trajectory $x(\cdot) \in \mathcal{C}^{\infty}(\mathbb{R}^+, \mathscr{X})$ solution of (1.1), such that $x(0) = x_0 \in \mathscr{X}$ and $x(t_1) = x_1 \in \mathscr{X}$.

Another well known result concerning state reachability is the one related with pole assignment. Indeed, the pair (A, B) is reachable iff for every symmetric (with respect to the real line) set of complex numbers Λ , of cardinality

37 n, there exists a proportional state feedback u = Fx such that the spectrum of

 $_{38}$ $(\lambda I - A_F)$ is Λ (see for example Theorems 2.1 and 9.3.1 of Wonham (1985)

³⁹ and Polderman & Willems (1998), respectively).

B. Behavioral reachability. Willems (1983, 1991) defined an input/state system as the triple $\Sigma_{i/s} = (\mathbb{R}^+, \mathscr{U} \times \mathscr{X}, \mathfrak{B}_{[A,B]})$, with behavior³

$$\mathfrak{B}_{[A,B]} = \left\{ (u,x) \in \mathcal{C}^{\infty}(\mathbb{R}^+, \mathscr{U} \times \mathscr{X}) \middle| \quad \left[\begin{array}{cc} (I \frac{\mathrm{d}}{\mathrm{d}t} - A) & -B \end{array} \right] \begin{bmatrix} x \\ u \end{bmatrix} = 0 \right\}. \quad (1.6)$$

In the behavioral framework of Willems (1983, 1991), the system $\Sigma_{i/s} = (\mathbb{R}^+, \mathscr{U} \times \mathscr{X}, \, \mathfrak{B}_{[A,B]})$ is called **reachable* if for any given (u_0, x_0) , $(u_1, x_1) \in \mathscr{U} \times \mathscr{X}$ and $t_1 > 0$, it is possible to find a trajectory $(u, x) \in \mathfrak{B}_{[A,B]}$, such that $(u(0), x(0)) = (u_0, x_0)$ and $(u(t_1), x(t_1)) = (u_1, x_1)$ (c.f. Polderman & Willems, 1998, Definition 5.2.2). In the following, this reachability concept is called **behavioral reachability*.

In (Polderman & Willems, 1998, Theorem 5.2.27) is proved that for the

In (Polderman & Willems, 1998, Theorem 5.2.27) is proved that for the case of state space representations $\mathfrak{R}^{ss}(A, B)$, state rechability is equivalent to behavioral reachability. Although the behavioral reachability is well characterized, it could be interesting to find an explicit control law $u(\cdot) \in \mathcal{C}^{\infty}(\mathbb{R}^+, \mathscr{U})$, similar to (1.5), which ensures $x(t_1) = x_1$ and $u(t_1) = u_1$. This will be done in Section 2.

54 1.2. Implicit Representations

As a generalization of proper linear systems, Rosenbrock (1970) introduced the *implicit representations* $\mathfrak{R}^{imp}(E, A, B)$, which are a set of differential

The original definition given by Willems (1983, 1991) is $\mathfrak{B}_{[A,B]} = \{(u,x) \in \mathcal{L}_1^{\mathrm{loc}}(\mathbb{R}^+, \mathcal{U} \times \mathcal{X}) \mid \exists x_0 \in \mathcal{X} \text{ s.t. } x(t) = \exp(At) x_0 + \int_0^t \exp(A(t-\tau)) Bu(\tau) d\tau \}$. But since we restrict our attention to infinitely differentiable functions (see footnote ¹), weak and strong solutions coincide (see Polderman & Willems, 1998, Theorem 2.3.11).

⁴ For consistency of this paper, we say reachable instead of controllable, as is stated in (Polderman & Willems, 1998) (see also footnote 2).

and algebraic equations (Brenan *et al*, 1996) of the following form (see also Lewis, 1992)

$$E dx/dt = Ax + Bu, (1.7)$$

where $E: \mathscr{X}_d \to \mathscr{X}_{eq}$, $A: \mathscr{X}_d \to \mathscr{X}_{eq}$ and $B: \mathscr{U} \to \mathscr{X}_{eq}$ are linear maps. The linear spaces $\mathscr{X}_d \approx \mathbb{R}^{n_d}$, $\mathscr{X}_{eq} \approx \mathbb{R}^{n_{eq}}$, and $\mathscr{U} \approx \mathbb{R}^m$ are called the descriptor, the equation, and the input spaces, respectively. In order to avoid redundant components in the input variable u, and linear dependence on the descriptor equations (1.7), as usually, we assume troughout the paper that the following hypotheses are verified:

65 **[H1]** Ker B = 0, and

66
$$[\mathbf{H2}] \operatorname{Im} E + \operatorname{Im} A + \mathscr{B} = \mathscr{X}_{eq}$$
.

For the case of regular implicit representations, *i.e.* representations where the linear transformations E and A are square and the pencil $[\lambda E - A]$ is invertible (Gantmacher, 1977), the *reachability* was studied by Verghese, Lévy and Kailath (1981) from a transfer function point of view, Yip and Sincovec (1981) in the time domain, Cobb (1984) from a distributional point of view, and by Özçaldiran (1985) from a geometric point of view.

In the case of implicit representations, where the linear transformations E and E are square and the pencil $[\lambda E - A]$ is not necessarily invertible, Özçaldiran (1986) extended his F reachability geometric characterization for the case of regular implicit representations (Özçaldiran, 1985), by means of the supremal (A, E, B) reachability subspace contained in \mathcal{X}_d , defined as

$$\mathscr{R}_{\mathscr{X}_d}^* = \mathscr{V}_{\mathscr{X}_d}^* \cap \mathscr{S}_{\mathscr{X}_d}^* . \tag{1.8}$$

This is a nice generalization of the classical case, $\mathfrak{R}^{ss}(A, B) = \mathfrak{R}^{imp}(I, A, B)$,

where the reachable space $\mathscr{R}^*_{\mathscr{X}_d}$ is equal to $\langle A \mid \mathscr{B} \rangle$, namely equal to $\mathscr{V}^*_{[\mathscr{X}:A,I,B]} \cap$

 $\mathscr{S}^*_{[I,A,B]}$. Thus, for representations $\mathfrak{R}^{imp}(E,A,B)$, with E and A not necessarily square, it was natural to associate its reachability with $\mathscr{R}^*_{\mathscr{X}_d}$.

Frankowska (1990) firmly established the pertinence of this reachability concept, using differential inclusions to relate it with behavioral properties.

One major difficulty when studying reachability for implicit systems (1.7) is that their solution set does not only depend on the initial conditions x(0) and on the external control input u, but also depends on a possible internal free variable (degree of freedom), which is completely unknown.

88 1.3. Outline

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In this paper, we study the reachability notion in the sense of Frankowska (1990), showing some connections with the important works of Willems (1991) and Geerts (1993), and we consider the relationships between the reachability property and the complete pole assignment ability.

The paper is organized as follows: In Section 2, we consider the behavioral reachability problem for state space representations, namely the ability of reaching the input-state pair $(u(\cdot), x(\cdot))$. In Section 3, we formalize the notion of *implicit systems*, following the behavioral point of view, and we also study the equivalences between the notions of existence of solution and impulse controllability. In Section 4, we study the reachability notion of Frankowska (1990) for implicit systems. In Section 5, we consider the existent relationships between the reachability property and the complete pole assignment ability, and in Section 6, we conclude the paper.

2. BEHAVIORAL REACHABILITY PROBLEM

We consider the following problem.

Problem 1. Let us consider an input/state system $\Sigma_{i/s} = (\mathbb{R}^+, \mathscr{U} \times \mathscr{X}, \mathfrak{B}_{[A,B]})$ represented by (1.1), and with the behavior (1.6). Given (u_0, x_0) , $(u_1, x_1) \in$ $B^{-1}\langle A | \mathscr{B} \rangle \times \langle A | \mathscr{B} \rangle$ and $t_1 > 0$, find a trajectory $(u, x) \in \mathfrak{B}_{[A,B]}$, such that $(u(0), x(0)) = (u_0, x_0)$ and $(u(t_1), x(t_1)) = (u_1, x_1)$.

This is the *behavioral reachability* problem, and in (Polderman & Willems, 1998, Theorem 5.2.27) is proved that for the case of *state space representations*, state reachability is equivalent to behavioral reachability. So, condition (1.4) guarantees the existence of a solution for Problem 1.

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One could think that the control law (1.5), proposed by Wonham (1985), solves Problem 1, but this proposition only guarantees the *reachability* of the state variable, $x(0) = x_0$ and $x(t_1) = x_1$, and nothing about the input variable u, which is let completely free at the end points u(0) and $u(t_1)$. An intermediary step towards the solution of Problem 1 is given by the next result proved in Appendix A.

Lemma 1. Let the state space representation (1.1) be reachable, with the reachability indices set (1.3). Let the linear map $F_{\mathcal{B}}: \mathbb{R}^n \to \mathbb{R}^m$ and the isomorphisms $T_{\mathcal{B}}: \mathbb{R}^n \to \mathbb{R}^n$ and $G_{\mathcal{B}}: \mathbb{R}^m \to \mathbb{R}^m$ be such that the pair $(A_{\mathcal{B}}, B_{\mathcal{B}})$, where $A_{\mathcal{B}} = T_{\mathcal{B}}^{-1}(A + BF_{\mathcal{B}})T_{\mathcal{B}}$ and $B_{\mathcal{B}} = T_{\mathcal{B}}^{-1}BG_{\mathcal{B}}$, is expressed in the Brunovsky canonical form (1.2). Let the reachability matrices, $\mathcal{R}_{[A_{\mathcal{B}}, B_{\mathcal{B}}]}$ and $\mathcal{R}_{[A_{\mathcal{B}_i}, b_{\mathcal{B}_i}]}$, $i \in \{1, \ldots, m\}$, of the pair $(A_{\mathcal{B}}, B_{\mathcal{B}})$ and the pairs $(A_{\mathcal{B}_i}, b_{\mathcal{B}_i})$, $i \in \{1, \ldots, m\}$, respectively, be defined as follows:

$$\mathcal{R}_{[A_{\mathcal{B}}, B_{\mathcal{B}}]} = \text{BDM} \left\{ \mathcal{R}_{[A_{\mathcal{B}_1}, b_{\mathcal{B}_1}]}, \dots, \mathcal{R}_{[A_{\mathcal{B}_m}, b_{\mathcal{B}_m}]} \right\},
\mathcal{R}_{[A_{\mathcal{B}_i}, b_{\mathcal{B}_i}]} = \begin{bmatrix} b_{\mathcal{B}_i} & A_{\mathcal{B}_i} b_{\mathcal{B}_i} & \cdots & A_{\mathcal{B}_i}^{\kappa_i - 1} b_{\mathcal{B}_i} \end{bmatrix}.$$
(2.1)

Let us assume that we have found trajectories $f_i \in \mathcal{C}^{\infty}(\mathbb{R}^+, \mathbb{R}^1)$, satisfying:

126 (i) for
$$j = 0, 1$$

$$\mathfrak{D}(\frac{\mathrm{d}}{\mathrm{d}t})f(t_j) = G_{\mathcal{B}}^{-1}(u_j - F_{\mathcal{B}}x_j),\tag{2.2}$$

where
$$\mathfrak{D}(d/dt) = DM \{d^{\kappa_1}/dt^{\kappa_1}, \ldots, d^{\kappa_m}/dt^{\kappa_m}\}, f(t) = \begin{bmatrix} f_1(t) & \cdots & f_m(t) \end{bmatrix}^T$$
and $t_0 = 0$.

129 (ii) If
$$\overline{w}_i(t) = \begin{bmatrix} \frac{\mathrm{d}^{\kappa_i - 1} f_i(t)}{\mathrm{d} t^{\kappa_i - 1}} & \dots & \frac{\mathrm{d} f_i(t)}{\mathrm{d} t} & f_i(t) \end{bmatrix}^T$$
, $1 \leq i \leq m$, and $\overline{w}(t) = \begin{bmatrix} \overline{w}_1^T(t) & \dots \\ \overline{w}_m^T(t) \end{bmatrix}^T$ then, for $j = 0, 1$,

$$\overline{w}(t_j) = \mathcal{R}_{[A_{\mathcal{B}}, B_{\mathcal{B}}]}^{-1} T_{\mathcal{B}}^{-1} x_j . \tag{2.3}$$

131 Then, applying the control law,

$$u(t) = F_{\mathcal{B}}x(t) + G_{\mathcal{B}}\mathfrak{D}(d/dt)f(t), \qquad (2.4)$$

to the system represented by (1.1), we get:

$$x(t) = T_{\mathcal{B}} \mathcal{R}_{[A_{\mathcal{B}}, B_{\mathcal{B}}]} \overline{w}(t) , \qquad (2.5)$$

133 with

$$(u(0), x(0)) = (u_0, x_0)$$
 and $(u(t_1), x(t_1)) = (u_1, x_1)$. (2.6)

Let us now propose the trajectories:⁵

$$f_{i}(t) = \begin{bmatrix} t^{2\kappa_{i}+1} & \cdots & t^{\kappa_{i}+1} \end{bmatrix} \mathfrak{a}_{i,1} + \begin{bmatrix} t^{\kappa_{i}} & \cdots & 1 \end{bmatrix} \mathfrak{a}_{i,0},$$

$$\mathfrak{a}_{i,1} = \begin{bmatrix} a_{i,2\kappa_{i}+1} & \cdots & a_{i,\kappa_{i}+1} \end{bmatrix}^{T} \in \mathbb{R}^{\kappa_{i}+1}, \text{ and } \mathfrak{a}_{i,0} = \begin{bmatrix} a_{i,\kappa_{i}} & \cdots & a_{i,0} \end{bmatrix}^{T} \in \mathbb{R}^{\kappa_{i}+1},$$

$$(2.7)$$

with $i \in \{1, ..., m\}$, and let us define the following auxiliary matrices:

⁵ Lewis (1986) did a similar proposition when he introduced a "fast" input in his reachability consideration.

$$X_{(i,0)}(t) = \begin{bmatrix} \kappa_i!/0! & 0 & \cdots & 0 & 0\\ (\kappa_i!/1!)t & (\kappa_i-1)!/0! & 0 \cdots & 0 & 0\\ \vdots & \vdots & \ddots & \ddots & \ddots\\ (\kappa_i!/\kappa_i!)t^{\kappa_i} & ((\kappa_i-1)!/(\kappa_i-1)!)t^{\kappa_i-1} & \cdots & (1!/1!)t & 0!/0! \end{bmatrix}, \quad (2.8)$$

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$$X_{(i,1)}(t) = \begin{bmatrix} ((2\kappa_i + 1)!/(\kappa_i + 1)!)t^{\kappa_i + 1} & \cdots & ((\kappa_i + 1)!/1!)t \\ \vdots & \ddots & \vdots \\ ((2\kappa_i + 1)!/(2\kappa_i + 1)!)t^{2\kappa_i + 1} & \cdots & ((\kappa_i + 1)!/(\kappa_i + 1)!)t^{\kappa_i + 1} \end{bmatrix}. \quad (2.9)$$

The following Lemma gives a selection of the coefficient vectors $\mathfrak{a}_{i,0}$ and $\mathfrak{a}_{i,1}$, $i \in \{1, ..., m\}$, for f in (2.7) to satisfy assumptions (2.2) and (2.3) of Lemma 139 1 (see Appendix B for the proof).

Lemma 2. For $i \in \{1, ..., m\}$, the determinants of the auxiliary matrices (2.8) and (2.9) satisfy

$$\det (X_{(i,0)}(t)) = \prod_{\ell=0}^{\kappa_i} \ell! \quad and \quad \det (X_{(i,1)}(t)) = t^{(\kappa_i + 1)^2} \prod_{\ell=0}^{\kappa_i} \ell!$$
 (2.10)

Moreover, if we select the coefficient vectors $\mathfrak{a}_{i,0}$ and $\mathfrak{a}_{i,1}$, as follows:

$$\mathfrak{a}_{i,0} = X_{(i,0)}^{-1}(0)v_0, \quad \mathfrak{a}_{i,1} = X_{(i,1)}^{-1}(t_1) \left(v_1 - X_{(i,0)}(t_1) \mathfrak{a}_{i,0} \right),
v_j = \left[\left((\underline{e}_m^i)^T G_{\mathcal{B}}^{-1} (u_j - F_{\mathcal{B}} x_j) \right)^T \quad \left(\mathcal{R}_{\left[A_{\mathcal{B}_i}, B_{\mathcal{B}_i} \right]}^{-1} P_i T_{\mathcal{B}}^{-1} x_j \right)^T \right]^T, \ j \in \{0, 1\},$$
(2.11)

143 where:

$$P_i = \begin{bmatrix} \underline{e}_n^{\hat{n}_i + 1} & \cdots & \underline{e}_n^{\hat{n}_i + \kappa_i} \end{bmatrix}^T, \quad \hat{n}_1 = 0 \quad and \quad \hat{n}_{i \ge 2} = \sum_{j=1}^{i-1} \kappa_j, \tag{2.12}$$

then the function f defined by (2.7) fulfills assumptions (2.2) and (2.3) of Lemma 1.

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Let us note from Lemma 2 that the proposed solutions only depend on
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     the set of reachability indices S_{\kappa}, and on the fixed final time t_1. Hence,
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     once S_{\kappa} and t_1 are given, the matrices X_{i,0}(0), X_{i,0}(t_1) and X_{i,1}(t_1) are uniquely
     determined. And thus, the values of \mathfrak{a}_{i,0} and \mathfrak{a}_{i,1} only depend on the boundary
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     points, (u_0, x_0) and (u_1, x_1), of the trajectory (u, x) \in \mathfrak{B}_{[A,B]}.
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          From the above observation, it is possible to track a given trajectory
151
     (\bar{u}, \bar{x}) \in \mathcal{C}^{\infty}(\mathbb{R}^+, \mathbb{R}^{m+n}), with a delayed time t_1. Indeed, we only need to fix a
152
     sampling time t_1 \in \mathbb{R}^{+*}, and to apply iteratively Lemma 1 with the settings
     (u_0, x_0) = (u(kt_1), x(kt_1)) and (u_1, x_1) = (\bar{u}(kt_1), \bar{x}(kt_1)).
154
          Otherwise written, in each sampling interval [kt_1, (k+1)t_1), we find a tra-
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     jectory (u, x) \in \mathfrak{B}_{[A,B]} \cap \mathcal{C}^{\infty}(\mathbb{R}^+ \cap [kt_1, (k+1)t_1), \mathscr{U} \times \mathscr{X}), such that (u(kt_1), x(kt_1))
156
     = (u_0, x_0) \text{ and } \lim_{\sigma \to t_1} (u(kt_1 + \sigma), x(kt_1 + \sigma)) = (u_1, x_1).
          We have proved in this way the following Theorem.
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Theorem 1. Let us consider an input/state system $\Sigma_{i/s} = (\mathbb{R}^+, \mathscr{U} \times \mathscr{X}, \mathfrak{B}_{[A,B]})$, represented by (1.1). If (1.4) is satisfied, then for any sequence $(\bar{u}_k, \bar{x}_k) \in \mathbb{R}^{m+n}$, $k \in \mathbb{Z}^+$, and a given sampling time $t_1 \in \mathbb{R}^{+*}$, there exists a control law $u \in \mathcal{C}^{\infty}(\mathbb{R}^+, \mathbb{R}^m)$, such that $(u(kt_1), x(kt_1)) = (\bar{u}_{k-1}, \bar{x}_{k-1})$.

3. IMPLICIT SYSTEMS

In this Section, we formalize the notion of *implicit system* following the behavioral point of view. For this, let us first state the following definition:

Definition 1. An implicit representation $\mathfrak{R}^{imp}(E, A, B)$ is called an input/descriptor system, when for all initial condition $x_0 \in \mathscr{X}_d$, there exists at least one
solution $(u, x) \in \mathcal{C}^{\infty}(\mathbb{R}^+, \mathscr{U} \times \mathscr{X}_d)$, such that $x(0) = x_0$. The input/descriptor

system is defined by the triple $\Sigma_{i/d} = (\mathbb{R}^+, \mathscr{U} \times \mathscr{X}_d, \mathfrak{B}_{[E,A,B]})$, with behavior:

$$\mathfrak{B}_{[E,A,B]} = \left\{ (u,x) \in \mathcal{C}^{\infty}(\mathbb{R}^+, \mathscr{U} \times \mathscr{X}) \middle| \quad \left[(E\frac{\mathrm{d}}{\mathrm{d}t} - A) \quad -B \right] \begin{bmatrix} x \\ u \end{bmatrix} = 0 \right\} \quad (3.1)$$

At this point, it is important to clarify what exactly means the sentence "there exists at least one solution". For this, we are going to recall hereafter the notions of existence of solution introduced by Geerts (1993) and Aubin & Frankowska (1991).

3.1. Existence of solution for every initial condition

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Following (Hautus, 1976) and (Hautus & Silverman, 1983), Geerts (1993) generalized the solvability results of (Geerts & Mehrmann, 1990). One advantage of this generalization is that the solvability is introduced in a very natural way, passing from the distributional framework (Schwartz, 1978) to the usual time domain with ordinary differential equations; this is precisely the starting point of the so called behavioral approach (Polderman & Willems, 1998), chosen in this paper.

Geerts (1993) considered the linear combinations of impulsive and smooth distributions, with μ coordinates, denoted by \mathscr{C}^{μ}_{imp} , as the signal sets. The set \mathscr{C}^{μ}_{imp} is a subalgebra and is also decomposed as $\mathscr{C}^{\mu}_{p-imp} \oplus \mathscr{C}^{\mu}_{sm}$, where $\mathscr{C}^{\mu}_{p-imp}$ and \mathscr{C}^{μ}_{sm} denote the subalgebras of pure impulses⁷ and smooth distributions⁸,

⁶ See also Polderman & Willems (1998) and Kuijper (1992b).

⁷ The unit element of this subalgebra is the Dirac delta distribution δ . Any linear combination of δ and its distributional derivatives $\delta^{(\ell)}$, $\ell > 1$, is called impulsive.

⁸ The set of regular distributions are distributions that are functions; namely piecewise continuous integrable, or measurable functions. In those papers, they assume that the regular distributions u(t) are smooth on $[0, \infty)$, *i.e.* that a function $v:[0, \infty) \to \mathbb{R}$ exists, ar-

- respectively (Schwartz, 1978). He introduced the following definitions for the distributional version of the *implicit representation* (1.7) $\mathfrak{R}_{dist}^{imp}(E,A,B)$: $p Ex = Ax + Bu + Ex_0$ (c.f. Definitions 3.1 and 4.1, Geerts, 1993)¹⁰.
- Definition 2. (Geerts, 1993) Given the solution set $S_C(x_0, u) := \{x \in \mathscr{C}^{n_d}_{imp} | pE A]x = Bu + Ex_0\}$, the implicit representation $\mathfrak{R}^{imp}_{dist}(E, A, B)$ is:
- C-solvable if $\forall x_0 \in \mathscr{X}_d \; \exists \; u \in \mathscr{C}_{imp}^m : S_C(x_0, u) \neq \emptyset$,
- C-solvable in the function sense if $\forall x_0 \in \mathscr{X}_d \exists u \in \mathscr{C}_{sm}^m : S_C(x_0, u) \cap \mathscr{C}_{sm}^n \neq \emptyset$.
- Given the "consistent initial conditions set" $\mathcal{I}_C := \{z_0 \in \mathscr{X}_d | \exists u \in \mathscr{C}_{\mathrm{sm}}^m \exists x \in S_C(z_0, u) \cap \mathscr{C}_{\mathrm{sm}}^{n_d} : x(0^+) = z_0 \}$, and the "weakly consistent initial conditions set" $\mathcal{I}_C^w := \{z_0 \in \mathscr{X}_d | \exists u \in \mathscr{C}_{\mathrm{sm}}^m \exists x \in S_C(z_0, u) \cap \mathscr{C}_{\mathrm{sm}}^{n_d} \}$, a point $x_0 \in \mathscr{X}_d$ is called C-consistent if $x_0 \in \mathcal{I}_C$, and weakly C-consistent if $x_0 \in \mathcal{I}_C^w$.
- Let us note that:
- (i) C-solvability is concerned with distributional solutions,
- (ii) C-solvability in the function sense is concerned with solutions only com-
- 200 posed by ordinary functions arbitrarily often differentiable,
- 201 (iii) the two notions of consistency, C-consistent and weakly C-consistent,
- lead to *smooth solutions*, namely with no impulsions, but

bitrarily often differentiable including at t = 0, such that u(t) = 0 for t < 0 and u(t) = v(t), for $t \ge 0$ (Hautus & Silverman, 1983). These distributions are identified as ordinary functions with support on \mathbb{R}^+ .

⁹ Ex_0 stands for $Ex_0 \delta$, $x_0 \in \mathscr{X}_d$ being the initial condition, and pEx stands for $\delta^{(1)} * Ex$ (* denotes convolution); if pEx is smooth and $E\dot{x}$ stands for the distribution that can be identified with the ordinary derivative Edx/dt, then $pEx = E\dot{x} + Ex_{0+}$.

¹⁰ He also considered the *B-free* case $\mathfrak{R}_{dist}^{imp}(E,A,f)$: $pEx = Ax + f + Ex_0$.

203 (iv) C-consistency avoids jumps at the origin, namely the smooth solutions

are continuous on the left, and

205 (iv) weakly C-consistent enables jumps at the origin, but they are piece-wise

206 continuous smooth solutions.

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Geerts (1993) characterized the existence of solutions for every initial condition in his Corollary 3.6, Proposition 4.2 and Theorem 4.5. Hereafter we summarize these results with their geometric equivalences.

Theorem 2. (Geerts, 1993) If [H2] is fulfilled, then

• $\mathfrak{R}_{dist}^{imp}(E,A,B)$ is C-solvable if and only if $[(\lambda E - A) - B]]$ is right invertible as a rational matrix, i.e. if and only if 11

$$E\mathcal{V}_{\mathcal{X}_d}^* + A\mathcal{S}_{\mathcal{X}_d}^* + \mathcal{B} = \mathcal{X}_{eq} . \tag{3.2}$$

• $\mathfrak{R}_{dist}^{imp}(E, A, B)$ is C-solvable in the function sense if and only if $\mathcal{I}_{C}^{w} = \mathscr{X}_{d}$, namely, if and only if $\operatorname{Im} E + A\mathscr{K}_{E} + \mathscr{B} = \mathscr{X}_{eq}$, i.e. if and only if 12

$$E\mathcal{V}_{\mathcal{X}_d}^* = \operatorname{Im} E . \tag{3.3}$$

• $\mathcal{I}_C = \mathscr{X}_d$ if and only if $\operatorname{Im} E + \mathscr{B} = \mathscr{X}_{eq}$ i.e. if and only if 13

$$E\mathcal{V}_{\mathcal{X}_d}^* + \mathcal{B} = \mathcal{X}_{eq} . \tag{3.4}$$

216 3.2. Existence of a viable solution

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In order to study the reachability for implicit systems, Frankowska (1990)
introduced the set-valued map (the set of all admissible velocities) $\mathbf{F}: \mathscr{X}_d \rightsquigarrow$ $\mathscr{X}_d, \mathbf{F}(x) = E^{-1}(Ax + \mathscr{B}) = \{v \in \mathscr{X} | Ev \in Ax + \mathscr{B}\}, \text{ and the differential inclusion}$

$$dx/dt \in \mathbf{F}(x)$$
, where $x(0) = x_0$, (3.5)

Frankowska (1990) showed that the solutions of (1.7) and the ones of (3.5) are the same. She also clarified the meaning of a viable solution and she characterized the largest subspace of such viable solutions.

Definition 3. (Frankowska, 1990, Aubin & Frankowska, 1991)

- An absolutely continuous function $x: \mathbb{R}^+ \to \mathscr{X}_d$ is called a *trajectory* of (3.5), if $x(0) = x_0$ and $dx/dt \in \mathbf{F}(x)$ for almost every $t \in \mathbb{R}^+$, that is to say, if there exists a measurable function $u: \mathbb{R}^+ \to \mathscr{U}$ such that $x(0) = x_0$ and Edx/dt = Ax + Bu, for almost every $t \in \mathbb{R}^+$.
 - Let \mathcal{K} be a subspace¹⁴ of \mathcal{X}_d . A trajectory x of (3.5) is called *viable* in \mathcal{K} , if $x(t) \in \mathcal{K}$ for all $t \geq 0$. The set of such trajectories is called

 $^{^{13}}$ Directly follows from (ALG–V) and [H2].

¹⁴ We restrict our discussion to subspaces of finite dimensional vector spaces. In (Frankowska, 1990) and in (Aubin & Frankowska, 1991) these definitions are stated in the more general framework of closed sets of normed vector spaces.

the set of viable solutions in \mathscr{K} . The subspace \mathscr{K} is called a viability domain of \mathbf{F} , if for all $x \in \mathscr{K} : \mathbf{F}(x) \cap \mathscr{K} \neq \emptyset$. The subspace \mathscr{K} is called the viability kernel of (3.5) when it is the largest viability domain of \mathbf{F} .

Theorem 3. (Aubin & Frankowska, 1991) The supremal (A, E, B)-invariant subspace contained in \mathcal{X}_d , $\mathcal{V}_{\mathcal{X}_d}^*$, is the viability kernel of \mathcal{X}_d for the set-valued map $\mathbf{F}: \mathcal{X}_d \leadsto \mathcal{X}_d$, $\mathbf{F}(x) = E^{-1}(Ax + \mathcal{B})$. Moreover, for all $x_0 \in \mathcal{V}_{\mathcal{X}_d}^*$ there exists a trajectory $x \in \mathcal{C}^{\infty}(\mathbb{R}^+, \mathcal{V}_{\mathcal{X}_d}^*)$ solution of (1.7) satisfying $x(0) = x_0$.

Frankowska (1990) called a singular system "strict" when the viability kernel coincides with the whole descriptor space \mathscr{X}_d , namely

$$\mathscr{V}_{\mathscr{X}_d}^* = \mathscr{X}_d . \tag{3.6}$$

In order to clarify ideas, let us extract from (Bonilla & Malabre, 1997, Section 2.1) the following result:

Result 1. There exists a subspace \mathscr{X}_1 such that:

$$\mathscr{X}_d = \mathscr{V}_{\mathscr{X}_d}^* \oplus \mathscr{X}_1, \quad \mathscr{X}_{eq} = (E\mathscr{V}_{\mathscr{X}_d}^* + \mathscr{B}) \oplus A\mathscr{X}_1, \quad and \quad \mathscr{X}_1 \approx A\mathscr{X}_1.$$
 (3.7)

Moreover, when projecting on \mathscr{X}_1 any trajectory $x \in C^{\infty}(\mathbb{R}^+, \mathscr{V}_{\mathscr{X}_d}^*)$ solution of (1.7), we always get a null trajectory.

Furthermore, for all $x_0 \in \mathcal{V}_{\mathscr{X}_d}^*$ there exists at least one trajectory $(u, x_\rho) \in \mathcal{C}^{\infty}(\mathbb{R}^+, \mathscr{U} \times \mathcal{V}_{\mathscr{X}_d}^*)$ solution of (1.7), satisfying $x_\rho(0) = x_0$.

PROOF OF RESULT 1. From the algorithm shown in (Fig. 1, Bonilla & Malabre, 1997) and from [**H2**], we get the geometric decompositions (3.7), and under these decompositions, (1.7) takes the following form:

$$\begin{bmatrix}
E_{\rho} & * \\
\hline
0 & \overline{X}_{\rho-1}
\end{bmatrix} \stackrel{d}{\underline{d}t} \begin{bmatrix} x_{\rho} \\
\overline{x}_{\rho-1}
\end{bmatrix} = \begin{bmatrix} A_{\rho} & 0 \\
\hline
0 & I_1 \end{bmatrix} \begin{bmatrix} x_{\rho} \\
\overline{x}_{\rho-1} \end{bmatrix} + \begin{bmatrix} B_{\rho} \\
\hline
0 \end{bmatrix} u, \quad (3.8)$$

where $x_{\rho} \in \mathscr{V}_{\mathscr{X}_{d}}^{*}$, $\bar{x}_{\rho-1} \in \mathscr{X}_{1}$, $I_{1} : \mathscr{X}_{1} \leftrightarrow A\mathscr{X}_{1}$ is an isomorphism, $\overline{X}_{\rho-1}$ is a nilpotent matrix (an upper triangular matrix with zeros in its diagonal). Then $\bar{x}_{\rho-1} \equiv 0$.

252 If we now apply the following geometric decompositions:

$$E\mathscr{V}_{\mathscr{X}_d}^* + \mathscr{B} = E\mathscr{V}_{\mathscr{X}_d}^* \oplus \mathscr{B}_C, \ \mathscr{B} = (\mathscr{B} \cap E\mathscr{V}_{\mathscr{X}_d}^*) \oplus \mathscr{B}_C, \ \mathscr{U} = B^{-1}E\mathscr{V}_{\mathscr{X}_d}^* \oplus B^{-1}\mathscr{B}_C,$$

$$(3.9)$$

where \mathscr{B}_C is some complementary subspace of $\mathscr{B} \cap E\mathscr{V}_{\mathscr{X}_d}^*$, we get for $\mathfrak{R}^{imp}(E_\rho, A_\rho, B_\rho)$ (recall (3.8)):

$$\left[\frac{\overline{E}_{\rho}}{0} \right] \frac{\mathrm{d}}{\mathrm{d}t} x_{\rho} = \left[\frac{\overline{A}_{\rho}}{\widehat{A}_{\rho}} \right] x_{\rho} + \left[\frac{\overline{B}_{\rho}}{0} \right] \left[\begin{array}{c} u_{1} \\ u_{2} \end{array} \right].$$
(3.10)

Since Im $\overline{E}_{\rho} = E \mathcal{V}_{\mathcal{X}_d}^*$, there exists $\overline{E}_{\rho}^r : E \mathcal{V}_{\mathcal{X}_d}^* \to \mathcal{V}_{\mathcal{X}_d}^*$ such that $\overline{E}_{\rho} \overline{E}_{\rho}^r = I$. Then, one solution of (3.10) is given by

$$x_{\rho}(t) = \exp\left(\overline{E}_{\rho}^{r}\overline{A}_{\rho}t\right)x_{0} + \int_{0}^{t} \exp\left(\overline{E}_{\rho}^{r}\overline{A}_{\rho}(t-\tau)\right)\overline{E}_{\rho}^{r}\overline{B}_{\rho}u_{1}(\tau)d\tau,$$

$$u_{2}(t) = -\widehat{A}_{\rho}x_{\rho}(t).$$

Thus, the subspaces $EV_{\mathcal{X}_d}^* + \mathcal{B} \subset \mathcal{X}_{eq}$ and $V_{\mathcal{X}_d}^* \subset \mathcal{X}_d$ characterize the set of all possible trajectories of (1.7) which are not identically zero for any input u. The projection of any trajectory solution of (1.7) on the quotient space $\mathcal{X}_d/V_{\mathcal{X}_d}^*$, in correspondence with the projection on $\mathcal{X}_{eq}/(EV_{\mathcal{X}_d}^* + \mathcal{B})$ for the equation space, results in an identically null function (see Bonilla & Malabre, 1995, Corollary 2.1). Let us note that when Assumption [H2] holds, the geometric conditions $EV_{\mathcal{X}_d}^* + \mathcal{B} = \mathcal{X}_{eq}$ and $V_{\mathcal{X}_d}^* = \mathcal{X}_d$ are equivalent 15.

From (ALG-V) and [**H2**]: $E\mathscr{V}_{\mathscr{X}_d}^* + \mathscr{B} = \mathscr{X}_{eq} \Rightarrow \mathscr{V}_{\mathscr{X}_d}^* = A^{-1}(E\mathscr{V}_{\mathscr{X}_d}^* + \mathscr{B}) = \mathscr{X}_d; \ \mathscr{V}_{\mathscr{X}_d}^* = \mathscr{X}_d$ $\Rightarrow \operatorname{Im} E = E\mathscr{V}_{\mathscr{X}_d}^* \& \operatorname{Im} A = \operatorname{Im} A \cap (E\mathscr{V}_{\mathscr{X}_d}^* + \mathscr{B}) \Rightarrow \mathscr{X}_{eq} = \operatorname{Im} E + \operatorname{Im} A + \mathscr{B} = E\mathscr{V}_{\mathscr{X}_d}^* + \mathscr{B}.$

3.3. Discussion about existence of solution

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which the distributional and time-domain frameworks lead to the same conclusions with respect to the shape of the resulting system's solution trajecto-267 ries (c.f. (3.4)) and (3.3), namely the resulting distributions are identified as 268 ordinary functions, with support on \mathbb{R}^+ , and the generalized derivatives can 260 be identified with ordinary derivatives. Also, it is well connected with the 270 viability discussion of Frankowska (1990) and Aubin & Frankowska (1991); indeed, a singular system is strict if and only if the consistent initial condition set \mathcal{I}_C coincides with the whole descriptor variable space \mathscr{X}_d (c.f. (3.6) 273 and (3.4), and recall Assumption [H2]). 274 Regarding the set of weakly consistent initial conditions Geerts (1993) 275 notes, in his abstract and conclusion, that the condition that this set equals to the whole state space (under the Assumption [H2]) is equivalent to the impulse controllability for regular systems (Cobb, 1984) (or controllability of 278 the infinite part in the sense of Verghese et al (1981)). This correspondence 279 has been generalized to non regular systems and one can note that the nowadays most commonly adopted definition for impulse controllability is the one cited by Ishihara & Terra (2001)¹⁶: a general singular system is impulse 282 controllable if for every initial condition there exists a smooth (impulse-free) 283 $control\ u(t)$, and a smooth (impulse-free, but with possible jumps, especially 284 at the origin) variable descriptor trajectory solution of the system.

An important contribution of Geerts (1993) is to give conditions under

More generally, one can verify that the paper of Geerts (1993) is the

¹⁶Notice that in this paper is stated that the definition comes from Geerts (1993).

main reference on solvability properties, consistency of initial conditions, the ability to find control such that no impulsive phenomenon appears (see for examples Hou & Müller (1999), Ishihara & Terra (2001), Hou (2004) and Zhang (2006)).

However, one should also cite Özçaldiran & Haliločlu (1993) who proved that there exists a pair of smooth distributions (without jumps), satisfying $\mathfrak{R}_{dist}^{imp}(E,A,B)$ if and only if $x(0_*) \in \mathscr{V}_{\mathscr{X}_d^*}$, namely $\mathscr{V}_{\mathscr{X}_d^*} = \mathscr{X}_d$ (see their Proposition 1.3), and Przyluski & Sosnowski (1994) who proved that the subspace $\mathscr{V}_{\mathscr{X}_d^*} + \mathscr{K}_E$ characterizes the set of initial conditions, for which there exists a pair of smooth distributions (with possible jumps) satisfying $\mathfrak{R}_{dist}^{imp}(E,A,B)$, namely $E\mathscr{V}_{\mathscr{X}_d^*} = \operatorname{Im} E$ (see their Proposition 1).

In Figure 1, we summarize all the above discussion.

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$$\begin{bmatrix}
\mathscr{V}_{\mathscr{X}_{d}^{*}} = \mathscr{X}_{d}
\end{bmatrix} \iff \begin{bmatrix}
\operatorname{Im} E + \mathscr{B} = \mathscr{X}_{eq} \\
\left(E\mathscr{V}_{\mathscr{X}_{d}^{*}} + \mathscr{B} = \mathscr{X}_{eq}\right)
\end{bmatrix} \implies \begin{bmatrix}
\operatorname{Im} E + A\mathscr{K}_{E} + \mathscr{B} = \mathscr{X}_{eq} \\
\left(E\mathscr{V}_{\mathscr{X}_{d}^{*}} = \operatorname{Im} E\right)
\end{bmatrix}$$
(a)
(b)
(c)

Figure 1: Connexions between the notions of existence of solution and impulse controllability (under the Asumption [**H2**]).

Fig. 1(a) is the condition of viable solution of Aubin & Frankowska (1991)
or smooth solution (without any jump) of Özçaldiran & Haliločlu (1993). Fig.
1(b) is the condition that the set of consistent initial condition equals the
whole space of Geerts (1993). Fig. 1(c) is the condition of C-solvability in the
function sense of Geerts (1993) or the condition of Przyluski & Sosnowski
(1994) that the set of initial conditions of smooth solutions (with possible
jumps) equals the whole space, or the impulse controllability condition of
Ishihara & Terra (2001), or the impulse-mode controllability with arbitrary

initial conditions of Hou (2004).

Finally let us note that if the notion of weakly consistent initial conditions
as defined by Geerts (1993) is associated to the notion of impulse controllability, the notion of consistent initial conditions as defined by Geerts (1993)
is associated to the notion of reachability of Frankowska (1990) (in the more
general non regular case) since the system must be strict to be reachable.
See also the controllability discussion found in Korotka et al (2011).

4. REACHABILITY FOR IMPLICIT SYSTEMS

For the case of implicit systems, Frankowska (1990) extended the classical reachability definition as follows.

Definition 4. (Frankowska, 1990) The implicit representation (1.7) is called reachable if for any pair of vectors $x_0, x_1 \in \mathcal{X}_d$ and for any pair of real numbers $t_1 > t_0 \ge 0$, there exists a trajectory $x(\cdot)$ solution of (1.7), such that $x(t_0) = x_0$ and $x(t_1) = x_1$.

Frankowska (1990) has established in her Theorem 4.4 that $\mathscr{R}_{\mathscr{X}_d}^*$ (see (1.8)) is the reachable space of implicit systems like (1.7), with E and A not necessarily square. Hereafter, we recall Corollary 2.4 of Aubin and Frankowska (1991) which is $ad\ hoc$ for our paper.

Theorem 4. (Aubin & Frankowska, 1991) For any $t_1 > 0$ and for a system like (1.7), with E and A not necessarily square, the reachable space of (1.7) at time t_1 from the initial descriptor variable x(0) is equal to $\mathscr{R}^*_{\mathscr{X}_d}$. Moreover, $\mathscr{R}^*_{\mathscr{X}_d}$ is the supremal subspace such that for all $x_0, x_1 \in \mathscr{R}^*_{\mathscr{X}_d}$ and $t_1 > 0$, there

- exists a trajectory $x \in \mathcal{C}^{\infty}(\mathbb{R}^+ \mathscr{R}^*_{\mathscr{X}_d})$ solution of (1.7) satisfying $x(0) = x_0$ and $x(t_1) = x_1$.
- In this Section we are interested in generalizing and solving Problem 1 332 in the case of an $input/descriptor\ system\ \Sigma_{i/d}=(\mathbb{R}^+,\mathscr{U}\times\mathscr{X}_d,\ \mathfrak{B}_{[E,A,B]}),$ with
- behavior (3.1).
- Problem 2. Let us consider a input/descriptor system $\Sigma_{i/d} = (\mathbb{R}^+, \mathscr{U} \times \mathscr{X}_d, \mathscr{X}_d)$
- 335 $\mathfrak{B}_{[E,A,B]}$) represented by (1.7), and with the behavior (3.1). Given (u_0,x_0) ,
- 336 $(u_1, x_1) \in B^{-1}E\mathscr{R}^*_{\mathscr{X}_d} \times \mathscr{R}^*_{\mathscr{X}_d}$ and $t_1 > 0$, find a trajectory $(u, x) \in \mathfrak{B}_{[E, A, B]}$, such
- that $(u(0), x(0)) = (u_0, x_0)$ and $(u(t_1), x(t_1)) = (u_1, x_1)$.
- For answering this question, we proceed as follows.
- 339 (i) We first apply some geometric decompositions to the subspaces \mathscr{X}_d and
- \mathscr{X}_{eq} , inspired by Proposition 2.2 of Aubin and Frankowska (1991); the aim
- of these decompositions is to point out a part of the implicit representation,
- more or less explicit, which is expressed as a state space representation.
- 343 (ii) We next show that such a state space representation is reachable in the
- 344 classical sense.
- 345 (iii) Finally, based on Section 2, we answer Problem 2.
- 346 4.1. State reachability
- The following Lemma is proved in Appendix C.
- Lemma 3. When $\mathscr{R}_{\mathscr{X}_d}^* = \mathscr{X}_d$, the implicit representation (1.7) can be re-
- stricted to $\mathscr{R}_{\mathscr{X}_d}^*$ in the domain, and to $A\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B}$ in the codomain.
- Moreover, the spaces $\mathscr{R}^*_{\mathscr{X}_d}$, \mathscr{B} , $A\mathscr{R}^*_{\mathscr{X}_d} + \mathscr{B}$ and \mathscr{U} can be decomposed as fol-
- $\text{351} \quad lows: \ \mathscr{R}_{\mathscr{X}_d}^* = \mathscr{R}_{\mathcal{C}} \oplus (\mathscr{R}_{\mathscr{X}_d}^* \cap \mathscr{K}_E), \ \mathscr{B} = (\mathscr{B} \cap E\mathscr{R}_{\mathscr{X}_d}) \oplus \mathscr{B}_{\mathcal{C}}, \ A\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B} = E\mathscr{R}_{\mathscr{X}_d}^* \oplus \mathscr{B}_{\mathcal{C}}$

and $\mathcal{U} = B^{-1}E\mathcal{R}_{\mathcal{X}_d}^* \oplus \mathcal{U}_{\mathcal{C}}$, where $\mathcal{R}_{\mathcal{C}}$ and $\mathcal{U}_{\mathcal{C}}$ are complementary subspaces such that $\mathcal{R}_{\mathcal{C}} \approx E\mathcal{R}_{\mathcal{X}_d}^*$ and $\mathcal{U}_{\mathcal{C}} = B^{-1}\mathcal{B}_{\mathcal{C}} \approx \mathcal{B}_{\mathcal{C}}$. Under these decompositions, the implicit representation (1.7), restricted to $\mathcal{R}_{\mathcal{X}_d}^*$ in the domain and to $A\mathcal{R}_{\mathcal{X}_d}^* + \mathcal{B}$ in the codomain, takes the following form:

$$\overline{E} dx/dt = \overline{A}x + \overline{B}u,$$

$$\overline{E} = \begin{bmatrix} I_{\mathcal{C}} & 0 \\ 0 & 0 \end{bmatrix}, \ \overline{A} = \begin{bmatrix} \overline{A}_{1,1} & \overline{A}_{1,2} \\ \overline{A}_{2,1} & \overline{A}_{2,2} \end{bmatrix}, \ \overline{B} = \begin{bmatrix} \overline{B}_1 & 0 \\ 0 & I_{\mathcal{U}_{\mathcal{C}}} \end{bmatrix},$$
(4.1)

where $I_{\mathcal{C}}: \mathscr{R}_{\mathcal{C}} \leftrightarrow E\mathscr{R}_{\mathscr{X}_d}^*$, and $I_{\mathscr{U}_{\mathcal{C}}}: \mathscr{U}_{\mathcal{C}} \leftrightarrow \mathscr{B}_{\mathcal{C}}$ are isomorphisms.

In order to locate the state reachability part of (4.1), let us first define the natural projections:

$$\begin{split} P_{\mathcal{C}}: \mathscr{R}_{\mathscr{X}_d}^* &\to \mathscr{R}_{\mathcal{C}} / / (\mathscr{R}_{\mathscr{X}_d}^* \cap \mathscr{K}_E), \ P_{\ell}: \mathscr{R}_{\mathscr{X}_d}^* \to (\mathscr{R}_{\mathscr{X}_d}^* \cap \mathscr{K}_E) / / \mathscr{R}_{\mathcal{C}}, \\ Q_1: \mathscr{U} &\to B^{-1} E \mathscr{R}_{\mathscr{X}_d}^* / / B^{-1} \mathscr{B}_{\mathcal{C}}, \ Q_2: \mathscr{U} \to B^{-1} \mathscr{B}_{\mathcal{C}} / / B^{-1} E \mathscr{R}_{\mathscr{X}_d}^*. \end{split}$$

Let us next apply to $\mathfrak{R}^{imp}(\overline{E}, \overline{A}, \overline{B})$ the reachability algorithm of $\overset{\circ}{O}$ zçaldiran (1985), $\overline{\mathscr{R}}^0 = \{0\}$, $\overline{\mathscr{R}}^{\mu+1} = \overline{E}^{-1} \left(\overline{A} \overline{\mathscr{R}}^{\mu} + (\overline{\mathscr{B}}_1 \oplus \mathscr{B}_{\mathcal{C}}) \right)$, whose limit is $\mathscr{R}^*_{\mathscr{X}_d}$; namely: $\overline{\mathscr{R}}^1 = I_{\mathcal{C}}^{-1} \overline{\mathscr{B}}_1 \oplus (\mathscr{R}^*_{\mathscr{X}_d} \cap \mathscr{K}_E) \text{ and } \overline{\mathscr{R}}^{\mu+1} = I_{\mathcal{C}}^{-1} \left(\overline{A}_{1,1} P_{\mathcal{C}} \overline{\mathscr{R}}^{\mu} + \operatorname{Im} \left[\overline{A}_{1,2} \overline{B}_1 \right] \right) \oplus (\mathscr{R}^*_{\mathscr{X}_d} \cap \mathscr{K}_E)$ 362 \mathscr{K}_E , for $\mu \geq 1$. We thus obtain $I_{\mathcal{C}} P_{\mathcal{C}} \overline{\mathscr{R}}^{\mu+1} = \overline{A}^{\mu}_{1,1} \operatorname{Im} \overline{B}_1 + \sum_{i=0}^{\mu-1} \overline{A}^i_{1,1} \operatorname{Im} \left[\overline{A}_{1,2} \overline{B}_1 \right]$,
which implies:

$$E\mathscr{R}_{\mathscr{X}_d}^* = \left\langle \overline{A}_{1,1} \mid \operatorname{Im}\left[\overline{A}_{1,2} \ \overline{B}_1\right] \right\rangle. \tag{4.2}$$

Thus, $(\overline{A}_{1,1}, [\overline{A}_{1,2} \overline{B}_1])$ is a state reachable pair.

365 4.2. Behavioral reachability

Given any initial condition $x_0 \in \mathscr{R}_{\mathscr{X}_d}^*$, the solution set of (4.1) is characterized by the following behavior

$$\mathfrak{B}_{[\overline{E},\overline{A},\overline{B}]} = \left\{ (u,x) \in \mathcal{C}^{\infty} \left(\mathbb{R}^{+}, \mathcal{U} \times \mathcal{R}_{\mathcal{X}_{d}}^{*} \right) \mid \exists x_{0} \in \mathcal{R}_{\mathcal{X}_{d}}^{*} \text{ s.t. } P_{\mathcal{C}}x(t) = \exp \left(I_{\mathcal{C}}^{-1} \overline{A}_{1,1} t \right) P_{\mathcal{C}}x_{0} + \int_{0}^{t} \exp \left(I_{\mathcal{C}}^{-1} \overline{A}_{1,1} (t-\tau) \right) I_{\mathcal{C}}^{-1} \left(\overline{A}_{1,2} P_{\ell} x(\tau) + \overline{B}_{1} Q_{1} u(\tau) \right) d\tau,$$

$$Q_{2}u(t) = -I_{\mathcal{U}_{\mathcal{C}}}^{-1} \left(\overline{A}_{2,1} P_{\mathcal{C}}x(t) + \overline{A}_{2,2} P_{\ell}x(t) \right) \right\},$$

$$(4.3)$$

which behavioral equations are

$$\frac{d}{dt}I_{\mathcal{C}}P_{\mathcal{C}}x = \overline{A}_{1,1}P_{\mathcal{C}}x + \begin{bmatrix} \overline{A}_{1,2} & \overline{B}_{1} \end{bmatrix} \begin{bmatrix} P_{\ell}x \\ Q_{1}u \end{bmatrix},$$

$$0 = \overline{A}_{2,1}P_{\mathcal{C}}x + \overline{A}_{2,2}P_{\ell}x + I_{\mathcal{B}_{\mathcal{C}}}Q_{2}u.$$
(4.4)

Let us note that

 $_{370}$ (i) the component $P_{\mathcal{C}}x$ is the part of the descriptor variable which needs a control law to reach the desired goal.

372 (ii) The component $P_{\ell}x$ is the free part of the descriptor variable which acts 373 as some kind of internal input variable, together with the component Q_1u 374 which is the effective external control input variable.

(iii) The component Q_2u of the external control variable must be equal to a component of the descriptor variable. This is because we have chosen a purely integral description. This part of the input corresponds to algebraic relationships linked with purely derivative actions.

From Lemmas 1 and 2, we get the following theorem which gives a solution to Problem 2.

Theorem 5. Consider the reachable part (4.4) of the implicit representation (1.7). Denote: $n = \dim (E\mathscr{R}_{\mathscr{X}_d}^*)$ and $m = \dim (\mathscr{R}_{\mathscr{X}_d}^* \cap \mathscr{K}_E) + \dim (B^{-1}E\mathscr{R}_{\mathscr{X}_d}^*)$.

Let $\{\kappa_1, \kappa_2, \dots, \kappa_m\} \subset \mathbb{Z}^+$ be the reachability indices of the pair $(\overline{A}_{1,1}, [\overline{A}_{1,2} \ \overline{B}_1])$,

with $\kappa_1 \geq \kappa_2 \geq \dots \geq \kappa_m \geq 1$ and $\kappa_1 + \kappa_2 + \dots + \kappa_m = n$.

Let the linear map $F: \mathscr{R}_{\mathcal{C}} \to (\mathscr{R}_{\mathscr{X}_d}^* \cap \mathscr{K}_E) \times (B^{-1}E\mathscr{R}_{\mathscr{X}_d}^*)$ and the isomorphisms $T: \mathscr{R}_{\mathcal{C}} \leftrightarrow \mathscr{R}_{\mathcal{C}}$ and $G: (\mathscr{R}_{\mathscr{X}_d}^* \cap \mathscr{K}_E) \times (B^{-1}E\mathscr{R}_{\mathscr{X}_d}^*) \leftrightarrow (\mathscr{R}_{\mathscr{X}_d}^* \cap \mathscr{K}_E) \times (B^{-1}E\mathscr{R}_{\mathscr{X}_d}^*)$ be such that the pair $(\overline{A}_{\mathcal{B}}, \overline{\Gamma}_{\mathcal{B}})$, where $\overline{A}_{\mathcal{B}} = (I_{\mathcal{C}}T)^{-1}(\overline{A}_{1,1} + [\overline{A}_{1,2} \overline{B}_1] F)$ 388 $(I_{\mathcal{C}}T)$ and $\overline{\Gamma}_{\mathcal{B}} = (I_{\mathcal{C}}T)^{-1}[\overline{A}_{1,2} \overline{B}_1] G$, is expressed in the Brunovsky canonical form (1.2). The reachability matrix $\mathcal{R}_{[\overline{A}_{\mathcal{B}_i}, \overline{\gamma}_{\mathcal{B}_i}]}$ is expressed in terms of the reachability matrices $\mathcal{R}_{[\overline{A}_{\mathcal{B}_i}, \overline{\gamma}_{\mathcal{B}_i}]}$ as in (2.1).

391 Let $x_0, x_1 \in \mathscr{R}_{\mathscr{X}_d}^*, Q_1u_0, Q_1u_1 \in B^{-1}E\mathscr{R}_{\mathscr{X}_d}^*, \text{ and } t_1 > 0 \text{ be given. If we apply}$

$$\begin{bmatrix} P_{\ell}x(t) \\ Q_1u(t) \end{bmatrix} = FP_{\mathcal{C}}x(t) + G\mathfrak{D}(d/dt)f(t) , \qquad (4.5)$$

where $f(t) \in \mathcal{C}^{\infty}(\mathbb{R}^+, \mathbb{R}^m)$ and $\mathfrak{D}(d/dt)$ are defined as in Lemmas 1 and 2, we get

$$P_{\mathcal{C}}x(t) = T\mathcal{R}_{\left[\overline{A}_{\mathcal{B}}, \overline{\Gamma}_{\mathcal{B}}\right]}\bar{w}(t) , \qquad (4.6)$$

394 and

$$(u(t_i), x(t_i)) = \left(\begin{bmatrix} Q_1 u_i \\ -I_{\mathscr{U}_C}^{-1} \begin{bmatrix} \overline{A}_{2,1} & \overline{A}_{2,2} \end{bmatrix} x_i \end{bmatrix}, x_i \right), i \in \{0, 1\}, t_0 = 0. \quad (4.7)$$

395 4.3. Comments on the reachability

For the general case of implicit systems, represented by (1.7) with E and A not necessarily square, Frankowska (1990) has been the first to give a functional interpretation of reachability. For this, she has used the Viability Theory. More precisely, she has shown that reachability is equivalent to finding a trajectory $x \in C^{\infty}(\mathbb{R}^+, \mathcal{X}_d)$ solution of (1.7), starting from the initial condition x_0 and reaching the desired x_1 in a given finite time t_1 , namely $x(0) = x_0$ and $x(t_1) = x_1$ (see Theorem 4). Moreover, Frankowska (1990) has shown that reachability is geometrically characterized by the well known

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reachable space \mathscr{R}^*_{\mathscr{X}_d}. Of course, \mathscr{R}^*_{\mathscr{X}_d} is contained in the viability kernel \mathscr{V}^*_{\mathscr{X}_d}.
    This guarantees the existence of at least one trajectory solution of (1.7),
405
    leaving from x_0. This is also clear from \mathscr{R}^*_{\mathscr{X}_d} = \mathscr{V}^*_{\mathscr{X}_d} \cap \mathscr{S}^*_{\mathscr{X}_d}.
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         One interesting thing found in the proof of (Aubin & Frankowska, 1991,
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    Proposition 2.2) was to put forward the importance of the state space rep-
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    resentation (4.4) of the implicit equation (1.7). This fact has enabled us to
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    apply systematically the results of the classical State Space Control The-
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    ory. More precisely, thanks to the reachability of the pair (\overline{A}_{1,1}, \lceil \overline{A}_{1,2} \overline{B}_1 \rceil)
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    (see (4.2)), it is possible to find trajectories f_i \in \mathcal{C}^{\infty}(\mathbb{R}^+, \mathbb{R}^1) (see (2.7), (2.11),
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     (2.12), (2.8), and (2.9) for synthesizing the control law (4.5) (see also (2.4))
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    which guarantees (4.7) (see Lemmas 1 and 2 and Theorem 5).
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         The aim of Theorem 5 was not to prove once more the sufficiency of The-
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    orem 4, but to interpret the reachability of (1.7) in the classical state space
    framework. This interpretation allows us to have a better understanding of
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    the existing mechanisms in the linear implicit systems reachability. Indeed,
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    there exist two control actions. The first one is due to the free variable P_{\ell x},
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    and another one is due to the control input Q_1u (see (4.3)). The control input
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    Q_2u is algebraically linked to the descriptor variable components, the state
    variable P_{\mathcal{C}}x and the free variable P_{\ell}x, by means of the algebraic restriction
     (4.4.b) (when it exists).
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For systems composed by infinite elementary divisors¹⁷, the matrix Q_1

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The Kronecker showed that any pencil $[\lambda E - A]$, $\lambda \in \mathbb{C}$, is strictly equivalent to a canonical matrix, composed by four kind of blocks: (i) finite elementary divisors (integral actions), e.g. $\begin{bmatrix} (\lambda - \alpha) & 1 \\ 0 & (\lambda - \alpha) \end{bmatrix}$, (ii) infinite elementary divisors (derivative actions), e.g.

is null and the square matrix Q_2 is invertible. In this case, the equations (4.5) and (4.6) describe the behavior of a system fed-back by the control law (4.4b). Indeed, from (4.5) and (4.6), we get:

$$x(t) = \begin{bmatrix} I \\ F \end{bmatrix} T \mathcal{R}_{\left[\overline{A}_{\mathcal{B}}, \overline{\Gamma}_{\mathcal{B}}\right]} \overline{w}(t) + \begin{bmatrix} 0 \\ G \end{bmatrix} \mathfrak{D}(d/dt) f(t).$$

428 And from (4.4b), (4.5) and (4.6), we have:

$$Q_2 u(t) = -I_{\mathscr{U}_{\mathcal{C}}}^{-1} \left(\left(\overline{A}_{2,1} + \overline{A}_{2,2} F \right) T \mathcal{R}_{\left[\overline{A}_{\mathcal{B}}, \overline{\Gamma}_{\mathcal{B}} \right]} \overline{w}(t) + \overline{A}_{2,2} G \mathfrak{D}(d/dt) f(t) \right).$$

It is remarkable that in the systems represented by *column minimal in*dices, it is possible to have reachable systems without any control. This phenomenon is possible because of the existence of the free variable $P_{\ell}x$, which acts as an internal control signal.

33 5. POLE ASSIGNMENT

One of the most important features of the reachability of a state space representation (1.1) is the complete assignability of the closed loop spectrum by means of a state feedback. This equivalence is no longer the case when dealing with implicit representations (1.7). For the implicit description case, a geometric condition has to be added in order to guarantee such a *pole* assignment ability. In the sequel we give geometric conditions, which enable

 $[\]begin{bmatrix} 1 & \lambda \\ 0 & 1 \end{bmatrix}$, (iii) column minimal indices (internal variable structure), e.g. $\begin{bmatrix} \lambda & 1 \end{bmatrix}$, and (iv) row minimal indices (internal behavioral restrictions), e.g. $\begin{bmatrix} \lambda \\ 1 \end{bmatrix}$; see Gantmacher (1977).

us to assign the closed loop spectrum of: (i) a reachable implicit description (1.7), and (ii) a reachable and observable implicit description with output equation, $\mathfrak{R}^{imp}(E,A,B,C)$:

$$Edx/dt = Ax + Bu, \quad y = Cx, \tag{5.8}$$

where $C: \mathscr{X}_d \to \mathscr{Y}$ is a linear map, and the linear space \mathscr{Y} is the output space.

At this point, it is useful to clarify what we mean by spectrum and observable part.

A. Spectrum. We distinguish between the finite spectrum, $\sigma_f(A, E) = \{\lambda \in \mathbb{C} \mid \exists w \neq 0 \text{ s.t. } Av = \lambda Ev \}$, and the infinite spectrum, $\sigma_{\infty}(E, A) = \{\mu \in \mathbb{C} \mid \exists w \neq 0 \text{ s.t. } Ew = \mu Aw \}$ (c.f. Gantmacher (1977), Wong (1974), Armentano (1986)); the elements of $\sigma_f(A, E)$ are called poles, and the elements of $\sigma_{\infty}(E, A)$ are called poles at infinity. Note that for the four kind of blocks of the Kronecker canonical form¹⁸: (i) $\sigma_f(A, E) = \emptyset$ and $\sigma_{\infty}(E, A) = \emptyset$ for its row minimal indices blocks, (ii) $\sigma_f(A, E) = \emptyset$ for its infinite elementary divisors blocks, (iii) card $\{\sigma_f(A, E)\} = \infty$ and card $\{\sigma_{\infty}(E, A)\} = \infty$ for its column minimal indices blocks.

B. Observable part. With respect to the observable part, let us recall that it was shown in (Bonilla & Malabre, 1995) that the third condition of Kuijper (1992a) $-\begin{bmatrix} sE-A \\ C \end{bmatrix}$ has full column rank for all $s \in \mathbb{C}$ – for getting a minimal implicit representation (among all externally equivalent representations

¹⁸ See footnote 17.

¹⁹ Two representations are called externally equivalent if the corresponding sets of all possible trajectories for the external variables, expressed in an input/output partition (u, y), are the same (Willems, 1983, Polderman & Willems, 1998).

of the same type), is equivalent to have a null unobservable space, namely: $\mathcal{N} = \{0\}$. Indeed, if we decompose the descriptor and equation spaces as: $\mathcal{X}_d = \mathcal{X}_{ob} \oplus \mathcal{N}$ and $\mathcal{X}_{eq} = \mathcal{W}_{ob} \oplus E \mathcal{N}$, where \mathcal{X}_{ob} and \mathcal{W}_{ob} are some complementary subspaces, (5.8) takes the following form:

$$\begin{bmatrix} E_{ob} & 0 \\ Z & E_{\mathcal{N}} \end{bmatrix} \stackrel{d}{=} \begin{bmatrix} x_{ob} \\ x_{\mathcal{N}} \end{bmatrix} = \begin{bmatrix} A_{ob} & 0 \\ X & A_{\mathcal{N}} \end{bmatrix} \begin{bmatrix} x_{ob} \\ x_{\mathcal{N}} \end{bmatrix} + \begin{bmatrix} B_{ob} \\ Y \end{bmatrix} u$$

$$y = \begin{bmatrix} C_{ob} & 0 \end{bmatrix} \begin{bmatrix} x_{ob} \\ x_{\mathcal{N}} \end{bmatrix}$$
(5.9)

And the implicit descriptions $\mathfrak{R}^{imp}(E,A,B,C)$ and $\mathfrak{R}^{imp}(E_{ob},A_{ob},B_{ob},C_{ob})$ are externally equivalents (c.f. Bonilla & Malabre, 1995, Theorem 2.1). The point we want to enlighten here is that, since $E_{\mathscr{N}}$ is epic, there then exists $E_{\mathscr{N}}^r$ such that $E_{\mathscr{N}}E_{\mathscr{N}}^r=I$, which implies that all the homogeneous trajectories of (5.9), beginning at any initial condition $\begin{bmatrix} 0 \\ x_0 \end{bmatrix} \in \mathscr{N}$, $x_{\mathscr{N}}(t) = \exp(E_{\mathscr{N}}^r A_{\mathscr{N}} t) x_0$, always remain inside $\mathscr{N} \subset \mathscr{K}_C$. Thus, like in the classical state representations, they are called unobservable trajectories; and since \mathscr{N} is the supremal (A, E) invariant subspace contained in \mathscr{K}_C with this property, $\mathfrak{R}^{imp}(E_{ob}, A_{ob}, B_{ob}, C_{ob})$ is called the observable part of $\mathfrak{R}^{imp}(E,A,B,C)$.

5.1. Pole Assignment for a Reachable Implicit Description

Theorem 6. (Bonilla & Malabre, 1993) Given an implicit system represented by (1.7), for every finite symmetric (with respect to the real line) set of complex numbers Λ of cardinality $\dim(\mathscr{X}_{\mathscr{X}_d}^*)$, there exists a proportional and derivative descriptor feedback $u = F_p x + F_d \mathrm{d}x/\mathrm{d}t$, such that $\sigma_f(A_{F_p}, E_{F_d}) = \Lambda$, if and only if

$$\mathscr{R}_{\mathscr{X}_d}^* = \mathscr{X}_d \,, \tag{5.10}$$

$$\dim(E\mathcal{V}_{\mathcal{X}_d}^* + \mathcal{B}) \ge \dim(\mathcal{V}_{\mathcal{X}_d}^*) . \tag{5.11}$$

Bonilla & Malabre (1993) named this property external reachability. In that paper, condition (5.11) is expressed in its equivalent form:

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$$\dim(\mathscr{B}/(\mathscr{B}\cap E\mathscr{V}_{\mathscr{X}_d}^*)) \ge \dim(\mathscr{V}_{\mathscr{X}_d}^*\cap \mathscr{K}_E) . \tag{5.12}$$

Let us note that the geometric condition (5.10) is the *reachability* condition of Frankowska (1990) (c.f. Theorem 4) and the geometric condition (5.11) is the *descriptor variable uniqueness* condition of Lebret (1991), namely the closed loop left invertibility property, which enables us to assign the poles by means of a proportional and derivative feedback.

Lemma 4. (Lebret, 1991) There exists a proportional and derivative descriptor feedback $u = F_p x + F_d dx/dt + v$, such that the fed-back implicit representation $\mathfrak{R}^{imp}(E_{F_d}, A_{F_p}, B)$ satisfies $\operatorname{Ker}(\lambda E_{F_d} - A_{F_p}) = \{0\}$ iff (5.11) is satisfied.

Let us also note that in the case of a strict singular system, the geometric condition (5.11) is translated to (c.f. (3.6), (3.4) and Fig. 1): $\dim(\mathscr{X}_{eq}) \geq \dim(\mathscr{X}_{d}).$ In other words, it is not possible to assign all the spectrum of an implicit system having one degree of freedom, as for example the
ones considered in (Bonilla & Malabre, 2003).

We have the following Corollary of Theorem 6, proved in Appendix D.

Corollary 1. Let the implicit representation (1.7) satisfy the geometric conditions (5.10) and (5.11). Then:

1. If $\mathscr{U}_{\mathcal{C}} = \{0\}$, the implicit representation (4.1) reduces to the following reachable state space representation $(\overline{\mathscr{B}}_1 = \operatorname{Im} \overline{B}_1)$:

$$dx/dt = \overline{A}_{1,1}x + \overline{B}_1u \quad with \quad \langle \overline{A}_{1,1} \mid \overline{\mathscr{B}}_1 \rangle = \mathscr{X}_d.$$
 (5.13)

2. If $\mathscr{U}_{\mathcal{C}} \neq \{0\}$, there exists a map $\overline{V}_{\ell} : \mathscr{R}_{\mathscr{X}_d}^* \cap \mathscr{K}_E \to \mathscr{B}_{\mathcal{C}}$ such that $\operatorname{Ker} \overline{V}_{\ell} = \{0\}$. Then, applying the proportional feedback

$$u = \begin{bmatrix} 0 & 0 \\ -I_{\mathcal{U}_{\mathcal{C}}}^{-1} \overline{A}_{2,1} & -I_{\mathcal{U}_{\mathcal{C}}}^{-1} (\overline{A}_{2,2} + \overline{V}_{\ell}) \end{bmatrix} x + v , \qquad (5.14)$$

we get

$$\begin{bmatrix} I_{\mathcal{C}} & 0 \\ 0 & 0 \end{bmatrix} dx/dt = \begin{bmatrix} \overline{A}_{1,1} & 0 \\ 0 & -I \end{bmatrix} x + \begin{bmatrix} \overline{B}_{1} & \overline{A}_{1,2} \overline{V}_{\ell}^{g} I_{\mathcal{U}_{\mathcal{C}}} \\ 0 & \overline{V}_{\ell}^{g} I_{\mathcal{U}_{\mathcal{C}}} \end{bmatrix} v, \quad (5.15)$$

where $\overline{V}_{\ell}^g: \mathscr{B}_{\mathcal{C}} \to \mathscr{R}_{\mathscr{X}_d}^* \cap \mathscr{K}_E$ is some left inverse of \overline{V}_{ℓ} , and

$$I_{\mathcal{C}}\mathscr{R}_{\mathcal{C}} = E\mathscr{R}_{\mathscr{X}_{d}}^{*} = \left\langle \overline{A}_{1,1} \mid \overline{\mathscr{B}}_{1} + \overline{A}_{1,2}(\mathscr{R}_{\mathscr{X}_{d}}^{*} \cap \mathscr{K}_{E}) \right\rangle \quad and \quad \mathscr{R}_{\mathscr{X}_{d}}^{*} \cap \mathscr{K}_{E} = \overline{V}_{\ell}^{g} I_{\mathscr{U}_{\mathcal{C}}} \mathscr{U}_{\mathcal{C}}.$$

$$(5.16)$$

Furthermore, applying the proportional and derivative feedback

$$u = \begin{bmatrix} 0 & 0 \\ -I_{\mathcal{U}_{\mathcal{C}}}^{-1} \overline{A}_{2,1} & -I_{\mathcal{U}_{\mathcal{C}}}^{-1} (\overline{A}_{2,2} + \overline{V}_{\ell}) \end{bmatrix} x + \begin{bmatrix} 0 & 0 \\ 0 & -I_{\mathcal{U}_{\mathcal{C}}}^{-1} \overline{V}_{\ell} \end{bmatrix} dx/dt + v ,$$
(5.17)

we get

$$dx/dt = \begin{bmatrix} \overline{A}_{1,1} & \overline{A}_{1,2} \\ 0 & 0 \end{bmatrix} x + \begin{bmatrix} \overline{B}_1 & 0 \\ 0 & \overline{V}_{\ell}^g I_{\mathscr{U}_c} \end{bmatrix} v, \tag{5.18}$$

505 with

$$\left\langle \begin{bmatrix} \overline{A}_{1,1} & \overline{A}_{1,2} \\ 0 & 0 \end{bmatrix} \middle| \operatorname{Im} \begin{bmatrix} \overline{B}_{1} & 0 \\ 0 & \overline{V}_{\ell}^{g} I_{\mathcal{U}_{\mathcal{C}}} \end{bmatrix} \right\rangle = \left\langle \overline{A}_{1,1} \middle| \overline{\mathcal{B}}_{1} + \overline{A}_{1,2} (\mathcal{R}_{\mathcal{X}_{d}}^{*} \cap \mathcal{K}_{E}) \right\rangle \oplus \mathcal{U}_{\mathcal{C}} = \mathcal{X}_{d}. \tag{5.19}$$

From this Corollary, we realize that with a proportional feedback, we can only modify the finite spectrum of $\overline{A}_{1,1} = R \overline{A}|_{(\mathscr{R}_{\mathscr{X}_d}^*/\mathscr{R}_{\mathscr{X}_d}^* \cap \mathscr{K}_E)}$, where $R: A\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B} \to E\mathscr{R}_{\mathscr{X}_d}^*/\mathscr{B}_{\mathscr{C}}$ is the natural projection. To assign all the finite spectrum of \overline{A} , we need a proportional and derivative feedback.

5.2. Pole Assignment for a Reachable and Observable Implicit Description

In this section, we are going to consider the reachability of the observable
part after feedback, of the implicit representation (5.8). For this, let us
recall that the supremal (A, E, B)-invariant subspace contained in Ker C, \mathcal{V}^* = $\sup\{\mathcal{V} \subset \mathcal{K}_C | A\mathcal{V} \subset E\mathcal{V} + \operatorname{Im} B\}$, that characterizes the biggest part of a given
implicit representation $\mathfrak{R}^{imp}(E, A, B, C)$, can be made unobservable by means
of a suitable proportional and derivative descriptor feedback (c.f. the early
Geometric Algorithms Section).

Given a proportional and derivative descriptor feedback $u = F_p^* x + F_d^* dx/dt$, where $(F_p^*, F_d^*) \in \mathbf{F}(\mathcal{V}^*)$, let us consider the *quotient implicit representation* $\mathfrak{R}^{imp}(E_*, A_*, B_*, C_*)$, where the linear applications E_* , A_* , B_* and C_* are the induced maps uniquely defined by

$$E_*\Phi = \Pi E_{F_d^*}, \quad A_*\Phi = \Pi A_{F_p^*}, \quad B_* = \Pi B \quad \text{, and} \quad C = C_*\Phi \; ,$$
 (5.20)

where $\Phi: \mathscr{X}_d \to \mathscr{X}_d/\mathscr{V}^*$ and $\Pi: E\mathscr{X}_d \to E\mathscr{X}_d/E_{F_d^*}\mathscr{V}^*$ are the canonical projections. In Appendix E, we prove the following Theorem.²⁰

Theorem 7. Given an implicit system represented by (5.8), for every symmetric (with respect to the real line) set of complex numbers Λ of cardinality dim $((\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{V}^*)/\mathscr{V}^*)$, there exists a proportional and derivative descriptor

feedback $u = F_p^*x + F_d^*dx/dt + v$, with $(F_p^*, F_d^*) \in \mathbf{F}(\mathscr{V}^*)$, such that $\sigma_f(A_*, E_*) = \Lambda$,

where E_* and A_* are the induced maps (5.20), if and only if:

$$(\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{V}^*)/\mathscr{V}^* = \mathscr{X}_d/\mathscr{V}^*, \tag{5.21}$$

$$\dim\left((E\mathscr{V}_{\mathscr{X}_d}^* + \mathscr{B})/(E\mathscr{V}^* + \mathscr{B})\right) + \dim(\mathscr{B}) \ge \dim\left(\mathscr{V}_{\mathscr{X}_d}^*/\mathscr{V}^*\right). \tag{5.22}$$

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²⁰For a related result for regular systems see Schumacher (1980).

Let us note that (5.22) is equivalent to:²¹

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$$\dim\left(\mathscr{B}/(\mathscr{B}\cap E\mathscr{V}_{\mathscr{X}_d}^*)\right) \ge \dim\left(\mathscr{V}_{\mathscr{X}_d}^*\cap\mathscr{K}_E\right) - \dim\left(\mathscr{V}^*\cap E^{-1}\mathscr{B}\right). \tag{5.23}$$

For the implicit representations (5.8), satisfying Theorem 7, we will say that
they have the externally reachable output dynamics property²². Theorem 7
is important because it enables us to tackle systems having an internal variable structure (see for example Bonilla & Malabre (1991), Bonilla & Malabre
(2003), and Bonilla & Malabre (2008)). Let us also note that the geometric condition (5.22) is the descriptor variable uniqueness property notion of
Lebret (1991), namely the closed loop left invertibility property of the observable part of the system.

Lemma 5. (Lebret, 1991) There exists a proportional and derivative descriptor feedback $u = F_p x + F_d dx/dt + v$, such that the fed-back implicit representation $\mathfrak{R}^{imp}(E_{F_d}, A_{F_p}, B)$ satisfies $\operatorname{Ker}(\lambda E_{F_d} - A_{F_p}) \subset \mathscr{N}_{(F_p, F_d)}$ iff (5.22) is satisfied.

Let us finally note that, when comparing (5.22) with (5.11), we realize that Theorem 7 is indeed establishing the external reachabilty of the observable part after feedback. Also note that in the case $\mathcal{V}^* = \{0\}$, (5.22) and

²¹ This equivalence follows the between from equivalence (5.11)and (5.12),and from the fact that $\mathscr{B} \cap E\mathscr{V}^* = E(\mathscr{V}^* \cap E^{-1}\mathscr{B})$ implies that $\dim (\mathscr{V}^* \cap E^{-1}\mathscr{B}) = \dim (\mathscr{V}^*) + \dim (\mathscr{B} \cap \operatorname{Im} E) - \dim (E\mathscr{V}^* + \mathscr{B} \cap \operatorname{Im} E).$

²² The externally reachable output dynamics notion is a simplification of the one of reachable with output dynamics assignment (see Bonilla et al, 1994, Definition 6).

546 (5.11) are the same; and in the case $\mathscr{V}^* = \mathscr{V}^*_{\mathscr{X}_d}$, we get the trivial condition $\dim(\mathscr{B}) \geq 0$.

Let us finish this Section with an academic example.

Academic Example. Let us consider a perturbed linear system represented by the state space representation, $\mathfrak{R}^{ss}(\overline{A}, [\overline{B} \ \overline{S}], \overline{C})$:

$$d\bar{x}/dt = \overline{A}\bar{x} + \begin{bmatrix} \overline{B} & \overline{S} \end{bmatrix} \begin{bmatrix} u \\ q \end{bmatrix}$$
 and $y = \overline{C}\bar{x}$, (5.24)

where $q \in \mathcal{Q} \approx \mathbb{R}^{\eta}$, $u \in \mathcal{U} \approx \mathbb{R}^{m}$, $y \in \mathcal{Y} \approx \mathbb{R}^{p}$ and $\bar{x} \in \overline{\mathcal{X}} \approx \mathbb{R}^{\bar{n}}$, are the disturbance, the input, the output, and the state variables, respectively. We assume that the three following assumptions hold true:

554 $[\mathbf{H1}] \operatorname{Ker} \overline{B} = \{0\} \text{ and } \operatorname{Ker} \overline{S} = \{0\},$

555 $[\mathbf{H2}] \ q(\cdot) \in \mathcal{C}^m(\mathbb{R}^+, \mathcal{Q}), \ q(t), \ \mathrm{d}q(t)/\mathrm{d}t, \ \ldots, \ \mathrm{d}^m q(t)/\mathrm{d}t^m \in \mathcal{L}^{\infty}, \ \forall \ t \geq 0,$

 $_{556}$ [H3] q is a measured disturbance.

We want to solve the Disturbance Decoupling Problem with a PD Feedback (DDP-PDF).

Problem 3 (DDP-PDF). Under which conditions does there exist a proportional and derivative feedback $u = (\overline{F}_{p1} + \overline{F}_{d1}d/dt)\overline{x} + (\overline{F}_{p2} + \overline{F}_{d2}d/dt)q + v$, such that the closed-loop transfer function matrix between q and y is identically zero, and the finite spectrum of the observable part of the closed loop system is assigned at will.

For solving this problem, let us rewrite (5.24) in the descriptor form (5.8) with

$$E = \begin{bmatrix} I_{\bar{n}} & 0 \end{bmatrix}, \quad A = \begin{bmatrix} \overline{A} & \overline{S} \end{bmatrix}, \quad B = \begin{bmatrix} \overline{B} \end{bmatrix}, \quad C = \begin{bmatrix} \overline{C} & 0 \end{bmatrix}, \tag{5.25}$$

where $x = \begin{bmatrix} \overline{x}^T & q^T \end{bmatrix}^T \in \mathscr{X}_d = \overline{\mathscr{X}} \oplus \mathscr{Q} \approx \mathbb{R}^{\overline{n}+\eta}$ and $\mathscr{X}_{eq} = \overline{\mathscr{X}} \approx \mathbb{R}^{\overline{n}}$. In this implicit representation, the perturbation q is acting as the free part of the descriptor variable x. Then from Theorem 7, the **DDP-PDF** is solvable if and only if the implicit representation (5.8) and (5.25) satisfies (5.21) and (5.22), namely if and only if both following conditions hold true (see Appendix F):

$$\langle \overline{A} \mid \operatorname{Im} \left[\overline{B} \ \overline{S} \right] \rangle + \overline{\mathscr{V}}_{\left[\overline{B} \ \overline{S} \right]}^* = \overline{\mathscr{X}},$$
 (5.26)

 $\dim\left(\overline{\mathcal{V}}_{\left[\overline{B}\ \overline{S}\right]}^* \cap \overline{\mathcal{B}}\right) \ge \dim\left(\frac{\operatorname{Im}\overline{S}}{\operatorname{Im}\overline{S} \cap \left(\overline{\mathcal{V}}_{\left[\overline{B}\ \overline{S}\right]}^* + \overline{\mathcal{B}}\right)}\right). \tag{5.27}$

Let us consider for example: $\overline{A} = T_u \{\underline{e}_3^2\}$, $\overline{S} = a\underline{e}_3^1 + b\underline{e}_3^2$, with $|a| + |b| \neq 0$, $\overline{B} = \underline{e}_3^3$ and $\overline{C} = (\underline{e}_3^1)^T$. We have for this case $\operatorname{Im} \overline{S} = \operatorname{span} \{a\underline{e}_3^1 + b\underline{e}_3^2\}$, $\overline{\mathscr{B}} = \operatorname{span} \{\underline{e}_3^3\}$, and $\operatorname{Im} [\overline{B} \overline{S}] = \operatorname{span} \{a\underline{e}_3^1 + b\underline{e}_3^2, \underline{e}_3^3\}$, then $\langle \overline{A} | \operatorname{Im} [\overline{B} \overline{S}] \rangle = \operatorname{span} \{\underline{e}_3^1, \underline{e}_3^2, \underline{e}_3^3\} = \overline{\mathscr{X}}$, $\overline{\mathscr{V}}_{[\overline{B} \overline{S}]}^* = \operatorname{span} \{\underline{e}_3^2, \underline{e}_3^3\}$, and $\operatorname{Im} \overline{S} \cap (\overline{\mathscr{V}}_{[\overline{B} \overline{S}]}^* + \overline{\mathscr{B}}) = \{0\}$. Therefore (5.26) and (5.27) are satisfied, and the DDP-PDF has solution. Indeed, applying to (5.24) and (5.25) the PD feedback

$$u = \begin{bmatrix} -1 & 0 & 1 \end{bmatrix} d\bar{x}/dt + \begin{bmatrix} -1/\tau & 0 & 0 \end{bmatrix} \bar{x} + \begin{bmatrix} 1/\tau \end{bmatrix} v,$$

we obtain the closed loop system described by:

$$\tau dy/dt + y = v, \ \bar{x}_1 = y, \ \bar{x}_2 = dy/dt - aq, \ \text{ and } \ \bar{x}_3 = d^2y/dt^2 - adq/dt - bq.$$

Let us note that $\overline{\mathscr{V}}_{\overline{B}}^* = \{0\}$, and that $\operatorname{Im} \overline{S} \cap \left(\overline{\mathscr{V}}_{\overline{B}}^* + \overline{\mathscr{B}}\right) = \{0\}$, so there is no purely proportional solutions (see for example Wonham (1985)).

6. CONCLUDING REMARKS

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The notion of reachability introduced by Frankowska (1990) generalizes the property introduced by Yip & Sincovec (1981) in the regular case. Furthermore, Cobb (1984) indicates that this last property is consistent with

that of Rosenbrock (1974) introduced in a purely structural framework. In the same paper, Cobb (1984) enlightens with time domain characterizations 579 the difference between the reachability in the sense of Rosenbrock (1974) and the reachability in the sense of Verghese et al (1981) based, once again, 581 on pure structural tools (Kronecker canonical forms and/or Smith canoni-582 cal forms). In the regular case, for which the system can be decomposed 583 into two parts, a finite or slow subsystem, and an infinite or fast subsys-584 tem, Cobb (1984) showed that Rosenbrock (1974) reachability is equivalent to the reachability of the finite part and controllability of the infinite part. He also showed that Verghese et al (1981) reachability is equivalent to the 587 reachability of the finite part associated to the impulse controllability of the 588 infinite part. The impulse controllability as defined by Cobb (1984), or the controllability of the infinite part in the sense of Verghese et al (1981) is not any more defined by the idea to reach a desired descriptor variable but by the 591 ability of the system to generate a maximal class of impulses using piecewise 592 smooth, non impulsive controls. 593

One can deduce from this analysis that if a regular system is reachable (reachability of the finite and controllability of the infinite part) in the sense of Cobb (1984), Yip & Sincovec (1981), Rosenbrock (1974) or Frankowska (1990) (the four notion are equivalent in this case) then any vector is a consistent initial condition in the sense of Geerts (1993). The converse implication is not true. In general, reachability is not a consequence of the fact that every vector of the descriptor space defines a consistent initial condition. The condition is necessary but not sufficient for reachability.

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In this paper we have given a geometric interpretation of the implicit sys-

tems reachability Theorem of Frankowska (1990) and we have also found some interesting connections between the works (Frankowska, 1990) and (Geerts, 1993). The geometric interpretation has enabled us to have a better understanding of the existing mechanisms in the linear implicit systems reachability. For this, we have first interpreted the *viability* notion from a geometric point of view. We have next solved Problem 2, with Theorem 5, which is a generalization of Problem 1, solved with Theorem 1.

We have also studied the existing relationships, between the reachabil-610 ity property and the capability of the complete pole assignment ability. In 611 Theorem 6, we have considered the pole assignment problem of a reachable 612 implicit description, $\mathfrak{R}^{imp}(E,A,B)$; we have also shown in Corollary 1, that 613 with a proportional feedback, we can only modify the spectrum of the re-614 striction to $\mathscr{R}_{\mathscr{X}_d}^*/\mathscr{R}_{\mathscr{X}_d}^* \cap \mathscr{K}_E$ in the domain and $E\mathscr{R}_{\mathscr{X}_d}^*$ in the co-domain; to assign all the spectrum, we need a proportional and derivative feedback. In 616 Theorem 7, we have considered the pole assignment problem of a reachable 617 and observable implicit description with output equation, $\mathfrak{R}^{imp}(E_*, A_*, B_*, C_*)$. 618

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738 Appendix A. Proof of Lemma 1

For the existence of such $F_{\mathcal{B}}$, $T_{\mathcal{B}}$ and $G_{\mathcal{B}}$, see for example Theorems 5.9 and 5.10 and Corollary 5.3 of Wonham (1985). Doing the change of state variable: $T_{\mathcal{B}}^{-1}x = \xi = \begin{bmatrix} \xi_1^T & \cdots & \xi_m^T \end{bmatrix}^T$, we obtain the following set of closed loop state space representations (see (1.1), (1.2), and (2.4)): $\mathrm{d}\xi_i/\mathrm{d}t = A_{\mathcal{B},i}\xi_i + b_{\mathcal{B},i}\mathrm{d}^{\kappa_i}f_i/\mathrm{d}t^{\kappa_i}$, $i \in \{1,\ldots,m\}$, which solutions are (integrate by parts n_i times each solution): $\xi_i(t) = \exp\left(A_{\mathcal{B},i}t\right)\xi_i(0) + \int_0^t \exp\left(A_{\mathcal{B},i}(t-\tau)\right)b_{\mathcal{B},i}\frac{\mathrm{d}^{\kappa_i}f_i(\tau)}{\mathrm{d}\tau^{\kappa_i}}\mathrm{d}\tau$ $= \exp\left(A_{\mathcal{B},i}t\right)\left(\xi_i(0) - \sum_{j=0}^{\kappa_i-1}A_{\mathcal{B},i}^jb_{\mathcal{B},i}\frac{\mathrm{d}^{\kappa_i-(j+1)}f_i(0)}{\mathrm{d}t^{\kappa_i-(j+1)}}\right) + \sum_{j=0}^{\kappa_i-1}A_{\mathcal{B},i}^jb_{\mathcal{B},i}\frac{\mathrm{d}^{\kappa_i-(j+1)}f_i(t)}{\mathrm{d}t^{\kappa_i-(j+1)}}$ $= \exp\left(A_{\mathcal{B},i}t\right)\left(\xi_i(0) - \mathcal{R}_{\left[A_{\mathcal{B}_i},b_{\mathcal{B}_i}\right]}\overline{w}_i(0)\right) + \mathcal{R}_{\left[A_{\mathcal{B}_i},b_{\mathcal{B}_i}\right]}\overline{w}_i(t), \quad i \in \{1,\ldots,m\},$ $\xi(t) = \exp\left(A_{\mathcal{B}}t\right)\left(\xi(0) - \mathcal{R}_{\left[A_{\mathcal{B},B_{\mathcal{B}}\right]}\overline{w}}(0)\right) + \mathcal{R}_{\left[A_{\mathcal{B},B_{\mathcal{B}}\right]}\overline{w}}(t),$ $x(t) = \exp\left((A + BF_{\mathcal{B}})t\right)\left(x(0) - T_{\mathcal{B}}\mathcal{R}_{\left[A_{\mathcal{B},B_{\mathcal{B}}\right]}\overline{w}}(0)\right) + T_{\mathcal{B}}\mathcal{R}_{\left[A_{\mathcal{B},B_{\mathcal{B}}\right]}\overline{w}}(t).$ (A1)

Therefore, (A1), (2.2) and (2.3) imply (2.5) and (2.6). \Box

745 Appendix B. Proof of Lemma 2

Let us first compute det $(X_{(i,1)}(t))$, for $\kappa_{i} \geq 2$. For this, we first do the decomposition $X_{(i,1)}(t) = D_{i,\ell}(t)\widetilde{X}_{(i,\kappa_{i}+1)}D_{i,r}(t), i \in \{1,\ldots,m\},$ where $D_{i,\ell}(t) = DM\{\frac{t}{(\kappa_{i}+1)!}, \cdots, t^{\kappa_{i}+1}/(2\kappa_{i}+1)!\}, D_{i,r}(t) = DM\{(2\kappa_{i}+1)!t^{\kappa_{i}}, \cdots, (\kappa_{i}+1)!\}$ and $\widetilde{X}_{(i,\kappa_{i}+1)} = \begin{bmatrix} (\kappa_{i}+1)!/(\kappa_{i}+1)! & \cdots & (\kappa_{i}+1)!/1! \\ \vdots & \cdots & \vdots \\ (2\kappa_{i}+1)!/(2\kappa_{i}+1)! & \cdots & (2\kappa_{i}+1)!/(\kappa_{i}+1)! \end{bmatrix}.$ (B1)

Defining the following column elementary matrices:

$$T_{i,1} = \left[\begin{array}{c|c} \underline{e}_{\kappa_{i}+1}^{1} & (\underline{e}_{\kappa_{i}+1}^{2} - \underline{e}_{\kappa_{i}+1}^{1}) & (\underline{e}_{\kappa_{i}+1}^{3} - \kappa_{i}\underline{e}_{\kappa_{i}+1}^{2}) & \cdots & (\underline{e}_{\kappa_{i}+1}^{\kappa_{i}+1} - 2\underline{e}_{\kappa_{i}+1}^{\kappa_{i}}) \end{array} \right], \quad T_{i,2} = \left[\begin{array}{c|c} \underline{e}_{\kappa_{i}+1}^{1} & \underline{e}_{\kappa_{i}+1}^{2} & (\underline{e}_{\kappa_{i}+1}^{3} - (\kappa_{i}+2)\underline{e}_{\kappa_{i}+1}^{2}) & (\underline{e}_{\kappa_{i}+1}^{4} - (\kappa_{i}+1)\underline{e}_{\kappa_{i}+1}^{3}) & \cdots & (\underline{e}_{\kappa_{i}+1}^{\kappa_{i}+1} - 4\underline{e}_{\kappa_{i}+1}^{\kappa_{i}}) \end{array} \right], \quad T_{i,2} = \left[\begin{array}{c|c} \underline{e}_{\kappa_{i}+1}^{1} & (\underline{e}_{\kappa_{i}+1}^{\kappa_{i}+1} - (2\kappa_{i}-1)\underline{e}_{\kappa_{i}+1}^{\kappa_{i}+1}) & (\underline{e}_{\kappa_{i}+1}^{\kappa_{i}+1} - (2\kappa_{i}-1)\underline{e}_{\kappa_{i}+1}^{\kappa_{i}+1}) & (\underline{e}_{\kappa_{i}+1}^{\kappa_{i}+1} - (2\kappa_{i})\underline{e}_{\kappa_{i}+1}^{\kappa_{i}}) \end{array} \right], \quad T_{i,2} = \left[\begin{array}{c|c} \underline{e}_{\kappa_{i}+1}^{1} & (\underline{e}_{\kappa_{i}+1}^{\kappa_{i}+1} - (2\kappa_{i})\underline{e}_{\kappa_{i}+1}^{\kappa_{i}}) & (\underline{e}_{\kappa_{i}+1}^{\kappa_{i}+1} - (2\kappa_{i})\underline{e}_{\kappa_{i}+1}^{\kappa_{i}}) \end{array} \right], \quad T_{i,2} = \left[\begin{array}{c|c} \underline{e}_{\kappa_{i}+1}^{1} & (\underline{e}_{\kappa_{i}+1}^{\kappa_{i}+1} - (2\kappa_{i})\underline{e}_{\kappa_{i}+1}^{\kappa_{i}}) & (\underline{e}_{\kappa_{i}+1}^{\kappa_{i}+1} - (2\kappa_{i})\underline{e}_{\kappa_{i}+1}^{\kappa_{i}}) \end{array} \right], \quad T_{i,2} = \left[\begin{array}{c|c} \underline{e}_{\kappa_{i}+1}^{1} & (\underline{e}_{\kappa_{i}+1}^{\kappa_{i}+1} - (2\kappa_{i})\underline{e}_{\kappa_{i}+1}^{\kappa_{i}}) & (\underline{e}_{\kappa_{i}+1}^{\kappa_{i}+1} - (2\kappa_{i})\underline{e}_{\kappa_{i}+1}^{\kappa_{i}}) \end{array} \right], \quad T_{i,2} = \left[\begin{array}{c|c} \underline{e}_{\kappa_{i}+1}^{1} & (\underline{e}_{\kappa_{i}+1}^{\kappa_{i}+1} - (2\kappa_{i})\underline{e}_{\kappa_{i}+1}^{\kappa_{i}}) & (\underline{e}_{\kappa_{i}+1}^{\kappa_{i}+1} - (2\kappa_{i})\underline{e}_{\kappa_{i}+1}^{\kappa_{i}}) \end{array} \right], \quad T_{i,2} = \left[\begin{array}{c|c} \underline{e}_{\kappa_{i}+1}^{1} & \underline{e}_{\kappa_{i}+1}^{1} & \underline{e}_{\kappa_{i}+1}^{1} & \underline{e}_{\kappa_{i}+1}^{1} & \underline{e}_{\kappa_{i}+1}^{1} & \underline{e}_{\kappa_{i}+1}^{1} \end{array} \right] \right]$$

$$\widetilde{X}_{(i,\kappa_{i}+1)} \prod_{j=1}^{\kappa_{i}} T_{i,j} = \begin{bmatrix}
0! & 0 & 0 & \cdots & 0 & 0 & 0 \\
1 & 1! & 0 & \cdots & 0 & 0 & 0 \\
\vdots & \vdots & \ddots & \ddots & \ddots & \vdots \\
1 & \prod_{\ell=\kappa_{i}}^{\kappa_{i}} \ell & \prod_{\ell=\kappa_{i}-1}^{\kappa_{i}} \ell & \cdots & \prod_{\ell=3}^{\kappa_{i}} \ell & \prod_{\ell=2}^{\kappa_{i}} \ell & \kappa_{i}!
\end{bmatrix}.$$
(B2)

which implies (2.10.b).

For the second statement, let us first note that (2.7)-(2.9), (2.2) and (2.3), imply:

$$X_{(i,1)}(t)\mathfrak{a}_{i,1} + X_{(i,0)}(t)\mathfrak{a}_{i,0} = \begin{bmatrix} d^{\kappa_i} f_i(t)/dt^{\kappa_i} \\ \overline{w}_i(t) \end{bmatrix},$$
(B3)

with $i \in \{1, ..., m\}$. And let us next note that (2.2) and (2.3) are equivalent to:

$$\frac{\mathrm{d}^{\kappa_i} f(t_j)}{\mathrm{d} t^{\kappa_i}} = (\underline{e}_m^i)^T G_{\mathcal{B}}^{-1}(u(t_j) - F_{\mathcal{B}} x(t_j)) \quad \text{and} \quad \overline{w}_i(t_j) = \mathcal{R}_{[A_{\mathcal{B}}, B_{\mathcal{B}}]}^{-1} P_i T_{\mathcal{B}}^{-1} x(t_j), \quad (B4)$$

with $i \in \{1, ..., m\}$ and $j \in \{0, 1\}$, and where $t_0 = 0$, $u(t_0) = u_0$, $u(t_1) = u_1$, $x(t_0) = u_0$, and $x(t_1) = x_1$. Therefore, (2.8)-(2.10), (B3) and (2.11) imply (B4).

761 Appendix C. Proof of Lemma 3

Let us first prove that the spaces \mathscr{X}_d , \mathscr{X}_{eq} and \mathscr{U} can be decomposed as follows:

$$\mathcal{X}_{d} = \mathcal{R}_{\mathcal{C}} \oplus (\mathcal{R}_{\mathcal{X}_{d}}^{*} \cap \mathcal{K}_{E}) \oplus \mathcal{X}_{2} \oplus \mathcal{X}_{1}, \quad \mathcal{X}_{eq} = E\mathcal{R}_{\mathcal{X}_{d}}^{*} \oplus \mathcal{B}_{\mathcal{C}} \oplus E\mathcal{X}_{2} \oplus A\mathcal{X}_{1},$$

$$\mathcal{U} = B^{-1}E\mathcal{R}_{\mathcal{X}_{d}}^{*} \oplus B^{-1}\mathcal{B}_{\mathcal{C}},$$
(C1)

764 where:

$$\mathcal{X}_{d} = \mathcal{V}_{\mathcal{X}_{d}}^{*} \oplus \mathcal{X}_{1}, \quad \mathcal{V}_{\mathcal{X}_{d}}^{*} = \mathcal{R}_{\mathcal{X}_{d}}^{*} \oplus \mathcal{X}_{2}, \quad \mathcal{R}_{\mathcal{X}_{d}}^{*} = \mathcal{R}_{\mathcal{C}} \oplus (\mathcal{R}_{\mathcal{X}_{d}}^{*} \cap \mathcal{K}_{E}),
\mathcal{X}_{eq} = (E\mathcal{V}_{\mathcal{X}_{d}}^{*} + \mathcal{B}) \oplus A\mathcal{X}_{1}, \quad E\mathcal{V}_{\mathcal{X}_{d}}^{*} + \mathcal{B} = (A\mathcal{R}_{\mathcal{X}_{d}}^{*} + \mathcal{B}) \oplus E\mathcal{X}_{2},
A\mathcal{R}_{\mathcal{X}_{d}}^{*} + \mathcal{B} = E\mathcal{R}_{\mathcal{X}_{d}}^{*} \oplus \mathcal{B}_{\mathcal{C}}, \quad \mathcal{B} = (\mathcal{B} \cap E\mathcal{R}_{\mathcal{X}_{d}}^{*}) \oplus \mathcal{B}_{\mathcal{C}}.$$
(C2)

765 And also:

$$\mathscr{R}_{\mathcal{C}} \approx E\mathscr{R}_{\mathscr{X}_{\mathcal{A}}}^*, \quad \mathscr{X}_2 \approx E\mathscr{X}_2, \quad \mathscr{X}_1 \approx A\mathscr{X}_1, \quad \mathscr{B}_{\mathcal{C}} \approx B^{-1}\mathscr{B}_{\mathcal{C}} = \mathscr{U}_{\mathcal{C}},$$
 (C3)

$$\mathscr{V}_{\mathscr{X}_d}^* \cap \mathscr{K}_E = \mathscr{R}_{\mathscr{X}_d}^* \cap \mathscr{K}_E, \quad \mathscr{B} \cap E \mathscr{V}_{\mathscr{X}_d}^* = \mathscr{B} \cap E \mathscr{R}_{\mathscr{X}_d}^*. \tag{C4}$$

766 **1.** From (1.8), (ALG-S) and (ALG-V), we get:

$$\mathscr{V}_{\mathscr{X}_d}^* = A^{-1}(E\mathscr{V}_{\mathscr{X}_d}^* + \mathscr{B}) \text{ and } \mathscr{R}_{\mathscr{X}_d}^* = \mathscr{V}_{\mathscr{X}_d}^* \cap E^{-1}(A\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B}).$$
 (C5)

Indeed:
$$\mathscr{V}_{\mathscr{X}_d}^* \cap E^{-1}\left(A\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B}\right) = \mathscr{V}_{\mathscr{X}_d}^* \cap E^{-1}\left(A\left(\left(A^{-1}\left(E\mathscr{V}_{\mathscr{X}_d}^* + \mathscr{B}\right)\right) \cap \mathscr{S}_{\mathscr{X}_d}^*\right) + \mathscr{B}\right)$$

$$768 \quad = \mathcal{V}_{\mathcal{X}_d}^* \cap E^{-1} \left(\left(E \mathcal{V}_{\mathcal{X}_d}^* + \mathcal{B} \right) \cap A \mathcal{S}_{\mathcal{X}_d}^* + \mathcal{B} \right) = \mathcal{V}_{\mathcal{X}_d}^* \cap E^{-1} \left(\left(E \mathcal{V}_{\mathcal{X}_d}^* + \mathcal{B} \right) \cap \left(A \mathcal{S}_{\mathcal{X}_d}^* + \mathcal{B} \right) \right)$$

$$\mathcal{S}_{\mathcal{X}_d} = \mathcal{V}_{\mathcal{X}_d}^* \cap \left(\mathcal{V}_{\mathcal{X}_d}^* + E^{-1}\mathcal{B}\right) \cap \mathcal{S}_{\mathcal{X}_d}^* = \mathcal{V}_{\mathcal{X}_d}^* \cap \mathcal{S}_{\mathcal{X}_d}^* = \mathcal{R}_{\mathcal{X}_d}^* \text{ (see also Özçaldiran, 1985, }$$

Malabre, 1987). From (C5) and Result 1, we get:

$$\mathcal{X}_{d} = \mathcal{V}_{\mathcal{X}_{d}}^{*} \oplus \mathcal{X}_{1}, \quad \mathcal{X}_{eq} = (E\mathcal{V}_{\mathcal{X}_{d}}^{*} + \mathcal{B}) \oplus A\mathcal{X}_{1},
\mathcal{V}_{\mathcal{X}_{d}}^{*} = \mathcal{R}_{\mathcal{X}_{d}}^{*} \oplus \mathcal{X}_{2}, \quad \mathcal{R}_{\mathcal{X}_{d}}^{*} = \mathcal{R}_{\mathcal{C}} \oplus (\mathcal{R}_{\mathcal{X}_{d}}^{*} \cap \mathcal{K}_{E}).$$
(C6)

771

$$E\mathscr{R}_{\mathscr{X}_d}^* = E\mathscr{V}_{\mathscr{X}_d}^* \cap (A\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B}) \text{ and } A\mathscr{R}_{\mathscr{X}_d}^* \subset A\mathscr{V}_{\mathscr{X}_d}^* \subset E\mathscr{V}_{\mathscr{X}_d}^* + \mathscr{B}.$$
 (C7)

2. From (C5.b), we get (C4.a), which implies together with (C6.c):

$$E\mathscr{V}_{\mathscr{X}_d}^* = E\mathscr{R}_{\mathscr{X}_d}^* \oplus E\mathscr{X}_2. \tag{C8}$$

Indeed, the direct sum comes from the fact that $\mathscr{X}_2 \cap \mathscr{K}_E \subset \mathscr{V}_{\mathscr{X}_d}^* \cap \mathscr{K}_E =$

774
$$\mathscr{R}_{\mathscr{X}_d}^* \cap \mathscr{K}_E$$
 implies that $(\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{X}_2) \cap \mathscr{K}_E = (\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{X}_2) \cap (\mathscr{V}_{\mathscr{X}_d}^* \cap \mathscr{K}_E) = (\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{X}_2)$

775
$$\cap (\mathscr{R}_{\mathscr{X}_d}^* \cap \mathscr{K}_E) = \mathscr{R}_{\mathscr{X}_d}^* \cap \mathscr{K}_E = \mathscr{R}_{\mathscr{X}_d}^* \cap \mathscr{K}_E + \mathscr{X}_2 \cap \mathscr{K}_E.$$

Moreover, since:
$$\mathscr{X}_2 \cap \mathscr{K}_E = (\mathscr{X}_2 \cap \mathscr{V}_{\mathscr{X}_d}^*) \cap \mathscr{K}_E = \mathscr{X}_2 \cap (\mathscr{V}_{\mathscr{X}_d}^* \cap \mathscr{K}_E) = \mathscr{X}_2$$

777
$$\cap (\mathscr{R}^*_{\mathscr{X}_d} \cap \mathscr{K}_E) = (\mathscr{X}_2 \cap \mathscr{R}^*_{\mathscr{X}_d}) \cap \mathscr{K}_E = \{0\}, \text{ we get: } \dim(E\mathscr{X}_2) = \dim(\mathscr{X}_2), \text{ thus}$$

778 (C3.b) follows.

779 **3.** From (C8) and (C7), we get:

$$E\mathscr{V}_{\mathscr{X}_d}^* + \mathscr{B} = (E\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B}) \oplus E\mathscr{X}_2. \tag{C9}$$

Indeed, since: $\{0\} = (E\mathscr{R}_{\mathscr{X}_d}^*) \cap (E\mathscr{X}_2) = E\mathscr{V}_{\mathscr{X}_d}^* \cap (A\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B}) \cap (E\mathscr{X}_2) = (A\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B}) \cap (E\mathscr{X}_2), \text{ we get: } E\mathscr{X}_2 \cap (E\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B}) \subset E\mathscr{X}_2 \cap (A\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B}) = \{0\}.$ Moreover, (C9), (C7) and (C8) imply:

$$E\mathscr{V}_{\mathscr{X}_{d}}^{*} + \mathscr{B} = (E\mathscr{R}_{\mathscr{X}_{d}}^{*} + \mathscr{B}) \oplus E\mathscr{X}_{2} = (E\mathscr{V}_{\mathscr{X}_{d}}^{*} \cap (A\mathscr{R}_{\mathscr{X}_{d}}^{*} + \mathscr{B}) + \mathscr{B}) \oplus E\mathscr{X}_{2}$$
$$= ((E\mathscr{V}_{\mathscr{X}_{d}}^{*} + \mathscr{B}) \cap (A\mathscr{R}_{\mathscr{X}_{d}}^{*} + \mathscr{B})) \oplus E\mathscr{X}_{2} = (A\mathscr{R}_{\mathscr{X}_{d}}^{*} + \mathscr{B}) \oplus E\mathscr{X}_{2}.$$
(C10)

⁷⁸³ 4. From (C7.a) and (3.9), there exist subspaces, $\mathcal{W}_{\mathcal{C}}$ and $\mathcal{B}_{\mathcal{C}}$, such that:

$$A\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B} = E\mathscr{R}_{\mathscr{X}_d}^* \oplus \mathscr{W}_{\mathcal{C}}, \quad \mathscr{B} = ((E\mathscr{R}_{\mathscr{X}_d}^*) \cap \mathscr{B}) \oplus \mathscr{B}_{\mathcal{C}}, \quad \mathscr{W}_{\mathcal{C}} \supset \mathscr{B}_{\mathcal{C}}$$
(C11)

From (C8), (C10), and (C11), we get: $E\mathscr{V}_{\mathcal{X}_d}^* + \mathscr{B} = (E\mathscr{R}_{\mathscr{X}_d}^* \oplus E\mathscr{X}_2) + \mathscr{B} =$ 785 $E\mathscr{R}_{\mathscr{X}_d}^* \oplus \mathscr{B}_{\mathcal{C}} \oplus E\mathscr{X}_2 = (A\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B}) \oplus E\mathscr{X}_2$, that is to say: $E\mathscr{R}_{\mathscr{X}_d}^* \oplus \mathscr{W}_{\mathcal{C}} = A\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B}$ 786 $\approx E\mathscr{R}_{\mathscr{X}_d}^* \oplus \mathscr{B}_{\mathcal{C}}$. Hence:

$$\mathcal{W}_{\mathcal{C}} = \mathcal{B}_{\mathcal{C}} \tag{C12}$$

- 5. From the geometric decompositions (C6), (C10), (C11), and (C12), the subspaces \mathscr{X}_d , \mathscr{X}_{eq} , and \mathscr{U} take the form (C1)-(C2).
- From (C2.c,a) and since: Ker $A \subset \mathcal{V}_{\mathcal{X}_d}^*$ and Ker $B = \{0\}$, we get (C3.a,c,d).
- 7. To prove (C4.b), note first that (C8) and (C9) imply $\mathscr{B} \cap E\mathscr{V}_{\mathscr{X}_d}^* = \mathscr{B} \cap (E\mathscr{R}_{\mathscr{X}_d}^* + E\mathscr{X}_2)$ and $(E\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B}) \cap E\mathscr{X}_2 = \{0\}$. Let $x \in \mathscr{B} \cap (E\mathscr{R}_{\mathscr{X}_d}^* + E\mathscr{X}_2)$, there then exist $z \in E\mathscr{R}_{\mathscr{X}_d}^*$, $y \in E\mathscr{X}_2$, and $b \in \mathscr{B}$ such that x = z + y = b,

which implies $y = b - z \in (E\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B}) \cap E\mathscr{X}_2 = \{0\}, i.e. \ x \in \mathscr{B} \cap E\mathscr{R}_{\mathscr{X}_d}^*$. Therefore: $\mathscr{B} \cap E\mathscr{V}_{\mathscr{X}_d}^* = \mathscr{B} \cap (E\mathscr{R}_{\mathscr{X}_d}^* + E\mathscr{X}_2) \subset \mathscr{B} \cap E\mathscr{R}_{\mathscr{X}_d}^* \subset \mathscr{B} \cap E\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B} \cap E\mathscr{X}_2 \subset \mathscr{B} \cap (E\mathscr{R}_{\mathscr{X}_d}^* + E\mathscr{X}_2) = \mathscr{B} \cap E\mathscr{V}_{\mathscr{X}_d}^*$.

Let us next note that under the geometric decompositions, (C1)-(C3), the implicit representation (1.7) takes the following form (recall (3.7) and (3.8)):

$$\begin{bmatrix} \overline{E} & 0 & * \\ 0 & I_2 & * \\ 0 & 0 & \overline{X}_{\rho-1} \end{bmatrix} \stackrel{d}{\underset{}}{dt} x = \begin{bmatrix} \overline{A} & \widehat{A} & 0 \\ 0 & \widehat{A}_3 & 0 \\ 0 & 0 & I_1 \end{bmatrix} x + \begin{bmatrix} \overline{B} \\ 0 \\ 0 \end{bmatrix} u, \quad (C13)$$

where $I_2: \mathscr{X}_2 \leftrightarrow E\mathscr{X}_2$ is an isomorphism, and the matrices \overline{E} , \overline{A} and \overline{B} , are the ones shown in (4.1). Then, when $\mathscr{R}_{\mathscr{X}_d} = \mathscr{X}_d$, we get (4.1).

800 Appendix D. Proof of Corollary 1

Let us first note that (5.10) implies that the implicit representation (C13) is only composed by the linear transformations (4.1).

Let us next note that Lemma 3 and (5.11) imply that (see (C2)-(C4)):

$$\mathcal{V}_{\mathcal{X}_d}^* \cap \mathcal{K}_E = \mathcal{R}_{\mathcal{X}_d}^* \cap \mathcal{K}_E \quad \text{and} \quad \mathcal{B}/(\mathcal{B} \cap E\mathcal{V}_{\mathcal{X}_d}^*) = \mathcal{B}/(\mathcal{B} \cap E\mathcal{R}_{\mathcal{X}_d}^*) \approx \mathcal{B}_{\mathcal{C}} \approx \mathcal{U}_{\mathcal{C}},$$

$$\tag{D1}$$

$$\dim(A\mathscr{R}_{\mathscr{X}_d}^* + \mathscr{B}) \ge \dim(\mathscr{R}_{\mathscr{X}_d}^*). \tag{D2}$$

Case 1. If $\mathscr{U}_{\mathcal{C}} = \{0\}$, then (5.12) and (D1) imply: $\mathscr{R}_{\mathscr{X}_d}^* \cap \mathscr{K}_E = \{0\}$. Thus, the blocks $\overline{A}_{1,2}$, $\overline{A}_{2,1}$, $\overline{A}_{2,2}$, and $I_{\mathscr{U}_{\mathcal{C}}}$ actually disappear from (4.1), corresponding to 0 row and 0 column. Moreover $\overline{B}_1 \neq 0$, because the pair $(\overline{A}_{1,1}, [\overline{A}_{1,2} \overline{B}_1])$ is reachable (see (4.2)). Namely, we get (5.13).

Case 2. The existence of \overline{V}_{ℓ} is implied by (D2). From (5.14) and (4.1), we get (5.15). From (5.15) and (4.2), we get (5.16). From (5.17) and (4.1), we get (5.18). From (5.18) and (4.2), we get (5.19).

Appendix E. Proof of Theorem 7

Let us first propose a PD descriptor feedback $u = F_p^* x + F_d^* dx/dt + v$, where the pair of linear transformations (F_p, F_d) is chosen such that:

$$(F_p^*, F_d^*) \in \mathbf{F}(\mathcal{V}^*) \text{ and } \mathscr{B} \cap E_{F_d^*} \mathcal{V}^* = \{0\}.$$
 (E1)

Let us next, consider the quotient implicit representation $\mathfrak{R}^{imp}(E_*, A_*, B_*, C_*)$ defined by (5.20). Let us note that $\operatorname{Ker} B_* = B^{-1}E_{F_d^*}\mathcal{V}^* \approx \mathcal{B} \cap E_{F_d^*}\mathcal{V}^*$ implies

Ker $B_* = \{0\}$, and that $\Phi \mathcal{N}_{(F_p^*, F_d^*)} = \Phi \sup \{\mathcal{V} \subset \operatorname{Ker} C \mid A_{F_p^*}\mathcal{V} \subset E_{F_d^*}\mathcal{V}\} = \Phi \mathcal{V}^* = \{0\}$ implies the observability of the quotient implicit representation $\mathfrak{R}^{imp}(E_*, A_*, B_*, C_*)$. The proof of Theorem 7 is done in 4 steps:

i) $\mathfrak{R}^{imp}(E_{F_d^*}, A_{F_p^*}, B, C)$ is externally equivalent to $\mathfrak{R}^{imp}(E_*, A_*, B_*, C_*)$. This fact follows from (Theorem 2.1, Bonilla & Malabre, 1995), which states, among others, the external equivalency between $\mathfrak{R}^{imp}(E_{F_d^*}, A_{F_p^*}, B, C)$ and $\mathfrak{R}^{imp}(E_*, A_*, B_*, C_*)$ (see also Kuijper & Schumacher, 1991).

ii) $\mathscr{V}_{\mathscr{X}_d/\mathscr{V}^*}^* = \Phi \mathscr{V}_{\mathscr{X}_d}^*$ and $\mathscr{S}_{\mathscr{X}_d/\mathscr{V}^*}^* = \Phi \mathscr{S}_{\mathscr{X}_d}^*$. For the case of the quotient implicit representation $\mathfrak{R}^{imp}(E_*, A_*, B_*, C_*)$ the corresponding algorithms (ALG-V) and (ALG-S), for computing $\mathscr{V}_{\mathscr{X}_d/\mathscr{V}^*}^*$ and $\mathscr{S}_{\mathscr{X}_d/\mathscr{V}^*}^*$, take the following form:

$$\mathcal{V}_{\mathcal{X}_d/\mathcal{V}^*}^0 = \mathcal{X}_d/\mathcal{V}_d^*, \quad \mathcal{V}_{\mathcal{X}_d/\mathcal{V}^*}^{\mu+1} = A_*^{-1}(E_*\mathcal{V}_{\mathcal{X}_d/\mathcal{V}^*}^{\mu} + \mathcal{B}_*),
\mathcal{S}_{\mathcal{X}_d/\mathcal{V}^*}^0 = \{0\}, \quad \mathcal{S}_{\mathcal{X}_d/\mathcal{V}^*}^{\mu+1} = E_*^{-1}(A_*\mathcal{S}_{\mathcal{X}_d/\mathcal{V}^*}^{\mu} + \mathcal{B}_*).$$
(E2)

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It is clear that: \mathscr{V}^0_{\mathscr{X}_d/\mathscr{V}^*} = \Phi\mathscr{V}^0_{\mathscr{X}_d} and \mathscr{S}^0_{\mathscr{X}_d/\mathscr{V}^*} = \Phi\mathscr{S}^0_{\mathscr{X}_d}. Let us assume that:
\mathcal{V}^{\mu}_{\mathcal{X}_d/\mathcal{V}^*} = \Phi \mathcal{V}^{\mu}_{\mathcal{X}_d} \quad \text{and} \quad \mathcal{S}^{\mu}_{\mathcal{X}_d/\mathcal{V}^*} = \Phi \mathcal{S}^{\mu}_{\mathcal{X}_d}, \quad \text{then from (E2) and from (5.20)}, \quad \text{we} \quad \mathcal{S}^{\mu}_{\mathcal{X}_d/\mathcal{V}^*} = \Phi \mathcal{S}^{\mu}_{\mathcal{X}_d}.
\text{get: } \mathcal{V}^{\mu+1}_{\mathcal{X}_d/\mathcal{V}^*} = (\Phi\mathcal{X}_d) \cap A^{-1}_* \left( E_* \Phi \mathcal{V}^{\mu}_{\mathcal{X}_d} + \Pi \mathscr{B} \right) = \Phi \Phi^{-1} A^{-1}_* \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right) = \Phi A_{F_d^*}^{-1} \Pi \left( E_{F_d^*} \mathcal{V}^{\mu}_{\mathcal{X}_d} + \mathscr{B} \right)
 \text{830} \quad \Pi^{-1}\Pi\left(E_{F_d^*}\mathscr{V}_{\mathscr{X}_d}^{\mu}+\mathscr{B}\right) = \Phi A_{F_d^*}^{-1}\left(E_{F_d^*}\mathscr{V}_{\mathscr{X}_d}^{\mu}+\mathscr{B}+E_{F_d^*}\mathscr{V}^*\right) = \Phi\left(\mathscr{X}_d\cap A_{F_d^*}^{-1}\left(E_{F_d^*}\mathscr{V}_{\mathscr{X}_d}^{\mu}+\mathscr{B}\right)\right) 
\text{831} = \Phi \mathscr{V}_{\mathscr{X}_d}^{\mu+1}, \text{ and } \mathscr{S}_{\mathscr{X}_d/\mathscr{V}^*}^{\mu+1} = (\Phi \mathscr{X}_d) \cap E_*^{-1} \left( A_* \left( (\Phi \mathscr{X}_d) \cap (\Phi \mathscr{S}_{\mathscr{X}_d}^{\mu}) \right) + \Pi \mathscr{B} \right) = \Phi \Phi^{-1} E_*^{-1}
832 \Pi\left(A_{F_d^*}\left(\mathscr{X}_d\cap\mathscr{S}_{\mathscr{X}_d}^{\mu}\right)+\mathscr{B}\right) = \Phi E_{F_*^*}^{-1}\Pi^{-1}\Pi\left(A_{F_d^*}\left(\mathscr{X}_d\cap\mathscr{S}_{\mathscr{X}_d}^{\mu}\right)+\mathscr{B}\right) = \Phi E_{F_*^*}^{-1}\left(A_{F_d^*}\left(\mathscr{X}_d\cap\mathscr{S}_{\mathscr{X}_d}^{\mu}\right)+\mathscr{B}\right)
833 \left(\mathscr{X}_d \cap \mathscr{S}^{\mu}_{\mathscr{X}_d}\right) + \mathscr{B} + E_{F_d^*}\mathscr{V}^*\right) = \Phi E_{F_d^*}^{-1}\left(A_{F_d^*}\left(\mathscr{X}_d \cap \mathscr{S}^{\mu}_{\mathscr{X}_d}\right) + \mathscr{B}\right) + \Phi \operatorname{Ker} E_{F_d^*} =
834 \Phi\left(\mathscr{X}_d \cap E_{F^*}^{-1}\left(A_{F_d^*} \left(\mathscr{X}_d \cap \mathscr{S}_{\mathscr{X}_d}^{\mu}\right) + \mathscr{B}\right)\right) = \Phi\mathscr{S}_{\mathscr{X}_d}^{\mu+1}.
835 iii) If (5.21) and (5.23) are satisfied, then \mathfrak{R}^{imp}(E_*, A_*, B_*, C_*) satisfies The-
836 orem 6. Since: (\mathcal{V}_{\mathcal{X}_d}^* + \mathcal{S}_{\mathcal{X}_d}^*) \cap \operatorname{Ker} \Phi = (\mathcal{V}_{\mathcal{X}_d}^* + \mathcal{S}_{\mathcal{X}_d}^*) \cap \mathcal{V}^* = \mathcal{V}^* = \mathcal{V}_{\mathcal{X}_d}^* \cap \mathcal{V}^* +
837 \mathscr{S}_{\mathscr{X}_d}^* \cap \mathscr{V}^*, we get from (5.21): \mathscr{R}_{\mathscr{X}_d/\mathscr{V}^*}^* = \mathscr{V}_{\mathscr{X}_d/\mathscr{V}^*}^* \cap \mathscr{S}_{\mathscr{X}_d/\mathscr{V}^*}^* = \Phi \mathscr{V}_{\mathscr{X}_d}^* \cap \Phi \mathscr{S}_{\mathscr{X}_d}^* = \Phi \mathscr{V}_{\mathscr{X}_d}^* \cap \Phi \mathscr{V}_{\mathscr{X}_d}^* = \Phi \mathscr{V}_{\mathscr{X}_d}^* \cap \Phi \mathscr{S}_{\mathscr{X}_d}^* = \Phi \mathscr{V}_{\mathscr{X}_d}^* \cap \Phi \mathscr{V}_{\mathscr{X}_d}^* = \Phi \mathscr{V}_{\mathscr{X}_d}^* \cap \Phi \mathscr{V}_{\mathscr{X}_d}^*
838 \Phi\left(\mathcal{V}_{\mathcal{X}_d}^*\cap\mathcal{S}_{\mathcal{X}_d}^*\right)=\Phi\mathcal{R}_{\mathcal{X}_d}^*=\Phi\left(\mathcal{R}_{\mathcal{X}_d}^*+\mathcal{V}^*\right)=\Phi\mathcal{X}_d=\mathcal{X}_d/\mathcal{V}^*, \text{ which is the first constants}
                 dition of Theorem 6. On the other hand, since for any F_d^*: \mathcal{X}_d \to \mathcal{U}, E^{-1}\mathcal{B} =
840 E_{F_d^*}^{-1}\mathscr{B}, we have: \dim(\mathscr{K}_E) + \dim(\operatorname{Im} E \cap \mathscr{B}) = \dim(E_{F_d^*}^{-1}\mathscr{B}), which together with
                         (5.23) imply:<sup>23</sup>
                                                             \dim(\mathscr{B}) \ge \dim\left(E_{F_*^*}^{-1}\mathscr{B}/(\mathscr{V}^* \cap E_{F_*^*}^{-1}\mathscr{B})\right) = \dim\left(\Phi E_{F_*^*}^{-1}\mathscr{B}\right) = \dim\left(E_*^{-1}\mathscr{B}_*\right),
                          then: \dim(\mathscr{B}_*) = \dim(\Pi\mathscr{B}) \geq \dim(E_*^{-1}\mathscr{B}_*) - \dim(\mathscr{B} \cap \operatorname{Ker}\Pi) = \dim(E_*^{-1}\mathscr{B}_*)
                          -\dim(B\operatorname{Ker} B_*) = \dim(E_*^{-1}\mathscr{B}_*) - \dim(\operatorname{Ker} B_*) = \dim(E_*^{-1}\mathscr{B}_*), \text{ that is to say:}
                        \dim(\mathscr{B}_*/(\mathscr{B}_* \cap \operatorname{Im} E_*)) \geq \dim(\mathscr{K}_{E_*}), which is the second condition<sup>24</sup> of Theorem
                        6.
 846
                        iv) If \mathfrak{R}^{imp}(E_*, A_*, B_*, C_*) satisfies Theorem 6, then (5.21) and (5.23) are sat-
                        is fied. From the first condition of Theorem 6, we have: \mathscr{X}_d/\mathscr{V}^* = \mathscr{R}^*_{\mathscr{X}_d/\mathscr{V}^*} = \mathscr{R}^*_{\mathscr{X}_d/\mathscr{V}^*}
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 $[\]begin{array}{ll} \overline{^{23}} \text{ Note that: } \mathscr{X}_d = \mathscr{R}^*_{\mathscr{X}_d} + \mathscr{V}^* \subset \mathscr{V}^*_{\mathscr{X}_d} \subset \mathscr{X}_d, \text{ and recall } (5.20). \end{array}$ $^{24} \text{ Note that: } \mathscr{X}_d/\mathscr{V}^* = \mathscr{R}^*_{\mathscr{X}_d/\mathscr{V}^*} \subset \mathscr{V}^*_{\mathscr{X}_d/\mathscr{V}^*} \subset \mathscr{X}_d/\mathscr{V}^*.$

849 $(\Phi \mathscr{V}_{\mathscr{X}_d}^*) \cap (\Phi \mathscr{S}_{\mathscr{X}_d}^*)$, which implies: $\mathscr{X}_d = \mathscr{V}_{\mathscr{X}_d}^* \cap (\mathscr{S}_{\mathscr{X}_d}^* + \mathscr{V}^*) = \mathscr{V}_{\mathscr{X}_d}^* \cap \mathscr{S}_{\mathscr{X}_d}^* + \mathscr{V}^*$ 850 $= \mathscr{R}_{\mathscr{X}_d}^* + \mathscr{V}^*$, which is the first condition of Theorem 7. From the sec851 ond condition of Theorem 6, we have: $^{24} \dim(\Pi \mathscr{B}) = \dim(\mathscr{B}_*) \geq \dim(\mathscr{K}_{E_*}) +$ 852 $\dim(\mathscr{B}_* \cap \operatorname{Im} E_*) = \dim(E_*^{-1}\mathscr{B}_*) = \dim(E_*^{-1}\Pi\mathscr{B}) = \dim(\Phi E_{F_d}^{-1}\mathscr{B})$. Then (recall

853 (E1)): $\dim(\mathscr{B}) \geq \dim(\Phi E_{F_d}^{-1}\mathscr{B}) + \dim(\mathscr{B} \cap \operatorname{Ker} \Pi) = \dim(\Phi E^{-1}\mathscr{B}) + \dim(\mathscr{B} \cap \operatorname{Im} E)$ 854 $E_{F_d}^* \mathscr{V}^*) = \dim(E^{-1}\mathscr{B}) - \dim(\mathscr{V}^* \cap E^{-1}\mathscr{B}) = \dim(\mathscr{K}_E) + \dim(\mathscr{B} \cap \operatorname{Im} E) - \dim(\mathscr{V}^* \cap E^{-1}\mathscr{B})$, which is the second condition 23 of Theorem 7.

Appendix F. Geometric Inequalities (5.26) and (5.27)

From (ALG–V), (ALG–S) and (5.25), we obtain: $\mathcal{V}_{\mathcal{X}_d}^* = \overline{\mathcal{X}} \oplus \mathcal{Q}$ and $\mathcal{S}_{\mathcal{X}_d}^*$ $= \langle \overline{A} \mid \operatorname{Im} \left[\overline{B} \, \overline{S} \right] \rangle \oplus \mathcal{Q}$, which imply: $E\mathcal{V}_{\mathcal{X}_d}^* = \overline{\mathcal{X}}$ and $\mathcal{R}_{\mathcal{X}_d}^* = \langle \overline{A} \mid \operatorname{Im} \left[\overline{B} \, \overline{S} \right] \rangle \oplus \mathcal{Q}$.

From (ALG–V) and (5.25), we get: $\mathcal{V}^0 = E^{-1}\operatorname{Im} E = E^{-1}\overline{\mathcal{V}}_{\left[\overline{B}\,\overline{S}\right]}^0$ and \mathcal{V}^1 $= E^{-1}\mathcal{K}_{\overline{C}} = E^{-1}\overline{\mathcal{V}}_{\left[\overline{B}\,\overline{S}\right]}^1$, then: $E\mathcal{V}^0 = \overline{\mathcal{V}}_{\left[\overline{B}\,\overline{S}\right]}^0$ and $E\mathcal{V}^1 = \overline{\mathcal{V}}_{\left[\overline{B}\,\overline{S}\right]}^1$. Let us assume that: $E\mathcal{V}^\mu = \overline{\mathcal{V}}_{\left[\overline{B}\,\overline{S}\right]}^\mu$, then: $\mathcal{V}^{\mu+1} = (E^{-1}\mathcal{K}_{\overline{C}}) \cap \left[\overline{A} \, \overline{S} \, \right]^{-1} \left(\overline{\mathcal{V}}_{\left[\overline{B}\,\overline{S}\right]}^\mu + \overline{\mathcal{B}} \right)$, which implies: $E\mathcal{V}^{\mu+1} = \mathcal{K}_{\overline{C}} \cap E \left[\overline{A} \, \overline{S} \, \right]^{-1} \left(\overline{\mathcal{V}}_{\left[\overline{B}\,\overline{S}\right]}^\mu + \overline{\mathcal{B}} \right) = \mathcal{K}_{\overline{C}} \cap \overline{A}^{-1} \left(\overline{\mathcal{V}}_{\left[\overline{B}\,\overline{S}\right]}^\mu + \operatorname{Im} \left[\overline{B} \, \overline{S} \, \right] \right) = \overline{\mathcal{V}}_{\left[\overline{B}\,\overline{S}\right]}^\mu$. Thus: $E\mathcal{V}^* = \overline{\mathcal{V}}_{\left[\overline{B}\,\overline{S}\right]}^*$.

From the previous paragraphs we have the following equivalences: $\mathcal{R}_{\mathcal{X}_d}^*$.

From the previous paragraphs we have the following equivalences: $\mathcal{R}_{\mathcal{X}_d}^*$.

Fighthalpoone is $\mathcal{R}_{\mathcal{X}_d}^* = \overline{\mathcal{X}} \oplus \mathcal{Q} \oplus \overline{\mathcal{X}} \oplus \overline{\mathcal{X}}$

From the two first paragraphs, (5.22) takes the form:

$$\dim\left(\frac{\overline{\mathscr{X}}}{\overline{\mathscr{V}}_{[\overline{B}\ \overline{S}]}^* + \overline{\mathscr{B}}}\right) + \dim(\overline{\mathscr{B}}) \ge \dim\left(\frac{\overline{\mathscr{X}}}{\overline{\mathscr{V}}_{[\overline{B}\ \overline{S}]}}\right) + \dim\left(\frac{\{0\} \oplus \mathscr{Q}}{\overline{\mathscr{V}}^* \cap \mathscr{K}_E}\right) \tag{F1}$$

From (ALG–V), (5.25) and the second paragraph, we obtain:

$$\mathcal{K}_{E} \cap \mathcal{V}^{*} = \mathcal{K}_{E} \cap (E^{-1}\mathcal{K}_{\overline{C}}) \cap \left[\overline{A} \ \overline{S} \right]^{-1} \left(\overline{\mathcal{V}}_{\left[\overline{B} \ \overline{S}\right]}^{*} + \overline{\mathcal{B}} \right)
= \mathcal{K}_{E} \cap \left[\overline{A} \ \overline{S} \right]^{-1} \left(\overline{\mathcal{V}}_{\left[\overline{B} \ \overline{S}\right]}^{*} + \overline{\mathcal{B}} \right) = \{0\} \oplus \overline{S}^{-1} \left(\overline{\mathcal{V}}_{\left[\overline{B} \ \overline{S}\right]}^{*} + \overline{\mathcal{B}} \right)
(F2)$$

From (F1) and (F2) we get (5.27) (recall that $\operatorname{Ker} \overline{S} = \{0\}$).