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Observability and Detectability of Singular Linear Systems with Unknown Inputs[☆]

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Abstract

In this paper the strong observability and strong detectability of a general class of singular linear systems with unknown inputs are tackled. The case when the matrix pencil is non-regular is comprised (i.e., more than one solution for the differential equation is allowed). It is shown that, under suitable assumptions, the original problem can be studied by means of a regular (non-singular) linear system with unknown inputs and algebraic constraints. Thus, it is shown that for purposes of analysis, the algebraic equations can be included as part of an extended system output. Based on this analysis, we obtain necessary and sufficient conditions guaranteeing the observability (or detectability) of the system in terms of the zeros of the system matrix. Corresponding algebraic conditions are given in order to test the observability and detectability. A formula is provided that expresses the state as high order derivative of a function of the output, which allows for the reconstruction of the actual state vector. It is shown that the unknown inputs may be reconstructed also.

Keywords: Singular systems, strong detectability, strong observability, algebraic observability.

1. Introduction

The problem of designing an observer for a multi-variable linear system partially driven by unknown inputs (UI) has been widely studied (Darouach et al. (1994); Floquet and Barbot (2004); Floquet et al. (2007); Guan and Saif (1992)). Such observers can be of important use for systems subject to disturbances or with inaccessible inputs, or when dealing with the problem of fault diagnosis. Observability and the problem of observer design have been widely studied for singular systems with perfectly known model. The problems of solvability, controllability and observability were studied in Yip and Sincovec (1981). There, the observability analysis is addressed and algebraic characterizations were found. By using the distributional framework, in Cobb (1984) the algebraic duality between controllability and observability is proven. Necessary and sufficient conditions allowing for the design of a Luenberger-like observer were found in Paraskevopoulos and Koumboulis (1992). In the three previous mentioned works, it was considered that the system has a regular matrix pencil which entails a unique state solution. Without any particular assumption over the matrix pencil of the system, the casual observability, which does not allow to use neither the derivatives of the input nor the derivatives of the output, is studied in Hou and Müller (1999a). The same authors suggest an observer design in Hou and Müller (1999b). In that work it is shown that by allowing the derivatives of the input and output to be involved in the

observer (called it there as a generalized observer), detectability is enough for the convergence of the observation error. A reduced order observer is designed in Darouach and Boutayeb (1995). In Darouach and Boutat-Baddas (2008) an observer for nonlinear singular systems is proposed. In spite of the extended literature regarding the observability analysis and synthesis of singular systems, there exist few results dealing with such problems when the system contains UI. In Paraskevopoulos et al. (1992), the observer design problem is considered for singular linear systems with UI (SLSUI) and necessary and sufficient conditions are given for the design of a Luenberger-like observer. In Darouach et al. (1996), a reduced order observer is proposed. Under some regularity conditions, the observer design is studied in Chu and Mehrmann (1999). Meanwhile, in Koenig (2005) a proportional multiple-integral observer is proposed. Using the graph-theory approach, observability conditions are found in Boukhobza and Hamelin (2007).

In this note, the observability problem of a general class of singular linear systems with UI is studied. The system is not required to have a regular matrix pencil. We obtain necessary and sufficient conditions for the strong observability and strong detectability. We show that the reconstruction of the state can be carried out by a formula that expresses the actual state as a high order derivative of a function of the output. The manuscript also includes Section 2, where the system is described and the strong observability and strong detectability are defined. In Section 3, we obtain the necessary and sufficient conditions which allow for the reconstruction of the state vector. Section 4 deals with an explicit formula that allows for the reconstruction of the state vector. The finite time reconstruction (observability) is considered in 4.1. The procedure for the asymptotic reconstruction (detectability) is given in 4.2. A summarized algo-

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rithm explaining all the estimation procedure is presented in 4.3. To reinforce the theoretical results, we present an example with simulations in Section 5. The following notation will be used throughout the paper. For a matrix X , we denote by X^\perp a full row rank matrix such that $X^\perp X = 0$, and by $X^{\perp\perp}$ a full row rank matrix such that $\text{rank } X^{\perp\perp} X = \text{rank } X$ (then the matrix $\begin{bmatrix} (X^\perp)^T & (X^{\perp\perp})^T \end{bmatrix}^T$ is nonsingular). The Moore-Penrose pseudo-inverse matrix of X is denoted by X^+ . The rank of X is denoted by means of ρ_X . By $\|\cdot\|$, we mean the Euclidean norm. \mathbb{C}^- denotes the set of complex numbers with strictly negative real part. I_r is the identity matrix of dimension r by r . $0_{r \times s}$ is the zero matrix of dimension r by s . For the limit from above, $x(0^+) = \lim_{t \rightarrow 0^+} x(t)$.

2. System Description and Problem Formulation

Let us consider the SLSUI governed by the following equations

$$\Sigma : \begin{cases} E\dot{x}(t) &= Ax(t) + D\mu(t) \\ y(t) &= Cx + F\mu(t) \end{cases}, \quad (1)$$

where $x(t) \in \mathbb{R}^n$ is the state vector, $y(t) \in \mathbb{R}^p$ is the system output, and $\mu(t) \in \mathbb{R}^m$ is the unknown input vector. Matrices $E, A \in \mathbb{R}^{n \times n}$, $D \in \mathbb{R}^{n \times m}$, $C \in \mathbb{R}^{p \times n}$, and $F \in \mathbb{R}^{p \times m}$ are all constant. The matrix E is assumed to be **singular**. Given a state $x_0 \in \mathbb{R}^n$ and a function $\mu(t)$, we denote by $x_\mu(x_0, t)$ the state of Σ at time t which results from taking the initial condition equal to x_0 and the input vector is equal to μ . Therefrom, in a straightforward manner we define the output $y_\mu(x_0, t) = Cx_\mu(x_0, t) + F\mu(t)$.

We are interested in the reconstruction of the (non-impulsive) trajectory of state vector $x(t)$ given the output information $y(\tau)_{\tau \in [0, t]}$. System Σ is *not assumed* to have a regular pencil (Kaczorek (2007)), i.e., it is allowed that $\det(sE - A) = 0$ for all $s \in \mathbb{C}$ (then $x_\mu(x_0, t)$ may have more than a solution). Nevertheless, $\mu(t)$ must be so that $x(t)$ be *piecewise continuous* for all $t > 0$; however, an impulse may occur at $t = 0$. In order to give algebraic conditions allowing the reconstruction of $x(t)$, we consider the following definitions, which are based on classical definitions for linear time invariant systems (see, e.g. Trentelman et al. (2001)).

Definition 1 (Strong observability). *The system Σ is strongly observable (SO) if for all $x_0 \in \mathbb{R}^n$ and for every input function μ , the following implication is satisfied*

$$y_\mu(x_0, t) = 0 \quad \forall t > 0 \text{ implies } x(0^+) = 0. \quad (2)$$

Definition 2 (Strong detectability). *The system Σ is strongly detectable (SD) if for all $x_0 \in \mathbb{R}^n$ and for every input function μ , the following implication holds*

$$y_\mu(x_0, t) = 0 \quad \forall t > 0 \text{ implies } \lim_{t \rightarrow \infty} x_\mu(x_0, t) = 0. \quad (3)$$

It is clear that strong observability is a necessary condition to reconstruct the entire trajectory of the state $x(t)$. Indeed, let us suppose that Σ is not strongly observable, then it means that there exist $\bar{\mu}$ and \bar{x}_0 such that $y_{\bar{\mu}}(\bar{x}_0, t) = 0 \quad \forall t > 0$, but

$x(0^+) \neq 0$. Then, since we assume that $x(t)$ is piecewise continuous, then $x_{\mu=0}(0, t) \equiv 0$ and $x_{\bar{\mu}}(\bar{x}_0, t) \neq 0$ in an open interval, however, both yield a system output identically equal to zero. Thereby, it would be impossible to reconstruct the entire state trajectory. Below it would be proven that SO is a structural necessary and sufficient condition for the reconstruction in finite time of $x(t)$. Analogously, it will be shown that SD is a structural necessary and sufficient condition for the asymptotic reconstruction of $x(t)$.

3. Observability Analysis

Since E is singular, there exist non-singular matrices $T \in \mathbb{R}^{n \times n}$ and $S \in \mathbb{R}^{n \times n}$ such that E can be transformed as follows¹

$$TES = \begin{bmatrix} I_{\rho_E} & 0 \\ 0 & 0 \end{bmatrix} \quad (4)$$

Thus, let us define the vector $z := \begin{bmatrix} z_1^T & z_2^T \end{bmatrix}^T = S^{-1}x$, where $z_1 \in \mathbb{R}^{\rho_E}$ and $z_2 \in \mathbb{R}^{n-\rho_E}$. In these new coordinates, Σ can be rewritten as follows

$$\Psi : \begin{cases} TES\dot{z}(t) &= TASz(t) + TD\mu(t) \\ y(t) &= CSz(t) + F\mu(t) \end{cases}. \quad (5)$$

In view of (4), Ψ takes the following form

$$\dot{z}_1(t) = T_1AS_1z_1(t) + T_1AS_2z_2(t) + T_1D\mu(t), \quad (6a)$$

$$0 = T_2AS_1z_1(t) + T_2AS_2z_2(t) + T_2D\mu(t), \quad (6b)$$

$$y = CS_1z_1(t) + CS_2z_2(t) + F\mu(t). \quad (6c)$$

where S_1 and S_2 matrices arise from the following partition of S , $S = \begin{bmatrix} S_1 & S_2 \end{bmatrix}$ with $S_1 \in \mathbb{R}^{\rho_E \times \rho_E}$ and $S_2 \in \mathbb{R}^{\rho_E \times (n-\rho_E)}$. Analogously, T_1 and T_2 matrices come from the partition $T^T = \begin{bmatrix} T_1^T & T_2^T \end{bmatrix}$ with $T_1 \in \mathbb{R}^{\rho_E \times n}$ and $T_2 \in \mathbb{R}^{(n-\rho_E) \times n}$. It is clear that Σ is SO (resp. SD)² if, and only if, Ψ is SO (resp. SD). Below we will see that a simple manner to study the observability of Ψ , and by extension of Σ , is by considering (6b) as part of the system output of a new pseudo system and considering z_2 as part of the vector of UI. Indeed, let us define the system Φ by means of the following equation,

$$\Phi : \begin{cases} \dot{z}_1(t) &= \bar{A}z_1(t) + \bar{D}v(t) \\ \bar{y}(t) &= \bar{C}z_1(t) + \bar{F}v(t) \end{cases}, \quad (7)$$

where $v(t) \in \mathbb{R}^{n-\rho_E+m}$, $\bar{y}(t) \in \mathbb{R}^{n-\rho_E+p}$ and the matrices \bar{A} , \bar{D} , \bar{C} , and \bar{F} are defined as follows

$$\bar{A} = T_1AS_1, \quad \bar{D} = \begin{bmatrix} T_1AS_2 & T_1D \end{bmatrix}, \\ \bar{C} = \begin{bmatrix} T_2AS_1 \\ CS_1 \end{bmatrix}, \quad \bar{F} = \begin{bmatrix} T_2AS_2 & T_2D \\ CS_2 & F \end{bmatrix}. \quad (8)$$

¹We might select $S = \begin{bmatrix} S_1 & S_2 \end{bmatrix}$ to be non-singular and so that $\text{im } S_2 = \ker E$. Thus, $ES = \begin{bmatrix} ES_1 & 0 \end{bmatrix}$ and $\text{rank } ES_1 = \text{rank } E$. Then a non-singular matrix T might be selected as $T^T = \begin{bmatrix} T_1^T & T_2^T \end{bmatrix}$ so that $T_1ES_1 = I$ and $T_2ES_1 = 0$, one possibility is to select $T_1 = (ES_1)^+ = \left[(ES_1)^T (ES_1) \right]^{-1} (ES_1)^T$.

²Respectively strongly detectable.

It is clear by (6) that Φ looks like system Ψ . In general, they do not represent identical systems. However, both systems are identical if these both identities hold: $\bar{y}^T = \begin{bmatrix} 0^T & y^T \end{bmatrix}$ and $v^T(t) = \begin{bmatrix} z_2^T(t) & \mu^T(t) \end{bmatrix}$. In the next theorem it is claimed that the fulfillment of the SO (resp. SD) of Σ is equivalent to the fulfillment of the SO (resp. SD) of Φ (condition needed for the reconstruction of z_1) plus a rank condition (required for the reconstruction of z_2).

Theorem 1. *System Σ is SO (resp. SD) if, and only if, Φ is SO (resp. SD) and the following rank condition holds*

$$\rho_B = n - \rho_E + \text{rank} \begin{bmatrix} D \\ F \end{bmatrix}, \text{ where } \bar{B} := \begin{bmatrix} \bar{D} \\ \bar{F} \end{bmatrix}. \quad (9)$$

Furthermore, this equivalence is independent of the choice of T and S .

Proof. Firstly, notice that, since $\text{rank} \begin{bmatrix} D^T & F^T \end{bmatrix} = \text{rank} \begin{bmatrix} (TD)^T & F^T \end{bmatrix}$, the fulfillment of (9) is equivalent to say that $\bar{B} \begin{bmatrix} z_2^T & \mu^T \end{bmatrix}^T = 0$ only if $z_2 = 0$.

Necessity. Let Σ be SO (resp. SD). Let us suppose that, for an input v and a state z_1 , $\bar{y}_v(z_1, t) = 0$. For the case when $z_2(t)$ and $\mu(t)$ satisfy the identity $\begin{bmatrix} z_2^T(t) & \mu^T(t) \end{bmatrix} = v^T(t)$, we obtain that Ψ and Φ represent the same system. Thus, defining $x_0 = Sz(0)$ we conclude that $y_\mu(x_0, t) = 0$ for all $t > 0$. Now, since, by assumption, (2) (resp. (3)) holds, $x(0^+) = 0$ (resp. $x(t)$ converges to zero), which in turn implies that $z(0^+) = 0$ (resp. $z(t)$ converges to zero), in particular $z_1(0^+) = 0$ (resp. $\lim_{t \rightarrow \infty} z_1(t) = 0$), i.e. Φ is SO (resp. SD).

Now, assume that (9) does not hold. Then, there exists a vector v which can be divided as $v^T = \begin{bmatrix} v_1^T & v_2^T \end{bmatrix}$ ($v_1 \in \mathbb{R}^{n-\rho_E}$, $v_2 \in \mathbb{R}^m$) so that $\bar{B}v = 0$ and $v_1 \neq 0$. By selecting $z_2(t) = v_1$ and $\mu(t) = v_2$, and $z_1(0^+) = 0$, eq. (6) is fulfilled and $y(t) = 0$ for all $t > 0$. Therefore $x(t) = Sz(t) = S \begin{bmatrix} 0 & v_1^T \end{bmatrix}^T = \text{const} \neq 0$. That is, in such a case Σ is not SO (resp. Σ is not SD).

Sufficiency. Firstly, let us notice that $y_\mu(x_0, t) \equiv 0$ implies $\bar{y}_v(z_{10}, t) \equiv 0$ for some z_{10} and v . Indeed, let us suppose that $y_\mu(x_0, t) = 0$ for a state $x_0 \in \mathbb{R}^n$ and an input function μ . Then, the function $z_\mu(z_0, t) = S^{-1}x_\mu(x_0, t)$ satisfies (6), specifically the algebraic constraint in (6b) is fulfilled. Therefore, we conclude that, for system Φ , $\bar{y}_v(z_{10}, t) = 0$, for all $t > 0$, with v as the extended vector of z_2 and μ . Furthermore, it is known, from linear system theory, that an input $v(t)$ zeroing the output must have the form $v(t) = K^*z_1(t) + Lw(t)$ ($t > 0$), for a particular matrix K^* , a matrix L such that $\bar{B}L = 0$, and a function $w(t)$.

Thus, assuming that Φ is SO (resp. SD), $y_\mu(x_0, t) = 0$ for all $t > 0$, implies, since also $\bar{y}_v(z_{10}, t) = 0$, for all $t > 0$, that $z_1(t) \equiv 0$ (resp. $z_1(t)$ converges to zero) and $\bar{B}v = 0$ (resp. $\bar{B}v = K^*z_1(t)$). The previous identity, assuming that (9) is true, implies that $z_2(t) \equiv 0$ (resp. $z_2(t)$ converges to zero). Therefore, in such a case, $z(t) \equiv 0$ and so $x(t) \equiv 0$. Therefore, we conclude that Σ is SO (resp. SD).

The independence from T and S is straightforward since we have proven Theorem 1 for arbitrary T and S satisfying (4). ■

As for Σ , we could expect that SO and SD can be completely characterized by the five-tuple (E, A, C, D, F) . Indeed, let $R(s)$ be the so-called system matrix of Σ , i.e.,

$$R(s) = \begin{bmatrix} sE - A & -D \\ C & F \end{bmatrix}, s \in \mathbb{C}.$$

We say that $s_0 \in \mathbb{C}$ is a zero of Σ if $\rho_{R(s_0)} < n + \text{rank} \begin{bmatrix} D \\ F \end{bmatrix}$. Let $\sigma_z(\Sigma)$ be defined as the set of zeros of Σ . Let us characterize SO and SD in terms of the zeros of Σ .

Corollary 2. *System Σ is SO (resp. SD) if, and only if, $\sigma_z(\Sigma) = \emptyset$ (resp. $\sigma_z(\Sigma) \subset \mathbb{C}^-$).*

Proof. Let us define $Q(s) = \begin{bmatrix} sI - \bar{A} & -\bar{D} \\ \bar{C} & \bar{F} \end{bmatrix}$, which is the system matrix of Φ . In view of (4) and (8), we obtain that,

$$\begin{bmatrix} T & 0 \\ 0 & I \end{bmatrix} R(s) \begin{bmatrix} S & 0 \\ 0 & I \end{bmatrix} = Q(s) \quad (10)$$

Thus, we deduce that $\rho_{R(s)} = \rho_{Q(s)}$. Then, Corollary 2 follows from the claims of Theorem 1 and the fact that Φ is SO (resp. SD) if, and only if, $\sigma_z(\Phi) = \emptyset$ (resp. $\sigma_z(\Phi) \subset \mathbb{C}^-$). ■

4. Algebraic Observability

As we might expect SO coincides with algebraic observability (AO): we say that Σ is AO if x can be expressed as an algebraic function of y and a finite number of its derivatives (see, e.g. Diop and Fliess (1991)). Let M_k ($k \geq 1$) be the matrices obtained by the following recursive algorithm (see, Molinari (1976)),

$$\begin{aligned} M_{k+1} &= N_{k+1}^{\perp\perp} N_{k+1}, \quad M_1 = (\bar{F}^{\perp\perp} \bar{C})^{\perp\perp} \bar{F}^{\perp\perp} \bar{C}, \\ N_{k+1} &= T_k \begin{pmatrix} M_k \bar{A} \\ \bar{C} \end{pmatrix}, \quad T_k = \begin{pmatrix} M_k \bar{D} \\ \bar{F} \end{pmatrix}^{\perp}. \end{aligned} \quad (11)$$

Let us denote by l , **the smallest integer such that $\rho_{M_l} = \rho_{M_{l+1}}$** . For our purposes, we point out that **Φ is SO if, and only if, $\rho_{M_l} = \rho_E$** . For the case of SD we have to work a bit more with system Φ . Indeed, let us assume that $\rho_{M_l} < \rho_E$. Let V be a full column rank matrix so that $M_l V = 0$. There exists a pair of matrices Q and K^* such that

$$\bar{A}V + \bar{D}K^* = VQ \quad \text{and} \quad \bar{C}V + \bar{F}K^* = 0. \quad (12)$$

From (12), it is clear that the $\text{im}((\bar{A} + \bar{D}K^*V^+)V) \subset \text{im} V$ and $(\bar{C} + \bar{F}K^*V^+)V = 0$. We can define a non-singular matrix P of dimension ρ_E as $P = \begin{bmatrix} M_l \\ V^+ \end{bmatrix}$, $P^{-1} = \begin{bmatrix} M_l^+ & V \end{bmatrix}$, where V^+ and M_l^+ are the Moore-Penrose pseudo-inverse of V and M_l , respectively: $V^+ = (V^T V)^{-1} V^T$ and $M_l^+ = M_l^T (M_l M_l^T)^{-1}$. By defining the vectors $w_1 = M_l z_1$ and $w_2 = V^+ z_1$, we have that

³Such a statement was proven in Hautus (1983).

$z_1 = M_l^+ w_1 + V w_2$. System Φ in these new coordinates can be rewritten as follows:

$$\dot{w}_1 = \bar{A}_1 w_1 + \bar{D}_1 (v - K^* w_2), \quad (13a)$$

$$\dot{w}_2 = \bar{A}_2 w_1 + \bar{A}_3 w_2 + \bar{D}_2 (v - K^* w_2), \quad (13b)$$

$$\bar{y} = \bar{C}_1 w_1 + \bar{F} (v - K^* w_2), \quad (13c)$$

where

$$\begin{aligned} \bar{A}_1 &= M_l (\bar{A} + \bar{D} K^* V^+) M_l^+, \quad \bar{D}_1 = M_l \bar{D}, \\ \bar{A}_2 &= V^+ (\bar{A} + \bar{D} K^* V^+) M_l^+, \quad \bar{D}_2 = V^+ \bar{D}, \\ \bar{A}_3 &= V^+ (\bar{A} + \bar{D} K^* V^+) V, \quad \bar{C}_1 = \bar{C} M_l^+. \end{aligned} \quad (14)$$

Thus, as for SD, it is known that system Φ is **SD if, and only if**, $\text{rank} \begin{bmatrix} \bar{D}_1 \\ \bar{F} \end{bmatrix} = \text{rank} \begin{bmatrix} \bar{D} \\ \bar{F} \end{bmatrix}$ and \bar{A}_3 is a Hurwitz matrix (see, e.g. Bejarano et al. (2009)).

Coming back to system Φ . Define $\xi_1 := (\bar{F}^\perp \bar{C})^{\perp\perp} \bar{F}^\perp \bar{y} = M_1 z_1$, with M_1 defined as in (11). Let us derive the vector ξ_1 :

$$\dot{\xi}_1(t) = M_1 \bar{A} z_1(t) + M_1 \bar{D} v(t). \quad (15)$$

Let us define a new vector ξ_2 as follows

$$\xi_2 := N_2^{\perp\perp} T_1 \begin{bmatrix} \xi_1 \\ \bar{y}(t) \end{bmatrix}, \quad (16)$$

with $N_2^{\perp\perp}$ and T_1 defined by (11). Thus, taking into account (7), (15), and (11), we have that

$$\frac{d}{dt} J_2 \begin{bmatrix} \bar{y} \\ \int_{t_0}^t \bar{y}(\tau) d\tau \end{bmatrix} = \xi_2 = M_2 z_1(t), \quad t > t_0 \geq 0, \quad (17)$$

where

$$J_2 = N_2^{\perp\perp} T_1 \begin{bmatrix} J_1 & 0 \\ 0 & I_{\bar{p}} \end{bmatrix}, \quad J_1 = (\bar{F}^\perp \bar{C})^{\perp\perp} \bar{F}^\perp.$$

In the first identity of (17), we take outside the differential operator from (16) and use the definition of ξ_1 . Thus, we can follow an iterative procedure to obtain the following set of equations, for $k \geq 1$,

$$\frac{d^k}{dt^k} J_{k+1} \begin{bmatrix} \bar{y}(t) \\ \vdots \\ \int_{t_0}^t \cdots \int_{t_0}^{\tau_k} \bar{y}(\tau_1) d\tau_1 \cdots d\tau_k \end{bmatrix} = M_{k+1} z_1, \quad (18)$$

where M_{k+1} is defined by (11), and J_{k+1} is defined by the following recursive algorithm, for $k \geq 1$,

$$J_1 = (\bar{F}^\perp \bar{C})^{\perp\perp} \bar{F}^\perp, \quad J_{k+1} = N_{k+1}^{\perp\perp} T_k \begin{bmatrix} J_k & 0 \\ 0 & I_{\bar{p}} \end{bmatrix}. \quad (19)$$

Thus $M_k z_1$ is expressed by a high order derivative of a function of $y(t)$. In such a way a real-time differentiator could be used, two of them frequently used due to their finite time convergence can be found in Levant (2003) and Mboup et al. (2009a). For instance if $\text{rank} M_l = \rho_E$, then z_1 is algebraically observable, i.e. it could be reconstructed by using a real-time differentiator.

In order to match system Σ with system Φ , from now on, we define $\bar{y} = \begin{bmatrix} 0_{1 \times n - \rho_E} & y^T \end{bmatrix} \in \mathbb{R}^{\bar{p}}$ ($\bar{p} := n - \rho_E + p$), and $v(t) = \begin{bmatrix} z_2^T(t) & \mu^T(t) \end{bmatrix}^T \in \mathbb{R}^q$ ($q = n - \rho_E + m$), then in view of (6), equations (5) and (7) are identical. Below, we consider two cases: when Σ is SO and when it is SD, but not SO. Of course, since Φ is a standard linear system, there might be other methods, besides the one proposed below, that might be used to carry out the algebraic reconstruction of the state.

4.1. Finite time reconstruction

Let us consider that system Σ is SO. Then, the reconstruction of entire state vector $x(t)$ can be achieved in finite time: by means of an algebraic formula. Let us proceed in the following way. Since Φ is SO, $\rho_{M_l} = \rho_E$. Then in this case, from (18), we obtain the equation

$$\frac{d^{l-1}}{dt^{l-1}} M_l^{-1} J_l \begin{bmatrix} \bar{y}(t) \\ \vdots \\ \int_{t_0}^t \cdots \int_{t_0}^{\tau_{l-1}} \bar{y}(\tau_1) d\tau_1 \cdots d\tau_{l-1} \end{bmatrix} = z_1, \quad (20)$$

where $M_l \in \mathbb{R}^{\rho_E \times \rho_E}$ and $J_l \in \mathbb{R}^{\rho_E \times \bar{p}^l}$. Let U be a matrix so that

$$\text{rank} \begin{bmatrix} D \\ F \end{bmatrix} U = \text{rank} \begin{bmatrix} D \\ F \end{bmatrix} =: \bar{m}, \quad U \in \mathbb{R}^{q \times \bar{m}} \quad (21)$$

Since (9) must be satisfied according to Theorem 1, we have that

$$z_2(t) = \begin{bmatrix} I_{n-\rho_E} & 0_{\bar{q}} \end{bmatrix} \left(\begin{bmatrix} \bar{D} \\ \bar{F} \end{bmatrix} \begin{bmatrix} I & 0 \\ 0 & U \end{bmatrix} \right)^+ \left(\begin{bmatrix} \dot{z}_1(t) \\ \bar{y} \end{bmatrix} \right) - \begin{bmatrix} \bar{A} \\ \bar{C} \end{bmatrix} z_1(t) \quad (22)$$

where $\bar{q} := n - \rho_E + \bar{m}$. Now, we are ready to give a formula to reconstruct x in finite time.

Theorem 3. *If system Σ is SO, then the state x can be expressed algebraically by the following formula:*

$$x(t) = \frac{d^l}{dt^l} \begin{bmatrix} S & 0_{n \times \bar{m}} \end{bmatrix} \begin{bmatrix} H_1 \\ H_2 \end{bmatrix} \begin{bmatrix} \bar{y}(t) \\ \vdots \\ \int_{t_0}^t \cdots \int_{t_0}^{\tau_{l-1}} \bar{y}(\tau_1) d\tau_1 \cdots d\tau_{l-1} \end{bmatrix}, \quad (23)$$

where $H_1 \in \mathbb{R}^{\rho_E \times \bar{p}^{(l+1)}}$ and $H_2 \in \mathbb{R}^{\bar{q} \times \bar{p}^{(l+1)}}$ are defined as

$$\begin{aligned} H_1 &:= \begin{bmatrix} 0_{\rho_E \times \bar{p}} & M_l^{-1} J_l \end{bmatrix}, \quad H_2 := \begin{bmatrix} \bar{D} U \\ \bar{F} U \end{bmatrix}^+ (G_1 - G_2), \\ G_1 &:= \begin{bmatrix} M_l^{-1} J_l & 0_{\rho_E \times \bar{p}} \\ 0_{\bar{p} \times \bar{p}^l} & I_{\bar{p} \times \bar{p}} \end{bmatrix}, \quad G_2 := \begin{bmatrix} 0_{\rho_E \times \bar{p}} & \bar{A} M_l^{-1} J_l \\ 0_{\bar{p} \times \bar{p}} & \bar{C} M_l^{-1} J_l \end{bmatrix}, \end{aligned}$$

$G_1, G_2 \in \mathbb{R}^{\rho_E + \bar{p} \times \bar{p}^{(l+1)}}$, and M_l and J defined recursively in (11) and (19), respectively, and U defined by (21).

Proof. Let us define the extended vector $Y_k \in \mathbb{R}^{\bar{p}^{(k+1)}}$ as follows

$$Y_k = \begin{bmatrix} \bar{y}(t) \\ \int_{t_0}^t \bar{y}(\tau_1) d\tau_1 \\ \vdots \\ \int_{t_0}^t \cdots \int_{t_0}^{\tau_k} \bar{y}(\tau_1) d\tau_1 \cdots d\tau_k \end{bmatrix}, \quad k = 1, 2, \dots \quad (24)$$

Then for a matrix X of suitable dimension the following holds:

$$\frac{d^k}{dt^k}XY_k = \frac{d^{k+1}}{dt^{k+1}} \begin{bmatrix} 0 & X \end{bmatrix} Y_{k+1}.$$

Thus, since $x = Sz$, by manipulating (22) and taking into account (20), (23) is obtained. ■

Remark 1. One might obtain a little more from the previous analysis, that is, one can express by an algebraic formula the part of μ that can be reconstructed (assuming Σ is SO). Indeed, let $\bar{\mu}$ be implicitly defined by the equation $\begin{bmatrix} D^T & F^T \end{bmatrix}^T \mu = \begin{bmatrix} D^T & F^T \end{bmatrix}^T U\bar{\mu}$. With the same procedure followed to obtain (23), $\bar{\mu}$ can be expressed by the formula

$$\bar{\mu}(t) = \frac{d^l}{dt^l} \begin{bmatrix} 0_{\bar{m} \times n} & I_{\bar{m}} \end{bmatrix} \begin{bmatrix} \bar{H}_1 \\ \bar{H}_2 \end{bmatrix} \begin{bmatrix} \bar{y}(t) \\ \vdots \\ \int_{t_0}^t \cdots \int_{t_0}^{\tau_1} \bar{y}(\tau_1) d\tau_1 \cdots d\tau_l \end{bmatrix}. \quad (25)$$

4.2. Asymptotic Reconstruction

Let us assume that Σ is SD, but not SO. Next we show how to carry out the estimation of $x(t)$.

Theorem 4. Assuming that Σ is SD, but not SO, we obtain that

$$\lim_{t \rightarrow \infty} \|x(t) - \hat{x}(t)\| = 0,$$

provided that $\hat{x}(t)$ is designed by following (26)-(27).

$$\hat{x}(t) = \frac{d^l}{dt^l} \begin{bmatrix} S & 0_{n \times \bar{m}} \end{bmatrix} \begin{bmatrix} \bar{H}_1 \\ \bar{H}_2 \end{bmatrix} Y_l(t) \begin{bmatrix} S & 0_{n \times m} \end{bmatrix} \begin{bmatrix} V \\ K^* \end{bmatrix} \hat{w}_2 \quad (26)$$

$$\dot{\hat{w}}_2 = \bar{A}_3 \hat{w}_2 + \bar{A}_2 \frac{d^l}{dt^l} \left(\begin{bmatrix} 0_{\rho_E - \rho_M \times \bar{p}} & J_l \end{bmatrix} + \bar{D}_2 U \bar{H}_2 \right) Y_l \quad (27)$$

where $\bar{H}_1 \in \mathbb{R}^{\rho_E \times \bar{p}(l+1)}$, $\bar{H}_2 \in \mathbb{R}^{\bar{q} + \bar{p} \times \bar{p}(l+1)}$, and $\bar{G}_1, \bar{G}_2 \in \mathbb{R}^{\rho_M + \bar{p} \times \bar{p}(l+1)}$ ($\rho_M = \rho_{M_l}$) satisfy the following identities,

$$\bar{H}_1 = \begin{bmatrix} 0_{\rho_E \times \bar{p}} & M_l^+ J_l \end{bmatrix}, \quad \bar{H}_2 := \begin{bmatrix} \bar{D}_1 U \\ \bar{F} U \end{bmatrix}^+ (\bar{G}_1 - \bar{G}_2),$$

$$\bar{G}_1 := \begin{bmatrix} J_l & 0_{\rho_M \times \bar{p}} \\ 0_{\bar{p} \times \bar{p}l} & I_{\bar{p}} \end{bmatrix}, \quad \bar{G}_2 := \begin{bmatrix} 0_{\rho_M \times \bar{p}} & \bar{A}_1 J_l \\ 0_{\bar{p} \times \bar{p}} & \bar{C}_1 J_l \end{bmatrix}.$$

Proof. In this case, as Σ is not SO, by differentiation, we are able to reconstruct $w_1 = M_l z_1$ only, where $\rho_{M_l} < \rho_E$. Since Φ is SD and by (21), we have that $\begin{bmatrix} M_l \bar{D} \\ \bar{F} \end{bmatrix} U$ has full column rank. Then, from (18) and from (13a) and (13c), we have the following expression for w_1 and z_2 ,

$$w_1(t) = \frac{d^{l-1}}{dt^{l-1}} J_l Y_{l-1}(t), \quad (28)$$

$$z_2(t) = \frac{d^l}{dt^l} \begin{bmatrix} I_{n-\rho_E} & 0_{n-\rho_E \times \bar{m}} \end{bmatrix} \bar{H}_2 Y_l(t) + \begin{bmatrix} I_{n-\rho_E} & 0_{n-\rho_E \times m} \end{bmatrix} K^* w_2. \quad (29)$$

Thus, since $x = Sz$ and $z_1 = M_l^+ w_1 + V w_2$, we obtain the identity

$$x(t) = \frac{d^l}{dt^l} \left(\begin{bmatrix} S & 0_{n \times \bar{m}} \end{bmatrix} \begin{bmatrix} \bar{H}_1 \\ \bar{H}_2 \end{bmatrix} Y_l(t) \right) + \begin{bmatrix} S & 0_{n \times m} \end{bmatrix} \begin{bmatrix} V \\ K^* \end{bmatrix} w_2. \quad (30)$$

Therefore, by comparing (26) with (30), we have that

$$\hat{x}(t) = x(t) - \begin{bmatrix} S & 0_{n \times \bar{m}} \end{bmatrix} \begin{bmatrix} V \\ K^* \end{bmatrix} (w(t) - \hat{w}_2(t)). \quad (31)$$

Furthermore, in view of (13b), (28), and (27), we obtain that

$$\dot{w}_2 - \dot{\hat{w}}_2 = \bar{A}_3 (w_2 - \hat{w}_2).$$

Since Φ is SD, \bar{A}_3 is Hurwitz; hence \hat{w}_2 converges to w_2 . Then by (31) the proof is finished. ■

Remark 2. If $\bar{\mu}$ needs to be reconstructed also, then it can be done by means of $\hat{\bar{\mu}}(t)$, defined as follows,

$$\hat{\bar{\mu}}(t) = \frac{d^l}{dt^l} \left(\begin{bmatrix} 0_{\bar{m} \times n} & I_{\bar{m}} \end{bmatrix} \begin{bmatrix} \bar{H}_1 \\ \bar{H}_2 \end{bmatrix} Y_l(t) \right) \begin{bmatrix} 0_{\bar{m} \times n} & I_{\bar{m}} \end{bmatrix} \begin{bmatrix} V \\ \bar{K}^* \end{bmatrix} \hat{w}(t) \quad (32)$$

where \bar{K}^* is implicitly defined by the equation

$$\begin{bmatrix} \bar{D} \\ \bar{F} \end{bmatrix} K^* = \begin{bmatrix} \bar{D} \\ \bar{F} \end{bmatrix} \begin{bmatrix} I_{n-\rho_E} & 0 \\ 0 & U \end{bmatrix} \bar{K}^*$$

So, we obtain straightforwardly that $\|\hat{\mu}(t) - \mu(t)\|$ goes to zero.

4.3. Summarized algorithm for the reconstruction of $x(t)$

Below, there is a step-by-step description that may be followed to carry out the estimation of $x(t)$.

Step 1. Choose S and T to bring E into the form (4).

Step 2. Calculate matrices \bar{A} , \bar{C} , \bar{D} , and \bar{F} according to (8).

Step 3. Calculate matrices M_l and J_l following the recursive algorithms (11) and (19), respectively. Let us remind that l is such that $\rho_{M_{l+1}} = \rho_{M_l} > \rho_{M_{l-1}}$.

Step 4A. Check if Σ is SO: test both conditions i) $\rho_{M_l} = \rho_E$ and ii) condition in (9). If they both are satisfied go to Step 5.A.

Step 4B. Check if Σ is SD: test the conditions i) $\text{rank} \begin{bmatrix} \bar{D}_1 \\ \bar{F} \end{bmatrix} =$

$n - \rho_E + \text{rank} \begin{bmatrix} D \\ F \end{bmatrix}$ and ii) \bar{A}_3 is Hurwitz (\bar{D}_1 and \bar{A}_3 defined in (14)). If they both are satisfied, go to Step 5.B. Otherwise, the exact estimation (even asymptotic) of x is impossible.

Step 5A. $x(t)$ is reconstructed by means of the formula (23).

Step 5B. $x(t)$ is estimated by using $\hat{x}(t)$ described in (26)-(27).

To carry out Step 5 one needs to use a real-time differentiator. In the example presented in the next section, we use algebraic numerical and sliding mode differentiators, which have been proposed in the last years and have a considerable acceptance in several observation and identification procedures thanks to their finite time convergence and robustness. A brief review of those differentiators is given in Appendix A.

5. Example

Let us consider that Σ has the following matrices values

$$E = \begin{bmatrix} 0 & 1 & 0 & 1 \\ 0 & 0 & 0 & 0 \\ 1 & 0 & 1 & 1 \\ 0 & 0 & 0 & 0 \end{bmatrix}, A = \begin{bmatrix} -1 & -1 & 1 & 1 \\ 2 & 1 & 2 & 0 \\ 1 & 1 & 1 & -1 \\ 0 & 0 & 0 & 0 \end{bmatrix}, D = \begin{bmatrix} 1 \\ -1 \\ 0 \\ 0 \end{bmatrix}$$

$$C = \begin{bmatrix} 0 & 1 & 0 & 1 \\ 0 & -1 & 0 & 2 \end{bmatrix}, F = \begin{bmatrix} 0 \\ 1 \end{bmatrix}.$$

It is easy to see that, in this example, $\det(sE - A) = 0$ for every $s \in \mathbb{C}$. Hence, many solutions for $x(t)$ are expected to satisfy the differential equation in (1). However, to each output $y(t)$ corresponds only one trajectory of $x(t)$ (a.e.). Indeed, we will see that, according to Theorem 1, Σ is SO. We describe the observer design following the step by step description.

Step 1. For this case matrices S and T are chosen as follows,

$$S = \begin{bmatrix} 0 & 0 & -1 & -1 \\ 0 & 0 & 0 & -1 \\ 0 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 \end{bmatrix}, T = \begin{bmatrix} 1 & 0 & 0 & 0 \\ -1 & 0 & 1 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 0 & 1 \end{bmatrix},$$

Step 2. Matrices \bar{A} , \bar{C} , \bar{D} , and \bar{F} take the following values:

$$\bar{A} = \begin{bmatrix} 1 & 1 \\ -2 & 0 \end{bmatrix}, \bar{D} = \begin{bmatrix} 2 & 3 & 1 \\ -2 & -6 & -1 \end{bmatrix},$$

$$\bar{C} = \begin{bmatrix} 0 & 2 \\ 0 & 0 \\ 1 & 0 \\ 2 & 0 \end{bmatrix}, \bar{F} = \begin{bmatrix} 0 & -3 & -1 \\ 0 & 0 & 0 \\ 0 & 0 & 0 \\ 0 & 3 & 1 \end{bmatrix}.$$

Step 3. Matrices M_2 and J_2 take the values

$$M_2 = \begin{bmatrix} 1 & 0 \\ \sqrt{2} & \sqrt{2} \end{bmatrix},$$

$$J_2 = \begin{bmatrix} 0 & 0 & 0 & 0 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & \sqrt{2}/2 & 0 & 0 & \sqrt{2}/2 \end{bmatrix}.$$

Step 4. Here, $\text{rank } M_2 = 2$, $\text{rank } \bar{B} = 3$, and $n - \rho_E = 2$. Therefore, both conditions of Theorem 1 are satisfied.

Step 5. The reconstruction of $x(t)$ can be done by means of the formula (23), once we define $\bar{y}^T = [0 \ 0 \ y^T]$. For this example, the reconstruction of μ is also possible following formula (25). Thus, we have

$$x_1 = -\frac{1}{3}y_1 + \frac{7}{12}y_2 - \frac{1}{2}\dot{y}_1 + \frac{1}{6}\dot{y}_2,$$

$$x_2 = \frac{2}{3}y_1 - \frac{1}{6}y_2 + \frac{1}{6}\dot{y}_2,$$

$$x_3 = -\frac{1}{4}y_2 + \frac{1}{2}\dot{y}_1,$$

$$x_4 = \frac{1}{3}y_1 - \frac{1}{6}y_2 + \frac{1}{6}\dot{y}_2,$$

$$\mu = \frac{1}{2}y_2 + \frac{1}{2}\dot{y}_2.$$

To estimate the state x and μ , we use two different differentiators, an algebraic numerical differentiator (ALND) and a high order sliding mode differentiator (HOSMD).

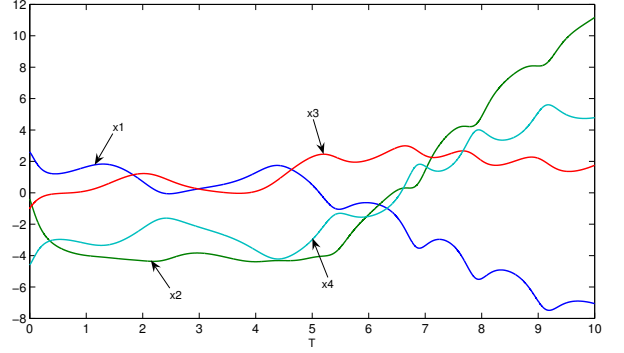


Figure 1: Trajectories of $x(t)$.

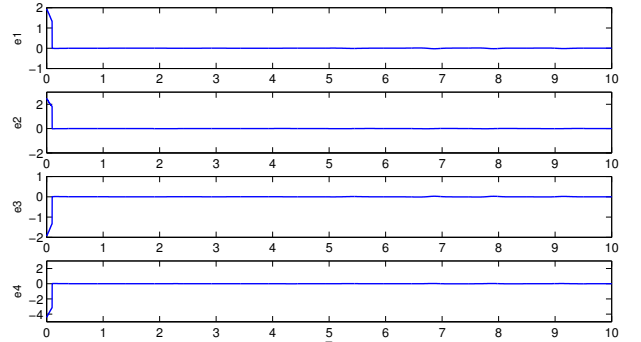


Figure 2: Estimation error $e = x - \hat{x}$ using ALND.

For simulation purposes, we have chosen $\mu = 2 \sin(x_1 - x_2 + x_3) + \cos(t)$ and $x_1 = x_2 - 2 - \cos(3t)$. The state trajectories are depicted in Figure 1. The estimation error is depicted in Figures 2 and 3. The corresponding estimation for μ , denoted by $\hat{\mu}$, is shown in Figure 4. We have also tested both differentiators in presence of output noise, for that we have used an approximation of white noise using the Matlab command `rand` with a uniform distribution on the interval $(-0.1, 0.1)$ for the first output and $(-0.05, 0.05)$ for the second one. The estimation error is depicted in Figure 5.

Conclusions

We have given necessary and sufficient conditions to estimate the slow (non-impulsive) trajectories, for singular systems

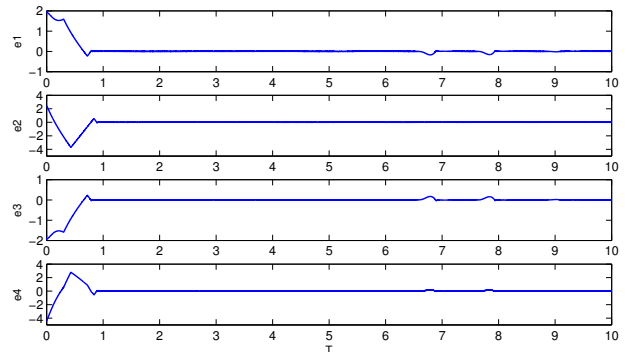


Figure 3: Estimation error $e = x - \hat{x}$ using HOSMD.

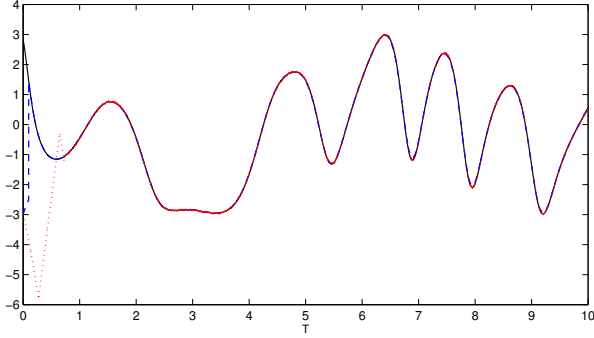


Figure 4: UI μ (solid) and its estimation $\hat{\mu}$ (dashed for ALND, dotted HOSMD).

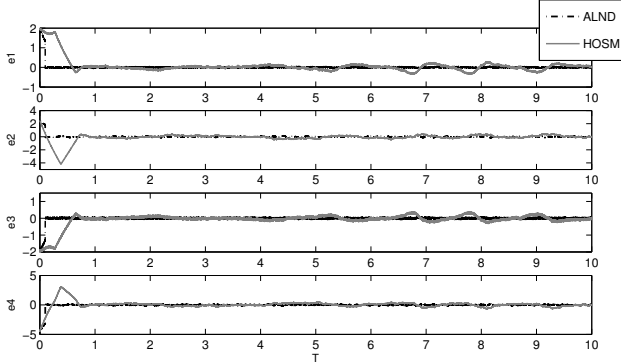


Figure 5: Estimation error $e = x - \hat{x}$ with a noisy output.

in which more than one solution of the differential equation is allowed, i.e. the pencil of the system is not required to be regular as it is assumed in most of the previous works where the observability is studied. We have given explicit formulas to reconstruct in finite time and asymptotically the states. Nevertheless, we have to notice that when the estimate of x needs an *excessive* number of derivatives, the error due to noise and sampling time could increase considerably. In that case, if an asymptotic estimation is enough, we should use the on-line differentiation only the needed times allowing after to design a Luenberger-like observer. In that case, a simple but cumbersome modified procedure might be followed in order to reduce at minimum the number of derivatives required to estimate the state, we refer the interested reader to Floquet et al. (2007) and Bejarano and Fridman (2010).

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Appendix A. Review of two differentiators

Appendix A.1. ALND

This algebraic setting for numerical differentiation of noisy signals is introduced in Fliess et al. (2004) and analyzed in Liu

et al. (2011a,b); Mboup et al. (2007, 2009b) (see also Nöthen (2007) for interesting discussions and comparisons). The reader may find additional theoretical foundations in Fliess (2006); Fliess and Sira-Ramírez (2003). Consider a signal $y(t) = \sum_{l=0}^{\infty} y^{(l)}(0) \frac{t^l}{l!}$ which is assumed to be analytic around $t = 0$ and its truncated Taylor expansion $y_N(t) = \sum_{l=0}^N y^{(l)}(0) \frac{t^l}{l!}$ at order N . The usual rules of symbolic calculus in Schwartz's distribution theory (Schwartz (1966)) yield $y_N^{(N+1)}(t) = y(0)\delta^{(N)} + \dots + y^{(N)}(0)\delta$, where δ is the Dirac measure at zero. Multiply both sides by $(-t)^l$: $(-t)^l y_N^{(N+1)}(t) = (-t)^l (y(0)\delta^{(N)} + \dots + y^{(N)}(0)\delta)$, and apply the rules $t\delta = 0$, $t\delta^{(l)} = -l\delta^{(l-1)}$, $l \geq 1$. We obtain a triangular system of linear equations from which the derivatives $y^{(l)}(0)$ can be obtained ($1 \leq l \leq N$)

$$(-t)^l y_N^{(N+1)}(t) = \frac{N!}{(N-l)!} \delta^{(N-l)} y(0) + \dots + \delta y^{(N-l)}(0). \quad (\text{A.1})$$

It means that the coefficients $y(0), \dots, y^{(N)}(0)$ are *linearly identifiable* Fliess and Sira-Ramírez (2003); Fliess and Sira-Ramírez (2008). The time derivatives of $y_N(t)$, the Dirac measures and its derivatives are removed by integrating with respect to time both sides of Eq. (A.1) at least ν times ($\nu > N$):

$$\int_0^t \int_0^{\tau_{\nu-1}} \dots \int_0^{\tau_1} (-\tau)^l y_N^{(N+1)} dt_{\nu-1} \dots dt_1 d\tau = \frac{N!}{(N-l)!} \frac{t^{\nu-N-l-1}}{(\nu-N-l-1)!} y(0) + \dots + \frac{t^{\nu-1}}{(\nu-1)!} y^{(N-l)}(0).$$

The iterated integrals may be replaced by

$$\int_0^t \int_0^{\tau_{\nu-1}} \dots \int_0^{\tau_1} \tau^\alpha x(\tau) dt_{\nu-1} \dots dt_1 d\tau = \int_0^t \frac{(t-\tau)^{\nu-1}}{(\nu-1)!} \tau^\alpha x(\tau) d\tau. \quad (\text{A.2})$$

It is clear that the numerical estimation rely on $\lim_{N \rightarrow +\infty} [y_N^{(l)}(0)]_{\text{estim}}(t) = y^{(l)}(0)$.

Remark 3. *These iterated integrals are low pass filters which attenuate the noises, which are viewed as highly fluctuating phenomena (see Fliess (2006) for more details). The above formulae may easily be extended to sliding time windows in order to obtain real time estimates (see Mboup et al. (2007, 2009b); Liu et al. (2011a,b) for further details). Moreover, according to the performed algebraic manipulation one can have some different formulae: $\forall t_0 \in I$,*

$$\tilde{x}_{t_0 \pm}^{(n)}(k, \mu, \beta T) = \frac{1}{(\beta T)^n} \int_0^1 \gamma_{k, \mu, n} (1-\tau)^\mu \tau^k P_n^{\mu, k}(\tau) x(t_0 + \beta T \tau) d\tau, \quad (\text{A.3})$$

where x is assumed to be analytic, $P_n^{\mu, k}$ is the n^{th} order Jacobi polynomial defined on $[0, 1]$ with $n \in \mathbb{N}$, $k, \mu \in]-1, \infty[$, $\beta = \pm 1$ and $T > 0$.

Appendix A.2. HOSMD

The HOSMD described in Levant (2003) can be expressed in a dynamic form by means of the equations (A.4), where the

signal to be differentiated is represented by a function $f \in C^{r+1}$.

$$\begin{aligned}
\dot{\alpha}_0 &= -\lambda_r L^{\frac{1}{r+1}} |\alpha_0 - f|^{\frac{r}{r+1}} \text{sign}(\alpha_0 - f) + \alpha_1 \\
\dot{\alpha}_1 &= -\lambda_{r-1} L^{\frac{1}{r}} |\alpha_1 - \dot{\alpha}_0|^{\frac{r-1}{r}} \text{sign}(\alpha_1 - \dot{\alpha}_0) + \alpha_2 \\
&\vdots \\
\dot{\alpha}_{r-1} &= -\lambda_1 L^{\frac{1}{2}} |\alpha_{r-1} - \dot{\alpha}_{r-2}| \text{sign}(\alpha_{r-1} - \dot{\alpha}_{r-2}) + \alpha_r \\
\dot{\alpha}_r &= -\lambda_0 L \text{sign}(\alpha_r - \dot{\alpha}_{r-1})
\end{aligned} \tag{A.4}$$

Then, in the cited paper it was shown that if λ_i 's gains are chosen properly, the differentiator converges in a finite time T , i.e., $\alpha_i(t) = f^{(i)}(t)$, for all $t \geq T$ and $i = 0, 1, \dots, r$. L is a constant such that $\|f^{(r+1)}(t)\| \leq L$. For the case when $r = 5$, the gains could be chosen as $\lambda_0 = 1.1$, $\lambda_1 = 1.5$, $\lambda_2 = 3$, $\lambda_3 = 5$, $\lambda_4 = 8$, $\lambda_5 = 12$.

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