

## Final observability of time-lag systems

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Final observability on  $[0, T]$  for constant time-lag systems is defined and studied. Observability conditions depending on the length of an observation interval are formulated and proved. Two types of observability are considered: one with initial function in the space  $L^p$ , second with continuous initial function.

### 1. Introduction

For constant time-lag systems described by linear differential-difference equations observability problem can be defined in several ways.

In [1] is considered the system given by equation.

$$\begin{aligned}\dot{x}(t) &= A_0 x(t) + A_1 x(t-1), \quad t \geq 0, \quad x(t) \in R^n, \\ y(t) &= Cx(t), \quad y(t) \in R^p, \\ x(t) &= \varphi(t), \quad t \in [-1, 0].\end{aligned}\tag{1.1}$$

The observability problem in classical sense is to determine the initial state  $\{\varphi(0), A_1 \varphi(t), t \in [-1, 0]\}$  where  $\varphi: [-1, 0] \mapsto R^n$  belongs to a preassumed class  $\mathcal{F} \subset L^1([-1, 0], R^n)$ , e.g.  $\varphi(\cdot) \in C([-1, 0], R^n)$  or  $\varphi(\cdot) \in L^p([-1, 0], R^n)$ ,  $1 \leq p \leq \infty$  knowing the output  $y(t)$ ,  $t \in [0, T]$ . In [1] necessary and sufficient conditions are given for observability of (1.1). Similar problem is considered in [9], where for system (1.1) general criteria are proved for determination of the initial function (not the state)  $\varphi(t)$ ,  $t \in [-1, 0]$ , with  $\varphi(\cdot)$  belonging to the set of piecewise continuous functions provided the output  $y(t)$ ,  $t \in [0, T]$ , is given. Spectral observability is discussed in [4], [8]. System (1.1) is spectral observable if and only if corresponding outputs do not vanish for all its eigensolutions.

Measurement and control problems assume quite often the knowledge of the vector  $x(T) \in R^n$  at  $T > 0$ , or the state of the system at  $T > 0$  given as a function  $x(T + \theta)$ ,  $\theta \in [-1, 0]$ , while only output information is available. In general case knowledge

of the state at  $T > 0$  is not an equivalent of information about the initial state, because behaviour of the state trajectory of (1.1) is described by a semigroup (not group) property. Hence it seems to be suitable to introduce the definition of observability as the determination from output  $y(t)$ ,  $t \in [0, T]$ , of vector  $x(T)$  or state  $x(T + \theta)$ ,  $T > 0$ ,  $\theta \in [-1, 0]$ .

Notation: In this paper the following notation will be used:  $A \in R^{m \times p}$  means that  $A$  is a real  $m$  by  $p$  matrix. The kernel and image of a matrix  $A$  will be denoted respectively by  $\text{Ker } A$  and  $\text{Im } A$ .  $A^T$  is the transpose of  $A$ .

The orthogonal complement of a subspace  $\mathcal{K} \subset R^n$  will be denoted by  $\mathcal{K}^\perp$ .

For given matrices  $A \in R^{n \times n}$ ,  $B \in R^{n \times m}$ ,  $\{A|B\} = \text{Im} [B|AB|\dots|A^{n-1}B]$  denotes controllability subspace. If  $\mathcal{B} = \text{Im } B$  the controllability subspace will be also written  $\{A|\mathcal{B}\}$ . Character  $\triangleq$  means equal by definition. Adjoint matrix of  $A$  will be denoted  $\text{adj } A$ . Identity  $p$  by  $p$  matrix is denoted by  $I_p$ .

## 2. Problem Statement and Preliminary Results

Consider the linear, time-invariant system given by the equations

$$\begin{aligned} \dot{x}(t) &= A_0 x(t) + A_1 x(t-1), \quad t \geq 0, \\ y(t) &= Cx(t), \end{aligned} \tag{2.1}$$

where  $A_0, A_1 \in R^{n \times n}$ ,  $C \in R^{p \times n}$ , with initial condition:  $x(t) = \varphi(t)$ ,  $t \in [-1, 0]$ , where  $\varphi(\cdot) \in \mathcal{F}$ . Two cases will be especially considered:  $\mathcal{F}$  denotes the space  $C([-1, 0], R^n)$  or  $L^p([-1, 0], R^n)$ ,  $1 \leq p \leq \infty$ . A function  $x_t: [-1, 0] \rightarrow R^n$ ,  $x_t(\theta) \triangleq x(t + \theta)$ ,  $\theta \in [-1, 0]$ , can be taken as a state for (2.1) at time  $t$ . In the case when  $\det A_1 = 0$  a state of (2.1) defined above possesses overabundance of information. Therefore we define the true state which contains necessary and sufficient information to solve system equation (2.1) as follows [1]:

DEFINITION 2.1. A pair  $(x(t), A_1 x_t) \triangleq (x(t), A_1 x(t + \theta))$ ,  $\theta \in [-1, 0]$  is said to be a state of system (2.1) at instant  $t$ .

The fundamental definition in this paper is the following.

DEFINITION 2.2. System (2.1) is  $\mathcal{F}$ -finally observable on  $[0, T]$  iff for all  $\varphi(\cdot) \in \mathcal{F}$  such that  $y(t) \equiv 0$ ,  $t \in [0, T]$ , the final trajectory value  $x(T) = 0$ . Definition 2.2 makes sense for  $T \geq 1$  and means that if the system (2.1) is  $\mathcal{F}$ -finally observable on  $[0, T]$  then a unique map  $\{y(t), t \in [0, T]\} \mapsto x(T)$  exists, where  $y(\cdot)$  is generated by the initial function  $\varphi(\cdot) \in \mathcal{F}$ . Hence it is possible to determine  $x(T)$  knowing the output  $y(t)$  on  $[0, T]$ . The importance of the Definition 2.2 lies in the fact that for  $\mathcal{F}$ -finally observable systems on  $[0, T]$  one is able to determine the final state  $(x(T+1), A_1 x_{T+1})$ . Indeed, for any function space  $\mathcal{F}$  such that  $x_t(\cdot) \in \mathcal{F}$ ,  $t > 0$ , provided  $\varphi \in \mathcal{F}$  we may shift the time variable and start with  $x_\tau(\cdot)$  for some  $\tau > 0$ , as an initial function. Utilizing the stationarity of system (2.1) we get the conclusion that  $\mathcal{F}$ -final observa-

bility on  $[0, T]$  implies  $\mathcal{F}$ -final observability on  $[0, T_1]$ ,  $T_1 > T$ . Hence one can determine  $x(T+\tau)$ ,  $\tau \in [0, 1]$ , from the data  $y(t)$ ,  $t \in [0, T+1]$  and, consequently, to get the function  $x_{T+1}(\cdot)$ . Summarizing these considerations we have:

**PROPOSITION 2.1.**

(i) If system (2.1) is  $\mathcal{F}$ -finally observable on  $[0, T]$ ,  $T \geq 1$ , then it is  $\mathcal{F}$ -finally observable on  $[0, T_1]$  for each  $T_1 > T$  provided  $\mathcal{F}$  is such that  $\varphi \in \mathcal{F}$  implies  $x_t(\cdot) \in \mathcal{F}$ ,  $t > 0$ .

(ii)  $\mathcal{F}$ -final observability of system (2.1) on  $[0, T]$  implies that the function  $x_{T+1}(\cdot)$  and the final state  $(x(T+1), A_1 x_{T+1}(\cdot))$  can be determined from the data  $y(t)$ ,  $t \in [0, T+1]$ , for any  $\varphi \in \mathcal{F}$ .

The statement (ii) is especially important when a linear-quadratic problem for controlled time delay systems is considered. In such a case the optimal control has the feedback form  $u(t) = \mathcal{K} x_t(\cdot)$  [7]. Therefore  $x_t(\cdot)$  has to be computed on-line if only output measurements are available in the system.

**PROPOSITION 2.2.**

(i) If system (2.1) is  $\mathcal{F}$ -finally observable on  $[0, T]$  and  $y(t) = 0$  for all  $t \in [0, T_1]$ ,  $T_1 > T$ , then  $x(t) = 0$  for all  $t \in [T, T_1]$ .

(ii)  $y(t) = 0$  for all  $t \in [0, n]$  implies  $y(t) = 0$  for all  $t \geq 0$

(iii)  $\mathcal{F}$ -final observability on  $[0, T]$ ,  $T > n$ , implies  $\mathcal{F}$ -final observability on  $[0, n]$ .

**Proof:**

(i) Follows directly from Proposition 2.1 (i).

(ii) Let  $\lambda$  denote differential operator:  $(\lambda x)(t) = \frac{dx}{dt}(t)$  and  $z$  denote the right shift operator:  $(zx)(t) = x(t-1)$ . State equations for system (2.1) can be written as follows:  $[(\lambda I - A_0 - A_1 z)x](t) = 0$ ,  $t \geq 0$  or, denoting  $\Delta(\lambda, z) \triangleq [\lambda I - A_0 - A_1 z]$ ,  $[\Delta(\lambda, z)x](t) = 0$ .

The polynomial in  $\lambda, z$  matrix  $\text{adj } \Delta(\lambda, z)$  is nonsingular and  $\text{adj } \Delta(\lambda, z) \cdot \Delta(\lambda, z) = \det \Delta(\lambda, z)$ . Hence the latter implies the equation

$$[\det \Delta(\lambda, z)x](t) = 0, \forall t > t_1 \quad (2.2)$$

where  $t_1$  is such that the left-hand side of (2.2) is well defined. We shall show that  $t_1 = n-1$ . In fact,

$$\det \Delta(\lambda, z) = \lambda^n + a_{n-1}(z)\lambda^{n-1} + \dots + a_1(z)\lambda + a_0(z) \quad (2.3)$$

where  $a_i(z)$ ,  $i=0, \dots, n-1$ , are polynomials in  $z$  and  $\deg(a_i(z)) \leq n-i$ .

Therefore,  $[a_i(z)x](t)$  is defined for  $t \geq n-i-1$ . Since  $(\lambda^n x)(t)$  is well defined for  $t > n-1$ , as it is easily seen by the method of steps, the conclusion follows that (2.2) holds with  $t_1 = n-1$ . Clearly, the output  $y(t)$  satisfies the similar equation:

$$[\det \Delta(\lambda, z)y](t) = 0, t \geq n \quad (2.4)$$



where  $\varphi(\cdot) \in \mathcal{F}$  is an initial state for (2.1). The space of functions  $f(\cdot)$  satisfying condition (3.6) will be denoted  $\mathcal{F}_1$ . Conversely, if  $z_k(s)$  is a solution to (3.1) and (3.2) then using formula (3.3) a continuous trajectory  $x(t)$ ,  $t \in [0, k]$  satisfying equation (2.1) can be obtained.

Using Lemma 3.1 any system (2.1) with delay may be transformed into system without delay equivalent of the latter on time interval  $[0, k]$ . Now, Definition 2.2 of  $\mathcal{F}$ -final observability can be expressed in terms of system (3.1), (3.2).

LEMMA 3.2. Let  $T > 1$ ,  $\mathcal{N} \triangleq \text{Ker } C$ ,  $\mathcal{N}_k \triangleq \mathcal{N} \times \dots \times \mathcal{N}$  ( $k$ -times),  $k$  be an integer satisfying  $k \geq T > k - 1$  and  $z_k(s)$  be a solution of (3.1), (3.2).

System (2.1) is  $\mathcal{F}$ -finally observable on  $[0, T]$  iff for all  $f \in \mathcal{F}_1$  such that for all  $s \in [0, 1]$   $z_k(s) \in \mathcal{N}_k$  implies  $z_k(T - k + 1) \in \mathcal{N}_{k-1} \times \{0\}$ .

Proof.  $z_k(s) \in \mathcal{N}_k$  if and only if  $y(t) \equiv 0$  for  $t \in [0, T]$ . Condition  $x(T) = 0$  is equivalent of  $z_k(T - k + 1) \in \mathcal{N}_{k-1} \times \{0\}$ .

Lemma 3.2 shows that  $\mathcal{F}$ -final observability problem can be regarded as a kind of controllability problem of existing a nonzero trajectory of system (3.1) satisfying (3.2) for which  $z_k(s) \in \mathcal{N}_k$ ,  $s \in [0, 1]$ , and  $z_k(T - k + 1) \notin \mathcal{N}_{k-1} \times \{0\}$ . The latter means that corresponding solution  $x(t)$  of system (2.1) is not equal to zero at  $t = T$ . So system (2.1) is not observable on  $[0, T]$ . To answer the question of the existence of such trajectory properties of some controllability subspaces will be considered.

DEFINITION 3.1 [5], [6]. Let  $A \in R^{p \times p}$ ,  $B \in R^{p \times l}$  and  $\mathcal{R}$  be a subspace in  $R^p$ . The greatest (as of inclusion order) subspace  $\mathcal{S}$  contained in  $\mathcal{R}$  satisfying

$$A\mathcal{S} \subset \mathcal{S} + \text{Im } B \quad (3.7)$$

is called maximal invariant controlled subspace and denoted  $\text{Mic}(A, B, \mathcal{R})$

LEMMA 3.3 [5], [6]. Condition (3.7) is satisfied iff a matrix  $D \in R^{l \times p}$  exists such that

$$(A + BD)\mathcal{S} \subset \mathcal{S} \quad (3.8)$$

DEFINITION 3.2 [5], [6]. Let  $A \in R^{p \times p}$ ,  $B \in R^{p \times l}$ ,  $\mathcal{R}$  be a subspace in  $R^p$  and  $\bar{A} \triangleq A + BD$  for arbitrary  $D$  satisfying (3.8),  $D \in R^{l \times p}$ . The subspace  $\{\bar{A} | \text{Mic}(A, B, \mathcal{R}) \cap \text{Im } B\}$  is called maximal controllability subspace and denoted  $\text{Mcs}(A, B, \mathcal{R})$ .

In [5], [6] are given algorithms for computing  $\text{Mic}(A, B, \mathcal{R})$  and  $\text{Mcs}(A, B, \mathcal{R})$  and properties of both subspaces are studied.

LEMMA 3.4 [2]. Each solution of (3.1) satisfying condition  $z_k(s) \in \mathcal{N}_k$ ,  $s \in [0, 1]$  has the form

$$z_k(s) = e^{\bar{A}^k} z_k(0) + \bar{z}_k(s), \quad s \in [0, 1] \quad (3.9)$$

where  $z_k(0) \in \text{Mic}(A^k, B^k, \mathcal{N}_k)$ ,  $\bar{z}_k(s) \in \text{Mcs}(A^k, B^k, \mathcal{N}_k)$  and  $\bar{A}^k \triangleq A^k + B^k D_k$  for arbitrary  $D_k \in R^{n \times nk}$  satisfying  $\bar{A}^k(\text{Mic}(A^k, B^k, \mathcal{N}_k)) \subset \text{Mic}(A^k, B^k, \mathcal{N}_k)$ .

Furthermore, the set of all points  $z_k(s)$  attainable from  $z_k(0) \in \text{Mic}(A^k, B^k, \mathcal{N}_k)$  equals

$$e^{\bar{A}^k} z_k(0) + \text{Mcs}(A^k, B^k, \mathcal{N}_k).$$

**COROLLARY 3.1.** If  $\text{Mic}(A^k, B^k, \mathcal{N}_k) = \{0\}$  then system (2.1) is  $\mathcal{F}$ -finally observable on  $[0, T]$ ,  $T \geq k$ , for any suitable  $\mathcal{F} \subset L^1$ .

**Proof.** If  $\text{Mic}(A^k, B^k, \mathcal{N}_k) = \{0\}$  then  $\text{Mcs}(A^k, B^k, \mathcal{N}_k) = \{0\}$  and only  $z_k(s) = 0$ ,  $s \in [0, 1]$ , is the solution of the form (3.9). Hence  $y(t) = 0$ ,  $t \in [0, k]$ , implies  $x(k) = 0$ , i.e., system (2.1) is  $\mathcal{F}$ -finally observable on  $[0, k]$ . By Proposition 2.1 (i) system (2.1) is  $\mathcal{F}$ -finally observable on  $[0, T]$ ,  $T > k$ .

The following Lemma will be useful in the sequel.

**LEMMA 3.5.** Let  $A \in R^{n \times p}$ ,  $\mathcal{P}$  be subspace of  $R^p$  and  $P^\perp$  is any matrix with independent rows spanning  $\mathcal{P}^\perp$  (if  $\dim \mathcal{P}^\perp = k$ , then  $P^\perp \in R^{k \times n}$ ). Define set  $A^{-1} \mathcal{P} \triangleq \{x \in R^n: \exists y \in \mathcal{P} \text{ and } y = Ax\}$ . Then the following equality holds:

$$A^{-1} \mathcal{P} = \text{Ker } P^\perp A.$$

**Proof.** Let  $x \in A^{-1} \mathcal{P}$ . Then  $y = Ax$  for some  $y \in \mathcal{P}$ . By definition of  $P^\perp$ ,  $P^\perp y = 0$  which implies  $P^\perp Ax = 0$ , i.e.,  $x \in \text{Ker } P^\perp A$ . Conversely, let  $x \in \text{Ker } P^\perp A$ . Then  $P^\perp Ax = 0$  implies  $Ax \in \mathcal{P}$  which means that  $y \in \mathcal{P}$  exists such that  $y = Ax$ , i.e.,  $x \in A^{-1} \mathcal{P}$ .

#### 4. $L^p$ -Final Observability on $[0, k]$

Now, it will be assumed that  $\mathcal{F} = L^p([-1, 0], R^n)$ ,  $1 \leq p \leq \infty$ . When  $T = 1$  then condition (3.2) is always fulfilled because  $J^1 z_1(1) = 0$ ,  $z_1(0) = x_1(0) = x(0)$  and  $z_1^0 = f(1) = x(0)$ . The system will be  $L^p$ -finally observable on  $[0, 1]$  if each trajectory of (3.1) contained in  $\mathcal{N}: z_1(s) \in \mathcal{N} \forall s \in [0, 1]$  has the property:

$$0 = z_1(1) = e^{\bar{A}^1} z_1(0) + \bar{z}_1(1)$$

$$\forall z_1(0) \in \text{Mic}(A_0, A_1, \mathcal{N}) \quad \text{and} \quad \forall \bar{z}_1(1) \in \text{Mcs}(A_0, A_1, \mathcal{N})$$

This implies, in virtue of nonsingularity of  $e^{\bar{A}^1}$ , that  $\text{Mic}(A_0, A_1, \mathcal{N}) = \{0\}$ . Conversely, if  $\text{Mic}(A_0, A_1, \mathcal{N}) = \{0\}$  then the property above holds since, by Definition 3.2,  $\text{Mcs}(A_0, A_1, \mathcal{N}) \subset \text{Mic}(A_0, A_1, \mathcal{N})$ .

Thus we get the following simple result.

**THEOREM 4.1:** System (2.1) is  $L^p$ -finally observable on  $[0, 1]$  iff

$$\text{Mic}(A_0, A_1, \mathcal{N}) = \{0\}.$$

Now the case when  $T=k$  and  $k$  is an integer  $k \geq 2$  will be considered. In accordance with Lemma 3.4 and Lemma 3.2 the question of  $L^p$ -final observability is a question whether all trajectories having the form:

$$z_k(s) = e^{\bar{A}^k} z_k(0) + \tilde{z}_k(s)$$

for

$$z_k(0) = z_k^0 + J^k z_k(1) \quad (4.1)$$

where  $z_k(0) \in \text{Mic}(A^k, B^k, \mathcal{N}_k)$ ,  $\tilde{z}_k(s) \in \text{Mcs}(A^k, B^k, \mathcal{N}_k)$ ,  $s \in [0, 1]$ , fulfil the condition:

$$z_k(1) \in \mathcal{N}_{k-1} \times \{0\}. \quad (4.2)$$

Using the equality  $z_k(1) = e^{\bar{A}^k} z_k(0) + \tilde{z}_k(1)$  condition (4.1) can be written as follows:

$$[I_{kn} - J^k e^{\bar{A}^k}; -J^k] \begin{bmatrix} z_k(0) \\ \tilde{z}_k(1) \end{bmatrix} = z_k^0. \quad (4.3)$$

Let  $M_k, P_k$  be matrices whose linearly independent columns span subspaces  $\text{Mic}(A^k, B^k, \mathcal{N}_k)$  and  $\text{Mcs}(A^k, B^k, \mathcal{N}_k)$  respectively. If  $m = \dim \text{Mic}(A^k, B^k, \mathcal{N}_k)$  and  $j = \dim \text{Mcs}(A^k, B^k, \mathcal{N}_k)$  then  $M_k \in R^{nk \times m}$ ,  $P_k \in R^{nk \times j}$ .

From (3.4) and definition of matrices  $M_k$  and  $P_k$  condition (4.3) can have the form

$$[(I_{kn} - J^k e^{\bar{A}^k}) M_k; -J^k P_k] \begin{bmatrix} z_1 \\ z_2 \end{bmatrix} = z_k^0 \in R^n \times \{0\}_{k-1}, \quad z^1 \in R^m, \quad z^2 \in R^j, \quad (4.4)$$

where  $\{0\}_{k-1}$  is the zero subspace in  $R^{n(k-1)}$ .

Similarly condition (4.2) is equivalent of:

$$z_k(1) = [e^{\bar{A}^k} M_k; P_k] \begin{bmatrix} z_1 \\ z_2 \end{bmatrix} \in \mathcal{N}_{k-1} \times \{0\}, \quad z^1 \in R^m, \quad z^2 \in R^j. \quad (4.5)$$

The system (2.1) will be  $L^p$ -finally observable on  $[0, k]$  iff for all vectors  $\begin{bmatrix} z_1 \\ z_2 \end{bmatrix}$  for which (4.4) holds, condition (4.5) is simultaneously fulfilled.

The set of vectors  $\begin{bmatrix} z_1 \\ z_2 \end{bmatrix} \in R^m \times R^j$  for which (4.4) holds is according to Lemma 3.5 given by

$$\left\{ \begin{bmatrix} z_1 \\ z_2 \end{bmatrix} \in R^m \times R^j : [(I_{kn} - J^k e^{\bar{A}^k}) M_k; -J^k P_k] \begin{bmatrix} z_1 \\ z_2 \end{bmatrix} \in R^n \times \{0\}_{k-1} \right\} = \\ = \text{Ker}(\hat{I}_k ([I_{nk} - J^k e^{\bar{A}^k}] M_k; -J^k P_k))$$

where  $\hat{I}_k \in R^{n(k-1) \times nk}$  is a matrix whose linearly independent rows span the subspace  $(R^n \times \{0\}_{k-1})^\perp = \{0\} \times R^{n(k-1)}$ . For example:

$$\hat{I}_k = \begin{bmatrix} 0 & & I_n & 0 & \dots & 0 \\ & & & & & \vdots \\ \vdots & & & & & 0 \\ 0 & \dots & 0 & & & I_n \end{bmatrix}$$

The  $L^p$ -final observability is equivalent of the following inclusion

$$[e^{\bar{A}^k} M_k \mid P_k] (\text{Ker} (\hat{f}_k [(I_{nk} - J^k e^{\bar{A}^k}] M_k \mid -J^k P_k))) \subset \mathcal{N}_{k-1} \times \{0\}. \tag{4.6}$$

Using again Lemma 3.5 the latter inclusion is fulfilled iff

$$\text{Ker} (\hat{f}_k [(I_{nk} - J^k e^{\bar{A}^k}] M_k \mid -J^k P_k]) \subset \text{Ker} (N_k^\perp (e^{\bar{A}^k} M_k \mid i P_k)) \tag{4.7}$$

where  $N_k^\perp \in R^{i \times nk}$  is a matrix whose linearly independent rows span the subspace  $(\mathcal{N}_{k-1} \times \{0\})^\perp = (\text{Im } C^T)_{k-1} \times R^n$  and  $i \triangleq \dim (\mathcal{N}_{k-1} \times \{0\})^\perp$ . If  $\text{rank } C = p$  (what is usually assumed)  $i = p(k-1) + n$  and the following matrix can be chosen as the

$$N_k^\perp = \begin{bmatrix} C & & & \\ & \ddots & & \\ & & C & \\ & & & 0 \\ 0 & & & & I_n \end{bmatrix}. \tag{4.8}$$

Denote:  $\Phi_k = \hat{f}_k ((I_{nk} - J^k e^{\bar{A}^k}] M_k \mid -J^k P_k)$ .

Inclusion (4.6) can be rewritten as the following rank condition:

$$\text{rank} [\Phi_k^T] = \text{rank} [\Phi_k^T \mid (N_k^\perp e^{\bar{A}^k} M_k \mid N_k^\perp P_k)^T].$$

Hence the following theorem was proved:

**THEOREM 4.2.** System (2.1) is  $L^p$ -finally observable on  $[0, k]$   $k$ - an integer  $k \geq 2$  iff

$$\text{rank} [\Phi_k^T] = \text{rank} [\Phi_k^T \mid (N_k^\perp e^{\bar{A}^k} M_k \mid N_k^\perp P_k)^T].$$

*Example.* As the simple example of the consideration given above the following system equation (2.1) will be examined.

$$\begin{aligned} \dot{x}(t) &= \begin{bmatrix} 0 & 1 \\ -1 & 2 \end{bmatrix} x(t) + \begin{bmatrix} -1 & 1 \\ 0 & 1 \end{bmatrix} x(t-1), & t \geq 0. \\ y(t) &= [0 \ 1] x(t). \end{aligned} \tag{4.9}$$

If  $k=1$  then  $\text{Ker } C = \mathcal{N} = \left\{ n: n = \begin{bmatrix} n_1 \\ 0 \end{bmatrix} \text{ and } n_1 \in R \right\}$ . It is easily to check that  $\text{Mic} (A_0, A_1, \mathcal{N}) \neq \{0\}$  then applying Theorem 4.1 system (4.9) is not  $L^p$ -finally observable on  $[0, 1]$ .

Now, let  $k=2$ . By (3.5) the matrices  $A^2, B^2$  are given by:

$$A^2 = \begin{bmatrix} 0 & 1 & 0 & 0 \\ -1 & 2 & 0 & 0 \\ -1 & 1 & 0 & 1 \\ 0 & 1 & -1 & 2 \end{bmatrix}, \quad B^2 = \begin{bmatrix} -1 & 1 \\ 0 & 1 \\ 0 & 0 \\ 0 & 0 \end{bmatrix}$$

$\mathcal{N}_2 = \mathcal{N} \times \mathcal{N} = \left\{ n: n = \begin{bmatrix} n_1 \\ 0 \\ n_3 \\ 0 \end{bmatrix} \text{ and } n_1, n_3 \in R \right\}$ . By inclusion (3.7) we can find that

$$\text{Mic} (A^2, B^2, \mathcal{N}_2) = \{0\}.$$

By Corollary 3.1 the system (4.9) is  $L^p$ -finally observable on  $[0, 2]$ .



### 5. C-Final Observability on $[0, k]$

In this part of the paper the case  $\mathcal{F} = C([-1, 0], R^n)$  will be considered. Let  $k$  denote a positive integer and  $k \geq 2$ . (The case  $k=1$  will be a simple corollary from theorem for  $k \geq 2$ ).

Each continuous function  $f(\cdot)$ , for which the trajectory of equation (3.1) is included in  $\mathcal{N}_k$  for all  $s \in [0, 1]$  has the form [6]:

$$f(s) = D_k z_k(s) + G_k v(s), \quad s \in [0, 1] \quad (5.1)$$

where  $D_k \in R^{n \times nk}$  is arbitrary satisfying  $(A^k + B^k D_k) (\text{Mic}(A^k, B^k, \mathcal{N}_k)) \subset \text{Mic}(A^k, B^k, \mathcal{N}_k)$  and  $G_k \in R^{n \times n}$  satisfies

$$\text{Im } B^k G_k = \text{Mic}(A^k, B^k, \mathcal{N}_k) \cap \text{Im } B^k.$$

$v(\cdot)$  is a function belonging to  $C([0, 1], R^n)$ .

According to Lemmas 3.2 and 3.4 the system (2.1) will be C-finally observable on  $[0, k]$  iff

$$z_k(1) \in \mathcal{N}_{k-1} \times \{0\}$$

for each

$$z_k(1) = e^{\bar{A}^k} z_k(0) + \bar{z}_k(1), \quad z_k(0) \in \text{Mic}(A^k, B^k, \mathcal{N}_k), \quad \bar{z}_k(1) \in \text{Mcs}(A^k, B^k, \mathcal{N}_k) \quad (5.2)$$

where from (3.2)

$$z_k(0) = z_k^0 + J^k z_k(1) \quad (5.3)$$

and from (5.1)

$$f(1) = D_k z_k(1) + G_k \cdot v(1). \quad (5.4)$$

Using the definition of  $z_k^0$  given by (3.4) and defining  $I^k \in R^{nk \times n}$  as

$$I^k \triangleq \begin{bmatrix} I_n \\ 0 \\ \vdots \\ 0 \end{bmatrix}$$

condition (5.4) can be formulated as follows:

$$z_k^0 = I^k D_k z(1) + I^k G_k v(1). \quad (5.5)$$

From (5.3) and (5.5) we have

$$[I_{kn} - (I^k D_k + J^k) e^{\bar{A}^k}] z_k(0) = (I^k D_k + J^k) \bar{z}_k(1) + I^k G_k v(1). \quad (5.6)$$

Let  $A_k \in R^{nk \times nk}$ ,  $\Omega_k \in R^{nk \times nk}$  be defined by:

$$\begin{aligned} A_k &\triangleq I_{nk} - (I^k D_k + J^k) e^{\bar{A}^k}, \\ \Omega_k &\triangleq I^k D_k + J^k, \end{aligned} \quad (5.7)$$

and as before  $M_k, P_k$  are matrices whose linearly independent columns span subspaces  $\text{Mic}(A^k, B^k, \mathcal{N}_k)$  and  $\text{Mcs}(A^k, B^k, \mathcal{N}_k)$  respectively. If  $m \triangleq \dim \text{Mic}(A^k, B^k, \mathcal{N}_k)$  and  $j \triangleq \dim \text{Mcs}(A^k, B^k, \mathcal{N}_k)$  then  $M_k \in R^{nk \times m}$ ,  $P_k \in R^{nk \times j}$ . Using this definitions (5.6) is equivalent of:

$$[A_k M_k \mid -\Omega_k P_k] \cdot \begin{bmatrix} z^1 \\ z^2 \end{bmatrix} = I^k G v \quad (1) \quad \text{where} \quad \begin{bmatrix} z^1 \\ z^2 \end{bmatrix} \in R^m \times R^j. \quad (5.8)$$

System (2.1) is  $C$ -finally observable iff condition (5.8) implies

$$[e^{\bar{A}^k} M_k \mid P_k] \cdot \begin{bmatrix} z^1 \\ z^2 \end{bmatrix} \in \mathcal{N}_{k-1} \times \{0\}. \quad (5.9)$$

By Lemma 3.5 this implication can be transformed into the form

$$\text{Ker} [\Sigma_k (A_k M_k \mid -\Omega_k P_k)] \subset \text{Ker} [N_k^\perp (e^{\bar{A}^k} M_k \mid P_k)] \quad (5.10)$$

where  $N_k^\perp$  is defined by (4.8) and  $\Sigma_k$  is a matrix whose linearly independent rows span subspace:

$$(\text{Im } I^k G_k)^\perp = \text{Ker} (I^k G_k)^T = \text{Ker } G_k^T (I^k)^T.$$

Inclusion (5.10) may be transformed into some equivalent rank condition:

$$\begin{aligned} \text{rank} [\Sigma_k A_k M_k \mid -\Sigma_k \Omega_k P_k]^T &= \\ &= \text{rank} [(\Sigma_k A_k M_k \mid -\Sigma_k \Omega_k P_k)^T \mid (N_k^\perp e^{\bar{A}^k} M_k \mid N_k^\perp P_k)^T] \end{aligned}$$

Hence the following theorem was proved:

**THEOREM 5.1.** System (2.1) is  $C$ -finally observable on  $[0, k]$ ,  $k$  an integer  $k \geq 2$  iff:

$$\begin{aligned} \text{rank} [\Sigma_k A_k M_k \mid -\Sigma_k \Omega_k P_k]^T &= \\ &= \text{rank} [(\Sigma_k A_k M_k \mid -\Sigma_k \Omega_k P_k)^T \mid (N_k^\perp e^{\bar{A}^k} M_k \mid N_k^\perp P_k)^T]. \quad (5.11) \end{aligned}$$

**COROLLARY 5.1.** System (2.1) is  $C$ -finally observable on  $[0, 1]$  iff

$$\begin{aligned} \text{rank} [\Sigma_1 M_1 - \Sigma_1 D_1 e^{\bar{A}^1} M_1 \mid -\Sigma_1 D_1 P_1]^T &= \\ &= \text{rank} [(\Sigma_1 M_1 - \Sigma_1 D_1 e^{\bar{A}^1} M_1 \mid -\Sigma_1 D_1 P_1)^T \mid (e^{\bar{A}^1} M_1 \mid P_1)^T] \quad (5.12) \end{aligned}$$

where  $\Sigma_1$  is a matrix whose linearly independent rows span  $\text{Ker } G_1$  and  $M_1, P_1$  are matrices whose linearly independent columns span  $\text{Mic}(A_0, A_1, \mathcal{N})$  and  $\text{Mcs}(A_0, A_1, \mathcal{N})$  respectively.

**Proof.** if  $k=1$  then  $A_1 = I_n - D_1 e^{\bar{A}^1}$ ,  $\Omega_1 = D_1$ , condition (5.9) has the following form:

$$[e^{\bar{A}^1} M_1 \mid P_1] \cdot \begin{bmatrix} z^1 \\ z^2 \end{bmatrix} = 0$$

what implies  $N_1^\perp = I_n$ .

*Example.* We will consider system given by equations (4.9). We have:

$$A_0 = \begin{bmatrix} 0 & 1 \\ -1 & 2 \end{bmatrix}, \quad A_1 = \begin{bmatrix} -1 & 1 \\ 0 & 1 \end{bmatrix}, \quad C = [0 \ 1].$$

Let  $k=1$ . Then  $\mathcal{N}_1 = \mathcal{N} = \left\{ n: n = \begin{bmatrix} n_1 \\ 0 \end{bmatrix} \text{ and } n_1 \in R \right\}$ . By definitions (3.5) of  $A^1, B^1$  and  $J^1$  we get:

$$A^1 = A_0, \quad B^1 = A_1, \quad J^1 = 0.$$

Using inclusion (3.8) we can find  $\text{Mic}(A_0, A_1, \mathcal{N}) = \mathcal{N} = \left\{ n: n = \begin{bmatrix} n_1 \\ 0 \end{bmatrix} \text{ and } n_1 \in R \right\}$ .

The matrix  $D_1$  satisfying (3.8) can be chosen as

$$D_1 = \begin{bmatrix} 0 & 0 \\ 1 & 0 \end{bmatrix}$$

what implies  $\bar{A}^1 = A_0 + A_1 D_1 = \begin{bmatrix} 1 & 1 \\ 0 & 2 \end{bmatrix}$  and  $e^{\bar{A}^1} = \begin{bmatrix} e & e^2 - e \\ 0 & e^2 \end{bmatrix}$ .

To find subspace  $\text{Mcs}(A_0, A_1, \mathcal{N})$  we observe that

$$\text{Im } A_1 \cap \text{Mic}(A_0, A_1, \mathcal{N}) = \text{Mic}(A_0, A_1, \mathcal{N})$$

and by inclusion (3.8) which shows that  $\text{Mic}(A_0, A_1, \mathcal{N})$  is  $\bar{A}^1$ -invariant we get:

$$\text{Mcs}(A_0, A_1, \mathcal{N}) = \text{Mic}(A_0, A_1, \mathcal{N}).$$

The matrix  $G_1$ , which fulfils condition

$$\text{Im } A_1 \cap \text{Mic}(A_0, A_1, \mathcal{N}) = \text{Im } A_1 G_1$$

can be chosen as:

$$G_1 = \begin{bmatrix} 1 & 0 \\ 0 & 0 \end{bmatrix}$$

$$\text{Ker } G_1^T = \text{Ker } G_1 = \left\{ g: g = \begin{bmatrix} 0 \\ g_2 \end{bmatrix} \text{ and } g_2 \in R \right\}.$$

Because subspace  $\text{Mic}(A_0, A_1, \mathcal{N})$ ,  $\text{Mcs}(A_0, A_1, \mathcal{N})$  and  $\text{Ker } G_1^T$  are determined we can choose matrices  $M_1, P_1, \Sigma_1$  as:

$$P_1 = M_1 = \begin{bmatrix} 1 \\ 0 \end{bmatrix}, \quad \Sigma_1 = [0 \ 1]$$

So:

$$[\Sigma_1 M_1 - \Sigma_1 D_1 e^{\bar{A}^1} M_1 \quad -\Sigma_1 D_1 P_1] = [-e \quad -1],$$

$$[e^{\bar{A}^1} M_1 \quad P_1] = \begin{bmatrix} e & 1 \\ 0 & 0 \end{bmatrix}.$$

The rank condition (5.12) is fulfilled because:

$$\text{rank} \begin{bmatrix} -e \\ -1 \end{bmatrix} = \text{rank} \begin{bmatrix} -e & e & 0 \\ -1 & 1 & 0 \end{bmatrix}.$$

Therefore system given by equation (4.9) is  $C$ -finally observable on  $[0, 1]$ .

## 6. Concluding Remarks

The observability criteria given in this paper can be generalized for system (2.1) with delay  $h \neq 1$  using substitution  $t = \tau h$  and  $\tilde{x}(\tau) = x(\tau h)$  or for many comensurable delays as well. Theorems 4.2 and 5.1 are applied for testing  $\mathcal{F}$ -final observability on  $[0, k]$  where  $k$  is a positive integer. By Lemmas given above one can prove conditions of  $\mathcal{F}$ -final observability on  $[0, T]$ ,  $T > 1$ . According to Proposition 2.1 (i) the necessary condition is  $\mathcal{F}$ -final observability on  $[0, k]$  where  $k$  fulfils inequality  $k - 1 < T \leq k$ , moreover the rank condition of some matrices must be fulfilled. For systems (2.1) with greater dimension  $n$  and when the observation interval is long it seems to be necessary to use computers to test observability. In this case some numerical problems with determination of subspaces and computation of exponents of matrices can appear.

The important problem which has not been solved yet is construction of observer determining  $x(T)$  using output information  $y(t)$ ,  $t \in [0, T]$ .

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### Obserwowalność końcowa układów z opóźnieniem

Zdefiniowano i zanalizowano pojęcie obserwowalności „końcowej”, na przedziale  $[0, T]$  dla układów ze stałym opóźnieniem. Sformułowano i udowodniono warunki obserwowalności zależne od długości przedziału obserwacji. Rozpatrzono dwa rodzaje obserwowalności: w pierwszym zakłada się, że funkcja początkowa jest klasy  $L^p$ , w drugim przyjmuje się, że funkcja ta jest ciągła.

### Окончательная наблюдаемость систем с запаздыванием

В работе определяется и анализируется понятие окончательной наблюдаемости в интервале  $[0, T]$  для систем с постоянным запаздыванием. Формулируются и доказываются условия наблюдаемости, зависящие от величины интервала наблюдений. Рассмотрены два вида наблюдаемости: в первом предполагается, что начальная функция определена в пространстве  $L^p$ , во втором — эта функция непрерывна.

