# A Study of the Observability of Multidimensional Hybrid Linear Systems 

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#### Abstract

A class of multidimensional hybrid linear systems is presented, with the time vector composed by $q$ continuous-time real components and by $r$ discrete-time integer ones, $q, r \geq 1$. The state equation is of multidimensional partial differential-difference type. A generalized variation-of-parameters formula is provided and it is used to obtain the state and the general response of the system. The fundamental concept of observability is studied for these systems. An observability Gramian is introduced, which is a generalization of the Gramians corresponding to the classical 1D continuous-time and 1D discrete-time systems. In the case of completely observable systems this Gramian is used to obtain a formula which provides the initial state of the system for any input-output pair. A list of observability criteria is given for time-invariant systems and the duality between the concepts of observability and reachability is emphasized.


Key-Words: multidimensional hybrid systems, observability, observability Gramian, time-varying systems, timeinvariant systems.

## 1 Introduction

In the past three decades a lot of published paper and books have been designed to the theory of twodimensional (2D) or, more generally, multidimensional ( $n \mathrm{D}$ ) systems, which become a distinct and important branch of the systems theory. The reasons for the increasing interest in this domain are on one side the richness in potential application fields and on the other side the richness and significance of the theoretical approaches. The list of application fields include circuits, control and signal processing, image processing, computer tomography, gravity and magnetic field mapping, seismology, control of multipass processes, etc. From the theoretical point of view, the domain of $n \mathrm{D}$ systems needs a specific approach, since many aspects of the 1 D systems do not generalize and there are many $n \mathrm{D}$ systems phenomena which have no 1D systems counterparts. Various state space 2D discrete-time models have been proposed in literature by Roesser [18], Fornasini-Marchesini [6] Attasi [4], Eising [5] and others. Apparently different, in fact some of them are equivalent.

A quite new subdomain of the 2 D systems theory is represented by the 2 D hybrid models, whose state equation is of differential-difference type, having a continuous-time variable and a discrete-time one
[9], [14], [16], [17]. These hybrid models have applications in various areas such as linear repetitive processes [8], [19], pollution modelling [7], long-wall coal cutting and metal rolling [20] or in iterative learning control synthesis [13].

The concept of observability which is fundamental in control theory was introduced for 1D systems by Kalman in [10], being imposed by engineering problems (see [11]). This concept and its dual (controllability) characterize the minimal systems. In the case of 2D systems some notions such as local and global controllability and observability were introduced, but they are not satisfactory from the point of view of minimality. In [12] the concepts of modal controllability/observability were defined and it was shown that a system is minimal iff it is modally controllable and observable, but these notions do not allow the richness of the characterizations of the 1D notions.

In the present paper the notion of observability is defined and it is analysed for a model of time-variable separable ( $q, r$ )-D hybrid (i.e. continuous-discrete) systems, which is a generalization of the class of 2D hybrid system studied in [15], this class being the continuous-discrete time-varying counterpart of Attasi's 2D discrete-time time-invariant model [4].

In Section 2 a variation-of-parameters formula is
provided and it is used to obtain the formula of the state of the system, determined by an initial state and a given control, as well as the formula of the general response of the system.

Section 3 analyses the notions of unobservable states and completely observable systems. A ( $q, r$ )-D (multidimensional) observability Gramian is defined, which is a natural generalization of the Gramians corresponding to the 1-D continuous-time and discretetime systems. Using this Gramian necessary and sufficient conditions of observability for time-varying systems are given, as well as the characterization of the subspace of unobservable states. In the case of completely observable systems a formula is provided which determines the initial state by knowing the control and the corresponding output.

Section 5 is devoted to time-invariant $(q, r)$-D hybrid systems. A list of necessary and sufficient conditions of observability is provided. The geometric characterization of the subspace of unobservable states is given in terms of invariant subspaces included in the kernel of the output matrix.

This study can be continued in many directions such as stability, positivity, 2D generalized systems, linear quadratic optimal control etc., by extending to the present approach some results presented for instance in [1], [2] and [3]. Another direction is represented by the application of this framework to the geometric approach of multitime systems studied in [21]-[25].

We shall use the following notations: $q \in \mathbf{N}$ and $r \in \mathbf{N}$ being the number of continuous and discrete variables respectively, a function $x\left(t_{1}, \ldots, t_{q}\right.$; $\left.k_{1}, \ldots, k_{r}\right), t_{i} \in \mathbf{R}, k_{i} \in \mathbf{Z}$ will be sometimes denoted by $x(t ; k)$, where $t=\left(t_{1}, \ldots, t_{q}\right), k=$ $\left(k_{1}, \ldots, k_{r}\right)$. By $\bar{m}$ with $m \in \mathbf{N}^{*}$ we denote the set $\{1,2, \ldots, m\}$ and by $\mathcal{P}(\bar{m})$ the family of all subsets of $\bar{m}$.

By $s \leq t,(s<t), s, t \in \mathbf{R}^{q}$ we mean $s_{i} \leq$ $t_{i},\left(s_{i}<t_{i}\right) \forall i \in \bar{q}$ and a similar signification has $l \leq k, l, k \in \mathbf{Z}^{r} ;(s ; l)<(t ; k)$ means $s \leq t, l \leq k$ and $(s ; l) \neq(t ; k)$. For $t^{0}, t^{1} \in \mathbf{R}^{q}$ and $k^{0}, k^{\overline{1}} \in$ $\mathbf{Z}^{r}, t^{0}<t^{1}, k^{0}<k^{1}$ we denote by $\left[t^{0}, t^{1}\right]$ and $\left[k^{0}, k^{1}\right]$ respectively the sets $\left[t^{0}, t^{1}\right]=\prod_{i=1}^{q}\left[t_{i}^{0}, t_{i}^{1}\right]$ and $\left[k^{0}, k^{1}\right]=\prod_{i=1}^{r}\left\{k_{j}^{0}, k_{j}^{0}+1, \ldots, k_{j}^{1}\right\}$.

If $\tau=\left\{i_{1}, \ldots, i_{l}\right\}$ is a subset of $\bar{m},|\tau|_{\sim}:=l$ and $\tilde{\tau}:=\bar{m} \backslash \tau ;$ for $i \in \bar{m}, \tilde{i}:=\bar{m} \backslash\{i\}$ and $\bar{i}:=\{i+$ $1, \ldots, m\}$. The notation $(\tau, \delta) \subset(\bar{q}, \bar{r})$ means that $\tau$ and $\delta$ are subsets of $\bar{q}$ and $\bar{r}$ respectively and $(\tau, \delta) \neq$ $(\bar{q}, \bar{r})$. For $\tau=\left\{i_{1}, \ldots, i_{l}\right\}$ and $\delta=\left\{j_{1}, \ldots, j_{h}\right\}$ the operators $\frac{\partial}{\partial \tau}$ and $\sigma_{\delta}$ are defined by
$\frac{\partial}{\partial \tau} x(t ; k)=\frac{\partial^{l}}{\partial t_{i_{1}} \ldots \partial t_{i_{l}}} x(t ; k), \sigma_{\delta} x(t ; k)=x\left(t ; k+e_{\delta}\right)$
where $e_{\delta}=e_{j_{1}}+\ldots+e_{j_{h}}, e_{j}=$ $(\underbrace{0, \ldots, 0}_{j-1}, 1,0, \ldots, 0) \in \mathbf{R}^{r}$; when $\tau=\bar{q}$ and $\delta=\bar{r}$ we denote $\partial / \partial \tau=\partial / \partial t$ and $\sigma_{\delta}=\sigma$.

If $A_{i}, i \in \bar{m}$ is a family of matrices, $\sum_{i \in \emptyset} A_{i}=0$ and $\prod_{i \in \emptyset} A_{i}=I$.

If $P$ is a positive definite matrix, one writes $P>$ 0.

## 2 State space representation of ( $q, r$ )-D CDSs

The time set of the Attasi-type multidimensional system is $T=\mathbf{R}^{q} \times \mathbf{Z}^{r}, q, r \in \mathbf{N}^{*}$.

Definition 1. A $(q, r)$-D continuous-discrete system (CDS) is a set $\Sigma=\left(\left\{A_{c i} \mid i \in \bar{q}\right\},\left\{A_{d j} \mid j \in\right.\right.$ $\bar{r}\}, B, C, D)$ with $A_{c i}(t ; k), i \in \bar{q}$ and $A_{d j}(t ; k), j \in \bar{r}$ commuting $n \times n$ matrices $\forall t \in \mathbf{R}^{q}, \forall k \in \mathbf{Z}^{r}$ and $B(t ; k), C(t ; k), D(t ; k)$ respectively $n \times m, p \times n$ and $p \times m$ real matrices, all these matrices being continuous with respect to $t \in \mathbf{R}^{q}$ for any $k \in \mathbf{Z}^{r}$; the state equation is

$$
\begin{gather*}
\frac{\partial}{\partial t} \sigma x(t ; k)=\sum_{(\tau, \delta) \subset(\bar{q}, \bar{r})}(-1)^{q+r-|\tau|-|\delta|-1} \times \\
\times\left(\prod_{i \in \tilde{\tau}} A_{c i}(t ; k)\right)\left(\prod_{j \in \tilde{\delta}} A_{d j}(t ; k)\right) \frac{\partial}{\partial \tau} \sigma_{\delta} x(t ; k)+  \tag{1}\\
+B(t ; k) u(t ; k)
\end{gather*}
$$

and the output equation is

$$
\begin{equation*}
y(t ; k)=C(t ; k) x(t ; k)+D(t ; k) u(t ; k) \tag{2}
\end{equation*}
$$

where

$$
x(t ; k)=x\left(t_{1}, \ldots, t_{q} ; k_{1}, \ldots, k_{r}\right) \in \mathbf{R}^{n}
$$

is the state, $u(t ; k) \in \mathbf{R}^{m}$ is the input and $y(t ; k) \in$ $\mathbf{R}^{p}$ is the output. $X=\mathbf{R}^{n}, U=\mathbf{R}^{n}$ and $Y=\mathbf{R}^{p}$ are respectively the state space, the input space and the output space .

For $\tau=\left\{i_{1}, \ldots, i_{l}\right\} \subset \bar{q}, \delta=\left\{j_{1}, \ldots, j_{h}\right\} \subset \bar{r}$ and $t_{i} \in \mathbf{R}, i \in \tau, t_{i}^{0} \in \mathbf{R}, i \in \tilde{\tau}, k_{j} \in \mathbf{Z}, j \in \delta$, $k_{j}^{0} \in \mathbf{Z}, j \in \tilde{\delta}$ we use the notation

$$
\begin{gathered}
x\left(t_{\tau}, t_{\tilde{\tau}}^{0} ; k_{\delta}, k_{\tilde{\delta}}^{0}\right):=x\left(t_{1}^{0}, \ldots, t_{i_{1}-1}^{0}, t_{i_{1}}\right. \\
t_{i_{1}+1}^{0}, \ldots, t_{i_{l}-1}^{0}, t_{i_{l}}, t_{i_{l}+1}^{0}, \ldots, t_{q}^{0} \\
k_{1}^{0}, \ldots, k_{j_{1}-1}^{0}, k_{j_{1}}, k_{j_{1}+1}^{0}, \ldots, k_{j_{h}-1}^{0}, k_{j_{h}} \\
\left.k_{j_{h}+1}^{0}, \ldots, k_{j_{r}}^{0}\right)
\end{gathered}
$$

Let $\Phi_{i}\left(t_{i}, t_{i}^{0} ; t_{\tilde{i}} ; k\right)$ be the (continuous) fundamental matrix of $A_{c i}(t ; k)$ with respect to the variables $t_{i}, t_{i}^{0}, i \in \bar{q}$, i.e. the unique matrix solution of the system $\frac{\partial Y}{\partial t_{i}}(t ; k)=A_{c i}(t ; k) Y(t ; k), Y\left(t_{\tilde{i}}, t_{i}^{0} ; k\right)=I$ for any $t_{l} \in \mathbf{R}, l \in \tilde{i}$ and $k \in \mathbf{R}^{r}$. If $A_{c i}$ is a constant matrix then $\Phi_{i}\left(t_{i}, t_{i}^{0} ; t_{\tilde{i}}, k\right)=e^{A_{c i}\left(t_{i}-t_{i}^{0}\right)}$.

The discrete fundamental matrix $F_{j}\left(t ; k_{j}, k_{j}^{0} ; k_{\tilde{j}}\right)$ of the matrix $A_{d j}(t ; k)$ is defined by

$$
\begin{aligned}
& F_{j}\left(t ; k_{j}, k_{j}^{0} ; k_{\tilde{j}}\right)= \\
& =\left\{\begin{array}{r}
A_{d j}\left(t ; k_{j}-1, k_{\tilde{j}}\right) A_{d j}\left(t ; k_{j}-2, k_{\tilde{j}}\right) \ldots A_{d j}\left(t ; k_{j}^{0}, k_{\tilde{j}}\right) \\
\text { for } k_{j}>k_{j}^{0} \\
I_{n} \quad \text { for } \quad k=k_{j}^{0},
\end{array}\right.
\end{aligned}
$$

for any $k_{h} \in \mathbf{Z}, h \in \tilde{j}$ and $t \in \mathbf{R}^{q}$.
$F_{j}\left(t ; k_{j}, k_{j}^{0} ; k_{j}\right)$ is the unique matrix solution of the difference system

$$
\begin{aligned}
& Y\left(t ; k_{j}+1, k_{\tilde{j}}\right)=A_{d j}(t ; k) Y\left(t ; k_{j}, k_{\tilde{j}}\right) \\
& Y\left(t ; k_{j}^{0} ; k_{\tilde{j}}\right)=I
\end{aligned}
$$

If $A_{d j}$ is a constant matrix then

$$
F_{j}\left(t ; k_{j}, k_{j}^{0} ; k_{\tilde{j}}\right)=A_{d j}^{k_{j}-k_{j}^{0}}
$$

Remark 1. Under the hypothesis: (H) "The matrices $A_{d j}(t ; k)$ are nonsingular for any $t \in \mathbf{R}^{q}$, $k \in \mathbf{Z}^{r}$ ", the discrete fundamental matrix $F_{j}$ can be defined for $k_{j}<k_{j}^{0}$ by

$$
F_{j}\left(t ; k_{j}, k_{j}^{0} ; k_{\tilde{j}}\right)=F_{j}\left(t ; k_{j}^{0}, k_{j} ; k_{\tilde{j}}\right)^{-1}
$$

In this case the semigroup property

$$
F_{j}\left(t ; k_{j}, k_{j}^{1} ; k_{\tilde{j}}\right) F_{j}\left(t ; k_{j}^{1}, k_{j}^{0} ; k_{\tilde{j}}\right)=F_{j}\left(t ; k_{j}, k_{j}^{0} ; k_{\tilde{j}}\right)
$$

holds for any $k_{j}^{0}, k_{j}^{1}, k_{j}$.
Definition 2. The vector $x^{0} \in \mathbf{R}^{n}$ is called an initial state of the system $\Sigma$ if

$$
\begin{align*}
& x\left(t_{\tau}, t_{\tilde{\tau}}^{0} ; k_{\delta}, k_{\tilde{\delta}}^{0}\right)=\left(\prod_{i \in \tau} \Phi_{i}\left(t_{i}, t_{i}^{0} ; t_{\tilde{i}} ; k\right)\right) \cdot \\
& \cdot\left(\prod_{j \in \delta} F_{j}\left(t ; k_{j}, k_{j}^{0} ; k_{\tilde{j}}\right)\right) x^{0} \tag{3}
\end{align*}
$$

for any $(\tau, \delta) \subset(\bar{q}, \bar{r})$; equalities (3) are called the initial conditions of $\Sigma$.

Proposition 1. The solution of the initial value problem

$$
\begin{gather*}
\frac{\partial}{\partial t} \sigma x(t ; k)=\sum_{(\tau, \delta) \subset(\bar{q}, \bar{r})}(-1)^{q+r-|\tau|-|\delta|-1} . \\
\cdot\left(\prod_{i \in \tilde{\tau}}\left(\sigma_{\delta} A_{c i}(t ; k)\right)\right)\left(\prod_{j \in \tilde{\delta}} A_{d j}(t, k)\right)  \tag{4}\\
\cdot \frac{\partial}{\partial \tau} \sigma_{\delta} x(t ; k)+f(t ; k)
\end{gather*}
$$

with the initial conditions (3) is given by the general-
ized variation-of-parameters formula

$$
\begin{align*}
& x(t ; k)=\left(\prod_{i=1}^{q} \Phi_{i}\left(t_{i}, t_{i}^{0} ; t_{\overline{i-1}}^{0}, t_{\tilde{i}} ; k\right)\right) \\
& \left(\prod_{j=1}^{r} F_{j}\left(t^{0} ; k_{j}, k_{j}^{0} ; k_{\overline{j-1}}^{0}, k_{\tilde{j}}\right)\right) x^{0}+  \tag{5}\\
& +\int_{t_{1}^{0}}^{t_{1}} \cdots \int_{t_{q}^{0}}^{t_{q}} \sum_{l_{1}=k_{1}^{0}}^{k_{1}-1} \ldots \sum_{l_{r}=k_{r}^{0}}^{k_{r}-1}\left(\prod_{i=1}^{q} \Phi_{i}\left(t_{i}, s_{i} ; s_{\overline{i-1}}, t_{\tilde{\tilde{i}}} ; k\right)\right) . \\
& \cdot\left(\prod_{j=1}^{r} F_{j}\left(s ; k_{j}, l_{j}+1 ; l_{\overline{j-1}}, k_{\tilde{\tilde{j}}}\right)\right) f(s ; l) d s_{1} \ldots d s_{q} ;
\end{align*}
$$

here $s=\left(s_{1}, \ldots, s_{q}\right), l=\left(l_{1}, \ldots, l_{r}\right)$ and if for instance $i=1$, then the corresponding variable $t_{\overline{i-1}}=$ $t_{\emptyset}^{0}$ lacks; $f: \mathbf{R}^{q} \times \mathbf{Z}^{r} \rightarrow \mathbf{R}^{n}$ is a continuous function with respect to $t \in \mathbf{R}^{q}$ for any $k \in \mathbf{Z}^{r}$.

Proof. We shall prove (5) firstly for the case $q=$ $r=1$, that is for the equation

$$
\begin{align*}
& \frac{\partial}{\partial t} x(t ; k+1)=A_{c}(t ; k+1) x(t ; k+1)+ \\
& +A_{d}(t ; k) \frac{\partial}{\partial t} x(t ; k)-A_{c}(t ; k) A_{d}(t ; k) x(t ; k)+f(t ; k) \tag{6}
\end{align*}
$$

with the initial conditions
$x\left(t ; k_{0}\right)=\Phi\left(t, t_{0} ; k_{0}\right) x_{0}, \quad x\left(t_{0} ; k\right)=F\left(t_{0} ; k, k_{0}\right) x_{0}$
where $\Phi\left(t, t_{0} ; k_{0}\right)$ is the fundamental matrix of $A_{c}(t ; k)$ and $F\left(t ; k, k_{0}\right)$ is the discrete fundamental matrix of $A_{d}(t ; k)$.

Let us consider the vector

$$
\begin{equation*}
z(t ; k)=\frac{\partial}{\partial t} x(t ; k)-A_{c}(t ; k) x(t ; k) \tag{8}
\end{equation*}
$$

From (7) and the first property of the fundamental matrix we have

$$
\begin{aligned}
z\left(t ; k_{0}\right)=\frac{\partial}{\partial t} x\left(t ; k_{0}\right)-A_{c}\left(t ; k_{0}\right) x\left(t ; k_{0}\right) & = \\
=\frac{\partial}{\partial t} \Phi\left(t ; t_{0}, k_{0}\right) x_{0}-A_{c}\left(t ; k_{0}\right) \Phi\left(t ; t_{0}, k_{0}\right) x_{0} & =0
\end{aligned}
$$

hence $z\left(t ; k_{0}\right)=0, \forall t \geq t_{0}$. Equation (6) can be written as the 1D difference equation $z(t ; k+1)=$ $A_{j}(t ; k) z(t ; k)+f(t ; k)$ whose solution, given by the discrete-time variation of parameters formula, is
$z(t ; k)=F\left(t ; k, k_{0}\right) z\left(t, k_{0}\right)+\sum_{l=k_{0}}^{k-1} F(t ; k, l+1) f(t, l)$ hence

$$
\begin{equation*}
z(t, k)=\sum_{l=k_{0}}^{k-1} F(t ; k, l+1) f(t, l) \tag{9}
\end{equation*}
$$

since $z\left(t, k_{0}\right)=0$. But the equation (8) can be written as the differential equation $\frac{\partial}{\partial t} x(t ; k)=$ $A_{c}(t ; k) x(t ; k)+z(t ; k)$ and by the variation of parameters formula its solution is
$x(t ; k)=\Phi\left(t, t_{0} ; k\right) x\left(t_{0} ; k\right)+\int_{t_{0}}^{t} \Phi(t, s ; k) z(s ; k) d s$.
By replacing $x\left(t_{0} ; k\right)$ and $z(s ; k)$ by their expressions (7) and (9) we get

$$
\begin{gather*}
x(t ; k)=\Phi\left(t, t_{0} ; k\right) F\left(t_{0} ; k, k_{0}\right) x_{0}+ \\
\int_{t_{0}}^{t}+\sum_{l=k_{0}}^{k-1} \Phi(t, s ; k) F(s ; k, l+1) f(s, l) \tag{10}
\end{gather*}
$$

hence formula (5) is true for $q=1, r=1$.
Assume that formula (5) is true for some $q, r \in \mathbf{N}$ and consider the equation (4) with $q+1$ instead of $q$. We introduce the function

$$
\begin{gather*}
z\left(t_{1} ; t_{2}, \ldots, t_{q+1} ; k_{1}, \ldots, k_{r}\right)=\frac{\partial x}{\partial t_{1}}\left(t_{1}, t_{2}, \ldots, t_{q+1}\right. \\
\left.k_{1}, \ldots, k_{r}\right)-A_{c 1}\left(t_{1}, t_{2}, \ldots, t_{q+1} ; k_{1}, \ldots, k_{r}\right) \\
\cdot x\left(t_{1}, t_{2}, \ldots, t_{q+1} ; k_{1}, \ldots, k_{r}\right) \tag{11}
\end{gather*}
$$

Then the equation (4) of order $(q+1, r)$ can be written as an equation $(q, r)$ of the same type with $z$ instead of $x$ and the initial conditions (3) give null initial conditions for $z$.

By the induction assumption the corresponding solution is

$$
\begin{align*}
& z\left(t_{1} ; t_{2}, \ldots, t_{q+1} ; k_{1}, \ldots, k_{r}\right)=\int_{t_{2}^{0}}^{t_{2}} \ldots \int_{t_{q+1}^{0}}^{t_{q+1}} \\
& \sum_{l_{1}=k_{1}^{0}}^{k_{1}-1} \ldots \sum_{l_{r}=k_{r}^{0}}^{k_{r}-1}\left(\prod_{i=2}^{q+1} \Phi_{i}\left(t_{i}, s_{i} ; t_{1} ; s_{\overline{i-1} \backslash\{1\}}, t_{\tilde{i}} ; k\right) .\right. \\
& \cdot\left(\prod_{j=1}^{r} F_{j}\left(t_{1} ; s_{2}, \ldots, s_{q+1} ; k_{j}, l_{j}+1 ; l_{\overline{j-1}} ; k_{\tilde{j}}\right)\right. \\
& \cdot f(s ; l) d s_{2} \ldots d s_{q+1} \tag{12}
\end{align*}
$$

But the variation of parameters formula gives the solution (11)

$$
\begin{aligned}
& x\left(t_{1}, t_{2}, \ldots, t_{q+1} ; k_{1}, \ldots, k_{r}\right)= \\
& =\Phi_{1}\left(t_{1}, t_{1}^{0} ; t_{2}, \ldots, t_{q+1} ; k_{1}, \ldots, k_{r}\right) \\
& \cdot x\left(t_{1}^{0} ; t_{2}, \ldots, t_{q+1} ; k_{1}, \ldots, k_{r}\right)+ \\
& +\int_{t_{1}^{0}}^{t_{1}} \Phi_{1}\left(t_{1}, s_{1} ; t_{2}, \ldots, t_{q+1} ; k_{1}, \ldots, k_{r}\right) \\
& \cdot z\left(s_{1} ; t_{2}, \ldots, t_{q+1} ; k_{1}, \ldots, k_{r}\right) d s_{1}
\end{aligned}
$$

and, by replacing $z$ (12) one obtains formula (5) for the case $(q+1, r)$.

Similarly one can derive the case $(q, r+1)$ from ( $q, r$ ), hence we proved by induction (5) for any $q, r \in$ N.

Theorem 1. The state of the system $\Sigma$ (1) determined by the initial state $x_{0} \in \mathbf{R}^{n}$ and the control $u$ is

$$
\begin{align*}
& x(t ; k)=\left(\prod_{i=1}^{q} \Phi_{i}\left(t_{i}, t_{i}^{0} ; t_{\overline{i-1}}^{0}, t_{\tilde{\tilde{i}}} ; k\right)\right) \\
& \left(\prod_{j=1}^{r} F_{j}\left(t^{0} ; k_{j}, k_{j}^{0} ; k_{\overline{j-1}}^{0}, k_{\tilde{j}}\right)\right) x^{0}+  \tag{13}\\
& +\int_{t_{1}^{0}}^{t_{1}} \ldots \int_{t_{q}^{0}}^{t_{q}} \sum_{l_{1}=k_{1}^{0}}^{k_{1}-1} \ldots \sum_{l_{r}=k_{r}^{0}}^{k_{r}-1}\left(\prod_{i=1}^{q} \Phi_{i}\left(t_{i}, s_{i} ; s_{\overline{i-1}}, t_{\tilde{i}} ; k\right)\right) \\
& \cdot\left(\prod_{j=1}^{r} F_{j}\left(s ; k_{j}, l_{j}+1 ; l_{\overline{j-1}}, k_{\tilde{\tilde{j}}}\right)\right) . \\
& \cdot B(s ; l) u(s ; l) d s_{1} \ldots d s_{q} .
\end{align*}
$$

Proof. Equation (1) has the form (4) with $f(t ; k)=B(u ; k) u(t ; k)$ and (13) results from (5) by replacing $f(t ; k)$.

By replacing the state $x(t ; k)$ (13) in the output equation (2) we obtain formula (14) below of the general response of the system.

Theorem 2. The input-output map of the ( $q, r$ )-D
$\operatorname{CDS} \Sigma(1),(2) i s$

$$
\begin{align*}
& y(t ; k)=C(t ; k)\left(\prod_{i=1}^{q} \Phi_{i}\left(t_{i}, t_{i}^{0} ; t_{\overline{i-1}}^{0}, t_{\tilde{i}}, k\right)\right) \\
& \left(\prod_{j=1}^{r} F_{j}\left(t^{0} ; k_{j} ; k_{\overline{j-1}}^{0}, k_{\tilde{j}}\right)\right) x^{0}+ \\
& +\int_{t_{1}^{0}}^{t_{1}} \ldots \int_{t_{q}^{0}}^{t_{q}} \sum_{l_{1}=k_{1}^{0}}^{k_{1}-1} \ldots \sum_{l_{r}=k_{r}^{0}}^{k_{r}-1} C(t ; k)  \tag{14}\\
& \cdot\left(\prod_{i=1}^{q} \Phi_{i}\left(t_{i}, s_{i} ; s_{\overline{i-1}}, t_{\tilde{i}} ; k\right)\right) \\
& \cdot\left(\prod_{j=1}^{r} F_{j}\left(s ; k_{j}, l_{j} ; l_{\overline{j-1}}, k_{\tilde{\tilde{j}}}\right)\right) \\
& \cdot B(s ; l) u(s ; l) d s_{1} \ldots d s_{q}+D(t ; k) u(t ; k)
\end{align*}
$$

## 3 Observability of time-varying ( $q, r$ )-D CDSs

For some $\left(t^{0} ; k^{0}\right),\left(t^{1} ; k^{1}\right),(t ; k) \in T$ with $\left(t^{0} ; k^{0}\right)<$ $\left(t^{1} ; k^{1}\right)$ we denote by $P$ the multiple interval $P=$ $\left[t^{0} ; t^{1}\right] \times\left[k^{0} ; k^{1}\right], \quad$ by $\int_{t^{0}}^{t}$ the multiple integral $\int_{t_{1}^{0}}^{t_{1}} \cdots \int_{t_{q}^{0}}^{t_{q}}$, by $\sum_{l=k^{0}}^{k-1}$ the sum $\sum_{l_{1}=k_{1}^{0}}^{k_{1}-1} \cdots \sum_{l_{r}=k_{r}^{0}}^{k_{r}-1}$ and $d s=$ $d s_{1} \cdots d s_{q}$.

A triplet $(t, k, \tilde{x}) \in \mathbf{R}^{q} \times \mathbf{Z}^{r} \times \mathbf{R}^{n}$ is said to be a phase of $\Sigma$ if $\exists u: T \rightarrow \mathbf{R}^{m}$ and $x^{0} \in \mathbf{R}^{n}$ such that $\tilde{x}=x(t ; k)$ where $x(t ; k)$ is given by (13). In this case one says that the control $u$ transfers the phase $\left(t^{0}, k^{0}, x^{0}\right)$ to the phase $(t, k, \tilde{x})$.

Definition 3. A phase $\left(t^{0}, k^{0}, x\right)$ is said to be unobservable/unobservable on $P$ if for any control $u(t ; k)$ it provides the same output $y(t ; k)$ as the phase $\left(t^{0}, k^{0}, 0\right)$ for any $(t ; k) \geq\left(t^{0} ; k^{0}\right) /$ for any $(t ; k) \in P$.

A state $x$ is said to be unobservable at $\left(t^{0} ; k^{0}\right)$ /unobservable on $P$ if the phase $\left(t^{0}, k^{0}, x\right)$ is unobservable/unobservable on $P$.

Proposition 2. The phase $\left(t^{0}, k^{0}, x\right)$ is unobservable/unobservable on $P$ if and only if

$$
\begin{align*}
& C(t ; k)\left(\prod_{i=1}^{q} \Phi_{i}\left(t_{i}, t_{i}^{0} ; t_{\overline{i-1}}^{0}, t_{\tilde{i}}, k\right)\right) \\
& \left(\prod_{j=1}^{r} F_{j}\left(t^{0} ; k_{j} ; k_{j-1}^{0}, k_{\tilde{j}}\right)\right) x^{0}=0 . \tag{15}
\end{align*}
$$

for any $(t ; k) \in T,(t ; k) \geq\left(t^{0} ; k^{0}\right)$ for any $(t ; k) \in$ $P$.

Proof. We replace $x^{0}$ by 0 in (14) and we obtain the output provided by the null initial state and by an arbitrary control $u$ :

$$
\begin{aligned}
& y_{0}(t ; k)=\int_{t_{1}^{0}}^{t_{1}} \cdots \int_{t_{q}^{0}}^{t_{q}} \sum_{l_{1}=k_{1}^{0}}^{k_{1}-1} \ldots \sum_{l_{r}=k_{r}^{0}}^{k_{r}-1} C(t ; k) \\
& \cdot\left(\prod_{i=1}^{q} \Phi_{i}\left(t_{i}, s_{i} ; s_{\overline{i-1}}, t_{\bar{i}} ; k\right)\right) \cdot \\
& \cdot\left(\prod_{j=1}^{r} F_{j}\left(s ; k_{j}, l_{j} ; l_{\overline{j-1}}, k_{\bar{j}}^{\bar{j}}\right)\right) B(s ; l) u(s ; l) d s_{1} \ldots d s_{q}+ \\
& +D(t ; k) u(t ; k)
\end{aligned}
$$

It results that the state $x^{0}$ is unobservable at $\left(t^{0} ; k^{0}\right)$ /unobservable on P if and only if $y_{0}(t ; k)=$ $y(t ; k)$ for any $(t ; k) \geq\left(t^{0} ; k^{0}\right) /$ for any $(t ; k) \in P$ (where the output $y(t ; k)$ is given by (14)). Obviously, this condition is equivalent to (15).

Formula (15) as well as all the following formulas concerning observability involve only the drift matrices $A_{c i}, A_{d j}$ and the matrix $C$, therefore in the sequel we shall consider $(q, r)$-D CDSs reduced to the form $\Sigma=\left(\left\{A_{c i} \mid i \in \bar{q}\right\},\left\{A_{d j} \mid j \in \bar{r}\right\}, C\right)$.

Definition 4. The system $\Sigma=\left(\left\{A_{c i} \mid i \in \bar{q}\right\}\right.$, $\left.\left\{A_{d j} \mid j \in \bar{r}\right\}, C\right)$ is said to be completely observable at $\left(t^{0} ; k^{0}\right)$ / completely observable on $P$ if there is no state $x \in \mathbf{R}^{n}, x \neq 0$ unobservable at $\left(t^{0} ; k^{0}\right)$ / unobservable on P .

Definition 5. The matrix

$$
\begin{align*}
& \mathcal{O}_{\Sigma}\left(t^{0}, t^{1} ; k^{0}, k^{1}\right)=\int_{t^{0}}^{t^{1}} \sum_{l=k^{0}}^{k^{1}-1} \\
& \cdot\left(C(t ; k)\left(\prod_{i=1}^{q} \Phi_{i}\left(t_{i}, t_{i}^{0} ; t_{\overline{i-1}}^{0}, t_{\tilde{\bar{i}}}, k\right)\right)\right. \\
& \left.\cdot\left(\prod_{j=1}^{r} F_{j}\left(t^{0} ; k_{j} ; k_{\overline{j-1}}^{0}, k_{\tilde{\bar{j}}}\right)\right)\right)^{T} \cdot  \tag{16}\\
& \cdot\left(C(t ; k)\left(\prod_{i=1}^{q} \Phi_{i}\left(t_{i}, t_{i}^{0} ; t_{\overline{i-1}}^{0}, t_{\tilde{\bar{i}}}, k\right)\right)\right. \\
& \left.\cdot\left(\prod_{j=1}^{r} F_{j}\left(t^{0} ; k_{j} ; k_{\frac{0}{j-1}}^{0}, k_{\tilde{j}}\right)\right)\right) d t
\end{align*}
$$

is called the observability Gramian of $\Sigma$ on $P$.
Obviously $\mathcal{O}_{\Sigma}=\mathcal{O}_{\Sigma}\left(t^{0}, t^{1} ; k^{0}, k^{1}\right)$ is a symmetrical positive semidefinite $n \times n$ matrix.

Proposition 3. The state $x$ is unobservable on $P$ if and only if

$$
\begin{equation*}
\mathcal{O}_{\Sigma}\left(t^{0}, t^{1} ; k^{0}, k^{1}\right) x=0 \tag{17}
\end{equation*}
$$

Proof. Necessity. If the state $x$ is unobservable on $P$ one obtains by (16) and by Proposition 2, for any $(t ; k) \in P$ :

$$
\begin{aligned}
& \mathcal{O}_{\Sigma}\left(t^{0}, t^{1} ; k^{0}, k^{1}\right) x= \\
& =\int_{t^{0}}^{t^{1}} \sum_{l=k^{0}}^{k^{1}-1}\left(C(t ; k)\left(\prod_{i=1}^{q} \Phi_{i}\left(t_{i}, t_{i}^{0} ; t_{\overline{i-1}}^{0}, t_{\tilde{i}}, k\right)\right)\right. \\
& \left.\cdot\left(\prod_{j=1}^{r} F_{j}\left(t^{0} ; k_{j} ; k_{\overline{j-1}}^{0}, k_{\tilde{j}}\right)\right)\right)^{T} \\
& \cdot\left(C(t ; k)\left(\prod_{i=1}^{q} \Phi_{i}\left(t_{i}, t_{i}^{0} ; t_{\overline{i-1}}^{0}, t_{\tilde{\bar{i}}}, k\right)\right)\right. \\
& \left.\cdot\left(\prod_{j=1}^{r} F_{j}\left(t^{0} ; k_{j} ; k_{j-1}^{0}, k_{\tilde{j}}\right)\right)\right) x d t=0
\end{aligned}
$$

Sufficiency. From (17) one obtains

$$
x^{T} \mathcal{O}_{\Sigma}\left(t^{0}, t^{1} ; k^{0}, k^{1}\right) x=0
$$

hence

$$
\begin{gathered}
\int_{t^{0}}^{t^{1}} \sum_{l=k^{0}}^{k^{1}-1} \| C(t ; k)\left(\prod_{i=1}^{q} \Phi_{i}\left(t_{i}, t_{i}^{0} ; t_{\overline{i-1}}^{0}, t_{\tilde{\bar{i}}}, k\right)\right) \\
\cdot\left(\prod_{j=1}^{r} F_{j}\left(t^{0} ; k_{j} ; k_{j-1}^{0}, k_{\tilde{j}}^{\tilde{j}}\right)\right) x \|^{2} d t=0
\end{gathered}
$$

Since the integrand is a sum of non-negative functions, it results that (17) holds a.e. on $P$, hence again by Proposition 2 the state $x$ is unobservable on $P$.

Similarly, we get
Proposition 4. The state $x$ is unobservable at $\left(t^{0} ; k^{0}\right)$ if and only if (17) holds for any $\left(t^{1} ; k^{1}\right) \in T$ with $\left(t^{0} ; k^{0}\right)<\left(t^{1} ; k^{1}\right)$.

Proposition 3 gives by paraphrase:
Corollary 1. The set of the states of the ( $q, r$ )-D CDS $\Sigma$ which are unobservable on $P$ is the subspace of $X$

$$
\begin{equation*}
X_{u o}=\operatorname{ker} \mathcal{O}_{\Sigma}\left(t^{0}, t^{1} ; k^{0}, k^{1}\right) \tag{18}
\end{equation*}
$$

By definition, the system $\Sigma$ is completely observable on $P$ if and only if $\operatorname{ker} \mathcal{O}_{\Sigma}\left(t^{0}, t^{1} ; k^{0}, k^{1}\right)=$ $\{0\}$, condition which is equivalent to the fact that $\mathcal{O}_{\Sigma}\left(t^{0}, t^{1} ; k^{0}, k^{1}\right)$ has the full rank $n$. We proved:

Theorem 3. The system $\Sigma$ is completely observable on $P$ if and only if

$$
\begin{equation*}
\operatorname{rank} \mathcal{O}_{\Sigma}\left(t^{0}, t^{1} ; k^{0}, k^{1}\right)=n \tag{19}
\end{equation*}
$$

The concept of observability is connected to the property of $\Sigma$ to allow the determining of the state of the system from the information about the exterior
signals (the control $u$ and the output $y$ ). The solution of this problem is given by the following result.

Theorem 4. Assume that the system $\Sigma=$ $\left(\left\{A_{c i} \mid i \in \bar{q}\right\},\left\{A_{d j} \mid j \in \bar{r}\right\}, C\right)$ is complete observable on $P$. If the control $u(t ; k)$ produces the output $y(t ; k), \forall(t ; k) \in P$, then the initial state $x^{0}$ is given by the formula

$$
\begin{align*}
& x^{0}=\mathcal{O}_{\Sigma}\left(t^{0}, t^{1} ; k^{0}, k^{1}\right)^{-1} \int_{t^{0}}^{t^{1}} \cdot \\
& \cdot \sum_{l=k^{0}}^{k^{1}-1}\left(C(t ; k)\left(\prod_{i=1}^{q} \Phi_{i}\left(t_{i}, t_{i}^{0} ; t_{\overline{i-1}}^{0}, t_{\tilde{i}}, k\right)\right) \cdot\right.  \tag{20}\\
& \left.\cdot\left(\prod_{j=1}^{r} F_{j}\left(t^{0} ; k_{j} ; k_{j-1}^{0}, k_{\tilde{j}}\right)\right)\right)^{T} \tilde{y}(t ; k) d t .
\end{align*}
$$

where

$$
\begin{aligned}
& \tilde{y}(t ; k)=y(t ; k)- \\
& +\int_{t_{1}^{0}}^{t_{1}} \ldots \int_{t_{q}^{0}}^{t_{q}} \sum_{l_{1}=k_{1}^{0}}^{k_{1}-1} \ldots \sum_{l_{r}=k_{r}^{0}}^{k_{r}-1} C(t ; k) . \\
& \cdot\left(\prod_{i=1}^{q} \Phi_{i}\left(t_{i}, s_{i} ; s \overline{i-1}, t_{\tilde{i}} ; k\right)\right) \cdot \\
& \cdot\left(\prod_{j=1}^{r} F_{j}\left(s ; k_{j}, l_{j} ; l \overline{j-1}, k_{\tilde{j}}\right)\right) B(s ; l) u(s ; l) d s_{1} \ldots d s_{q}- \\
& -D(t ; k) u(t ; k) .
\end{aligned}
$$

Remark 2. If we replace the initial state $x^{0}$ given by (20) and (21) in (13) we obtain the state $x(t ; k)$ for any $(t ; k) \in P$; therefore if a system $\Sigma$ is completely observable one can recover the whole trajectory $x(t ; k),(t ; k) \in P$ of the system from the exterior data.

## 4 Observability of time-invariant ( $q, r$ )-D CDSs

The system $\Sigma=\left(\left\{A_{c i} \mid i \in \bar{q}\right\},\left\{A_{d j} \mid j \in \bar{r}\right\} ; \mathbf{C}\right)$ is said to be time-invariant (or stationary) if all its matrices are constant. Since in this case the initial moment $\left(t^{0}, k^{0}\right)$ is not relevant, we can consider it $\left(t^{0}, k^{0}\right)=(0,0) \in T_{+}:=\mathbf{R}_{+}^{q} \times \mathbf{Z}_{+}^{r}$. The fundamental matrices become $\Phi_{i}\left(t_{i}, 0 ; t_{\tilde{i}} ; h\right)=e^{A_{c i} t_{i}}$ and $\left(F_{j}\left(t ; h_{j}, 0 ; h_{\tilde{j}}\right)=A_{d j}^{k_{j}}\right.$ and the input-output map (14)
can be written in the form

$$
\begin{gather*}
y(t, k)=C\left(\exp \left(\sum_{i=1}^{q} A_{c i} t_{i}\right)\right)\left(\prod_{j=1}^{r} A_{d j}^{k_{j}}\right) x^{0}+ \\
+\int_{0}^{t} \sum_{l=j}^{k-1} C\left(\exp \left(\sum_{i=1}^{q} A_{c i}\left(t_{i}-s_{i}\right)\right)\right)  \tag{22}\\
\cdot\left(\prod_{j=1}^{r} A_{d j}^{k_{j}-l_{j}-1}\right) B u(s, l) d s_{1} \ldots d s_{q}
\end{gather*}
$$

In this case, the observability Gramian (which will be denoted $\mathcal{O}_{\Sigma}(t, k)$ instead of $\mathcal{O}_{\Sigma}\left(t^{0}, t ; k^{0}, k\right)$ becomes

$$
\begin{gather*}
\mathcal{O}_{\Sigma}(t, k)=\int_{0}^{t} \sum_{l=0}^{k}\left(\prod_{j=1}^{r}\left(A_{d j}^{T}\right)^{l_{j}}\right) \\
\cdot\left(\exp \left(\sum_{i=1}^{q} A_{c i}^{T} s_{i}\right)\right) C^{T} C  \tag{23}\\
\cdot\left(\exp \left(\sum_{i=1}^{q} A_{c i} s_{i}\right)\right)\left(\prod_{j=1}^{r} A_{d j}^{l_{j}}\right) d s_{1} \ldots d s_{q}
\end{gather*}
$$

If the state $x$ is unobservable at $\left(t^{0}, k^{0}\right)=(0,0)$ we will say that $x$ is unobservable. From Proposition 2 and Theorem 3 we obtain

Proposition 6. The state $x \in X$ is unobservable if and only if

$$
\begin{gather*}
C\left(\exp \left(\sum_{i=1}^{q} A_{c i} t_{i}\right)\right)\left(\prod_{j=1}^{r} A_{d j}^{k_{j}}\right) x=0  \tag{24}\\
\forall(t, k) \in T_{+}
\end{gather*}
$$

Theorem 5. The system $\Sigma=\left(\left\{A_{c i} \mid i \in\right.\right.$ $\left.\bar{q}\},\left\{A_{d j} \mid j \in \bar{r}\right\}, C\right)$ is completely observable if and only if rank $\mathcal{O}_{\Sigma}(t, k)=n$ for any $(t, k) \in T_{+}$.

In this case of time-invariant systems we can use a simpler controllability matrix instead of the controllability Gramian.

Definition 6. The observability matrix of the system $\Sigma$ is

$$
\begin{gather*}
O_{\Sigma}=\left[C^{T} A_{c 1}^{T} C^{T} \ldots\left(A_{c 1}^{T}\right)^{n-1} C^{T} \ldots\right. \\
\ldots\left(\prod_{\alpha \in \gamma}\left(A_{c \alpha}^{T}\right)^{i_{\alpha}}\right)\left(\prod_{\beta \in \delta}\left(A_{d \beta}^{T}\right)^{j_{\beta}}\right) C^{T} \ldots  \tag{25}\\
\left.\ldots\left(\prod_{\alpha=1}^{q}\left(A_{c \alpha}^{T}\right)^{n-1}\right)\left(\prod_{\beta=1}^{r}\left(A_{d \beta}^{T}\right)^{n-1}\right) C^{T}\right]^{T}
\end{gather*}
$$

where we consider subsets $\gamma \subset \bar{q}$ and $\delta \subset \bar{r}$ and numbers $i_{\alpha}$ and $j_{\beta}$ verifying $0 \leq i_{\alpha}, j_{\beta} \leq n-1, \forall \alpha \in \gamma$ and $\forall \beta \in \delta ; \prod_{\alpha \in \emptyset} A_{\alpha}=I$.

Theorem 6. The system $\Sigma=\left(\left\{A_{c i} \mid i \in\right.\right.$ $\left.\bar{q}\},\left\{A_{d j} \mid j \in \bar{r}\right\}, C\right)$ is completely observable if and only if

$$
\begin{equation*}
\operatorname{rank} O_{\Sigma}=n \tag{26}
\end{equation*}
$$

Proof. Necessity. Let us assume that (26) fails, i.e. $\quad \operatorname{rank} \mathcal{O}_{\Sigma}<n$ since $\mathcal{O}_{\Sigma}$ has $n$ columns. Then there exists $x \in X \backslash\{0\}$ such that $\mathcal{O}_{\Sigma} x=0$, which implies $C x=0$, $C A_{c 1} x=0, \ldots, C\left(\prod_{\beta \in \delta} A_{d \beta}^{j_{\beta}}\right)\left(\prod_{\alpha \in \gamma} A_{c \alpha}^{i_{\alpha}}\right) x=$ $0, \ldots, C\left(\prod_{\beta=1}^{r} A_{d \beta}^{n-1}\right)\left(\prod_{\alpha=1}^{q} A_{c \alpha}^{n-1}\right) x=0$.

By the Hamilton-Cayley Theorem applied to the matrices $A_{c i}, i \in \bar{q}$ and $A_{d j}, j \in \bar{r}$ ( and taking into account the commutativity of these matrices), we can prove that

$$
\begin{equation*}
C\left(\prod_{i=1}^{q} A_{c i}^{a_{i}}\right)\left(\prod_{j=1}^{r} A_{d j}^{k_{j}}\right) x=0, \forall a_{i}, k_{j} \in \mathbf{N} . \tag{27}
\end{equation*}
$$

Then

$$
\begin{aligned}
& C\left(\exp \left(\sum_{i=1}^{q} A_{c i} t_{i}\right)\right)\left(\prod_{j=1}^{r} A_{d j}^{k_{j}}\right) x= \\
& =C\left(\prod_{i=1}^{q}\left(\sum_{a_{i}=0}^{\infty} \frac{A_{c i}^{a_{i}} t_{i}^{a_{i}}}{a_{i}!}\right)\right)\left(\prod_{j=1}^{r} A_{d j}^{k_{j}}\right) x= \\
& =\sum_{a_{i}=0}^{\infty}\left(\prod_{i=1}^{q} \frac{t_{i}^{a_{i}}}{a_{i}!}\right) C\left(\prod_{i=1}^{q} A_{c i}^{a_{i}}\right)\left(\prod_{j=1}^{r} A_{d j}^{k_{j}}\right) x=0 .
\end{aligned}
$$

It results by Proposition 6 that the state $x$ is unobservable, hence $\Sigma$ is not completely observable.

Sufficiency. Assume that $\Sigma$ is not completely observable, hence by Proposition 6 equality (24) holds for some $x \in X \backslash\{0\}$. By deriving successively this equality with respect to $t_{i}, i \in \bar{q}$ and by taking $t=\left(t_{1}, \ldots, t_{q}\right)=(0, \ldots, 0)$ we obtain (27). Then $\mathcal{O}_{\Sigma} x=0$, hence $\operatorname{rank} \mathcal{O}_{\Sigma}<n$.

From the proof of Theorem 6 we also deduce
Corollary 2. A state $x \in X$ is unobservable if and only if $x$ verifies (27).

Corollary 3. The set of all unobservable states of $\Sigma$ is $X_{u o}=\operatorname{Ker} O_{\Sigma}$.

Theorem 7. The system $\Sigma$ is completely observable if and only if equality (27) implies $x=0$.

The following theorem emphasizes the duality relation between the concepts of observability and reachability.

Definition 7. The system $\hat{\Sigma}=\left(\left\{\hat{A}_{c i} \mid i \in\right.\right.$ $\left.\bar{q}\},\left\{\hat{A}_{d j} \mid j \in \bar{r}\right\}, \hat{B}, \hat{C}, \hat{D}\right)$ is called the dual of the
system $\Sigma=\left(\left\{A_{c i} \mid i \in \bar{q}\right\},\left\{A_{d j} \mid j \in \bar{r}\right\}, B, C, D\right)$ if $\forall i \in \bar{q}, \hat{A}_{c i}=A_{c i}^{T}, \forall j \in \bar{r}, \hat{A}_{d j}=A_{d j}^{T}, \hat{B}=C^{T}$, $\hat{C}=B^{T}, \hat{D}=D^{T}$.

Theorem 8. The system $\Sigma$ is completely observable if and only if its dual $\hat{\Sigma}$ is completely controllable.

Proof. Obviously, the observability matrix $\mathcal{O}_{\Sigma}$ (18) of $\Sigma$ and the controllability matrix of $\hat{\Sigma}, C_{\hat{\Sigma}}$ (see Definition 4 in [17]) verify $\mathcal{O}_{\Sigma}^{T}=C_{\hat{\Sigma}}$, hence rank $\mathcal{O}_{\Sigma}=\operatorname{rank} C_{\hat{\Sigma}}$. The proof is complete by Theorem 6 (and by Theorem 4.3 in [14], which states that $\hat{\Sigma}$ is completely reachable iff rank $C_{\hat{\Sigma}}=n$ ).

Now we can give a geometric characterization of the set of unobservable states of $\Sigma$.

Theorem 9. The set $X_{u}$ of all unobservable states of $\Sigma$ is the greatest subspace of $X$ which is $\left(\left\{A_{c i} \mid i \in \bar{q}\right\},\left\{A_{d j} \mid j \in \bar{r}\right\}\right)$ - invariant and it is contained in $\operatorname{Ker} C$.

Proof. By Corollary 2,

$$
\begin{gathered}
X_{u o}=\left\{x \in X \mid C\left(\prod_{i=1}^{q} A_{c i}^{a_{i}}\right)\left(\prod_{j=1}^{r} A_{d j}^{k_{j}}\right) x=0\right. \\
\left.\forall a_{i}, k_{j} \in \mathbf{N}\right\}
\end{gathered}
$$

Obviously, for $a_{i}=0, \forall i \in \bar{q}$ and $k_{j}=0, \forall j \in \bar{r}$ it results than $C x=0 \forall x \in X_{u o}$, hence $X_{u o} \subset \operatorname{Ker} C$. For any $\iota \in \bar{q}, x \in X_{u o}$ implies

$$
\begin{aligned}
& C\left(\prod_{i=1}^{q} A_{c i}^{a_{i}}\right)\left(\prod_{j=1}^{r} A_{d j}^{k_{j}}\right)\left(A_{c \iota} x\right)= \\
& =C\left(\prod_{i=1}^{q} A_{c i}^{b_{i}}\right)\left(\prod_{j=1}^{r} A_{d j}^{k_{j}}\right) x=0
\end{aligned}
$$

where $b_{i}=\left\{\begin{array}{ll}a_{i}+1 & \text { for } i=\iota \\ a_{i} & \text { for } i=\tilde{\iota}\end{array}\right.$, hence $A_{c \iota} x \in$ $X_{u o}$.

Similarly, we can prove that $A_{d j} x \in X_{u o}, \forall j \in \bar{r}$, hence $X_{u o}$ is invariant with respect to $A_{c i}, \forall i \in \bar{q}$ and $A_{d j}, \forall j \in \bar{r}$. Now let us assume that $\mathcal{V}$ is a subspace of $X\left(\left\{A_{c i} \mid i \in \bar{q}\right\},\left\{A_{d j} \mid j \in \bar{r}\right\}\right)$-invariant included in $\operatorname{Ker} C$. Let $x$ be any element of $\mathcal{V}$. Since $\mathcal{V} \subset \operatorname{Ker} C, C x=0 . \mathcal{V}$ being invariant with respect to the drift matrices we obtain $A_{c i} x \in \mathcal{V}$ and $A_{d j} \in \mathcal{V}, \forall i \in \bar{q}, \forall j \in \bar{r}$ and by recurrence, $\left(\prod_{i=1}^{q} A_{c i}^{a_{i}}\right)\left(\prod_{j=1}^{r} A_{d j}^{k_{j}}\right) x \in \mathcal{V} \subset \operatorname{Ker} C, \forall a_{i}, k_{j} \in \mathbf{N}$, hence $C\left(\prod_{i=1}^{q} A_{c i}^{a_{i}}\right)\left(\prod_{j=1}^{r} A_{d j}^{k_{j}}\right) x=0$, i.e. $x \in X_{u o}$. Therefore $\mathcal{V} \subset X_{u o}$.

From Theorem 5 and Theorem 9 we obtain
Theorem 10. The system $\Sigma$ is completely observable if $\{0\}$ is the greatest subspace of $X$ which is $\left(\left\{A_{c i} \mid i \in \bar{q}\right\},\left\{A_{d j} \mid j \in \bar{r}\right\}\right)$-invariant and included in Ker $C$.

We denote by $\sigma(A)$ the spectrum of a matrix $A$.
Theorem 11. The system $\Sigma$ is completely observable if and only if there is no vector $x \in \mathbf{R}^{n} \backslash\{0\}$ such that

$$
\begin{gather*}
C\left(\prod_{i=1}^{q}\left(s_{i} I-A_{c i}\right)^{-1}\right) \\
\left(\prod_{j=1}^{r}\left(z_{j} I-A_{d j}\right)^{-1}\right) x=0 \tag{28}
\end{gather*}
$$

$\forall s_{i} \in \mathbf{C} \backslash \sigma\left(A_{c i}\right), i \in \bar{q}, \forall z_{j} \in \mathbf{C} \backslash \sigma\left(A_{d j}\right), j \in \bar{r}$.
Proof. Since for a square matrix $A$

$$
(z I-A)^{-1}=\sum_{k=0}^{\infty} A^{k} z^{-k-1}, \forall z \in \mathbf{C} \backslash \sigma(A)
$$

equality (28) is equivalent to

$$
\begin{gathered}
\sum_{a \geq 0} \sum_{k \geq 0}\left(C\left(\prod_{i=1}^{q} A_{c i}^{a_{i}}\right)\left(\prod_{j=1}^{r} A_{d j}^{k_{j}}\right) x\right) \\
\left(\prod_{i=1}^{q} s_{i}^{-a_{i}-1}\right)\left(\prod_{j=1}^{r} z_{j}^{-k_{j}-1}\right)=0 .
\end{gathered}
$$

This multiple Laurent series is null if and only if all its coefficients are equal to zero, condition which is equivalent to (27) hence Theorem 11 is a consequence of Theorem 7.

Definition 8. Two systems $\Sigma=\left(\left\{A_{c i} \mid i \in\right.\right.$ $\left.\bar{q}\},\left\{A_{\tilde{A}} \mid j \in \bar{r}\right\}, \underset{\sim}{B}, \underset{\sim}{C}, \underset{\sim}{D}\right)$ and $\tilde{\Sigma}=\left(\left\{\tilde{A}_{c i} \mid i \in\right.\right.$ $\left.\bar{q}\},\left\{\tilde{A}_{d j} \mid j \in \bar{r}\right\}, \tilde{B}, \tilde{C}, \tilde{D}\right)$ are isomorphic if there exists a nonsingular matrix $U$ such that $\tilde{A}_{c i}=$ $U^{-1} A_{c i} U, \forall i \in \bar{q}, \tilde{A}_{d j}=U^{-1} A_{d j} U, \forall j \in \bar{r}$, $\tilde{B}=U^{-1} B, \tilde{C}=C U, \tilde{D}=D$.

The next result gives the canonical form of the unobservable systems.

Theorem 12. The system $\Sigma$ is not completely observable if and only if it is isomorphic to a system $\tilde{\Sigma}$ with

$$
\begin{align*}
& \tilde{A}_{c i}=\left[\begin{array}{cc}
A_{11, c i} & 0 \\
A_{21, c i} & A_{22, c i}
\end{array}\right], \forall i \in \bar{q},  \tag{29}\\
& \tilde{A}_{d j}=\left[\begin{array}{cc}
A_{11, d j} & 0 \\
A_{21, d j} & A_{22, d j}
\end{array}\right], \forall j \in \bar{r}, \tilde{C}=\left[C_{1} 0\right],
\end{align*}
$$

where $A_{11, c i}, A_{11, d j} \in \mathbf{R}^{\tilde{n} \times \tilde{n}}$ and $\tilde{n}<n$. For $\tilde{n}=\operatorname{rank} \mathcal{O}_{\Sigma}$, the subsystem $\Sigma_{1}=\left(\left\{A_{11, c i} \mid i \in\right.\right.$ $\left.\bar{q}\},\left\{A_{11, d j} \mid j \in \bar{r}\right\}, C_{1}\right)$ is completely observable.

Proof. Let $\tilde{n}=\operatorname{rank} \mathcal{O}_{\Sigma}$. Since $\Sigma$ is not completely observable $\operatorname{rank} \mathcal{O}_{\Sigma}<n$, hence $\tilde{n}<n$ and by Corollary $3, \tilde{n}=\operatorname{dim} X_{u o}$. If we consider a basis of $X_{u o}$ and we complete it to a basis $\mathcal{B}$ of $X=\mathbf{R}^{n}$, we obtain the direct sum decomposition $X=X_{1} \oplus X_{2}$, where $X_{2}=X_{u o}$. Let $U$ be the transition matrix. Let us denote by $\tilde{A}_{c i}, \tilde{A}_{d j}, \tilde{C}$ the matrices corresponding to $A_{c i}, A_{d j}$ and $C$ in the basis $B$. A state $x$ becomes $\tilde{x}=U^{-1} x$ and since $X_{u o}=X_{2}$ we obtain

$$
X_{2}=\left\{\left.\tilde{x}=\left[\begin{array}{c}
0 \\
x_{2}
\end{array}\right] \right\rvert\, x_{2} \in \mathbf{R}^{n-\tilde{n}}\right\} .
$$

Let us partition the matrices corresponding to the direct sum decomposition, i.e. the matrices of the isomorphic system $\tilde{\Sigma}$ as

$$
\begin{aligned}
& \tilde{A}_{c i}=\left[\begin{array}{ll}
A_{11, c i} & A_{12, c i} \\
A_{21, c i} & A_{22, c i}
\end{array}\right], \\
& \tilde{A}_{d j}=\left[\begin{array}{ll}
A_{11, d j} & A_{12, d j} \\
A_{21, d j} & A_{22, d j}
\end{array}\right], \tilde{C}=\left[C_{1} C_{2}\right],
\end{aligned}
$$

hence $A_{11, c i}, A_{11, d j}$ are $\tilde{n} \times \tilde{n}$ matrices and $C_{1}$ is a $p \times \tilde{n}$ matrix. Since by Theorem $9 X_{2}$ is $A_{c i}$-invariant, we have $\tilde{A}_{c i} \tilde{x} \in X_{2}$, hence

$$
\begin{aligned}
& {\left[\begin{array}{ll}
A_{11, c i} & A_{12, c i} \\
A_{21, c i} & A_{22, c i}
\end{array}\right]\left[\begin{array}{c}
0 \\
x_{2}
\end{array}\right]=} \\
& =\left[\begin{array}{ll}
A_{12, c i} & x_{2} \\
A_{22, c i} & x_{2}
\end{array}\right] \in X_{2}, \forall x_{2} \in \mathbf{R}^{n-\tilde{n}} .
\end{aligned}
$$

Then $A_{12, c i} x_{2}=0, \forall x_{2} \in \mathbf{R}^{n-\tilde{n}}$, hence $A_{12, c i}=0$, $\forall i \in \bar{q}$. Similarly, we can prove that $A_{12, d j}=0, \forall j \in$ $\bar{r}$. Since again by Theorem $9 X_{2}$ is included in $\operatorname{Ker} C$, it results that (in the basis $\mathcal{B}$ ) $\left[C_{1} C_{2}\right]\left[\begin{array}{c}0 \\ x_{2}\end{array}\right]=0$, $\forall x_{2} \in \mathbf{R}^{n-n_{2}}$, hence $C_{2}=0$ and the matrices of the system have the form (29).

Now, by (29) we can prove that

$$
\begin{gathered}
\tilde{C}\left(\prod_{i=1}^{q} \tilde{A}_{c i}\right)\left(\prod_{j=1}^{r} \tilde{A}_{d j}\right)= \\
=\left[C_{1}\left(\prod_{i=1}^{q} \tilde{A}_{11, c i}\right)\left(\prod_{j=1}^{r} \tilde{A}_{11, d j}\right) \quad 0\right]
\end{gathered}
$$

where 0 is a $p \times(n-\tilde{n})$ null matrix. Then the controllability matrix of $\tilde{\Sigma}$ is formed by blocks having this structure.

By Hamilton-Cayley Theorem applied to matrices $A_{11, c i}$ and $A_{11, d j}$ it results that $A_{11, c i}^{\tilde{n}+k}$ are linear combination of matrices $A_{11, c i}^{l}, l=0, \tilde{n}-1$, for any $k \geq 0$ and similarly for $A_{11, d j}^{\tilde{\tilde{n}}+k}$. Then, if $\mathcal{O}_{\Sigma_{1}}$ is the observability matrix of the subsystem $\Sigma_{1}$, since $\mathcal{O}_{\tilde{\Sigma}}=\mathcal{O}_{\Sigma} U$ and $U$ is nonsingular, we get

$$
\operatorname{dim} \Sigma_{1}=\tilde{n}=\operatorname{rank} \mathcal{O}_{\Sigma}=\operatorname{rank} \mathcal{O}_{\tilde{\Sigma}}=\operatorname{rank} \mathcal{O}_{\Sigma_{1}}
$$

hence $\Sigma_{1}$ is completely observable.
Conversely if $\Sigma$ is isomorphic to a system $\tilde{\Sigma}$ (29), since the blocks of $\mathcal{O}_{\tilde{\Sigma}}$ have the above structure, obviously rank $\mathcal{O}_{\Sigma}=\operatorname{rank} \mathcal{O}_{\tilde{\Sigma}} \leq \tilde{n}<n$, hence $\Sigma$ is not completely observable.

We can restate this theorem as
Theorem 13. The system $\Sigma$ is completely observable if and only if it is not isomorphic to a system of the form (29).

Theorem 14. The system $\Sigma$ is completely observable if and only if there is no common eigenvector of the matrices $A_{c i}, i \in \bar{q}$ and $A_{d j}, j \in \bar{r}$, belonging to $\operatorname{Ker} C$.

Proof. Necessity. Let us assume that $\exists x \in \mathbf{R}^{n} \backslash$ $\{0\}$ such that $A_{c i} x=\lambda_{i} x, A_{d j} x=\mu_{j} x$ for some $\lambda_{i}, \mu_{j} \in \mathbf{C}, \forall i \in \bar{q}, \forall j \in \bar{r}$ and $C x=0$. Then we can prove by induction that

$$
\begin{aligned}
& C\left(\prod_{i=1}^{q} \tilde{A}_{c i}^{a_{i}}\right)\left(\prod_{j=1}^{r} \tilde{A}_{d j}^{k_{j}}\right) x= \\
= & \left(\prod_{i=1}^{q} \lambda_{i}^{a_{i}}\right)\left(\prod_{j=1}^{r} \mu_{j}^{k_{j}}\right) C x=0,
\end{aligned}
$$

$\forall a_{i}, k_{j} \geq 0, i \in \bar{q}, i \in \bar{r}$, therefore $\mathcal{O}_{\Sigma} x=0$ which implies that rank $\mathcal{O}_{\Sigma}<n$ and $\Sigma$ is not completely observable.

Sufficiency. Let us assume that $\Sigma$ is not completely observable. Then there exists $v \in \mathbf{R}^{n} \backslash\{0\}$ such that $\mathcal{O}_{\Sigma} v=0$. Let us denote by $S_{1}$ the subspace $X_{u o}=\operatorname{Ker} \mathcal{O}_{\Sigma}$, hence $S_{1}$ is a proper subspace of $\mathbf{R}^{n}$, and by Theorem $9 S_{1}$ is $\left(\left\{A_{c i} \mid i \in \bar{q}\right\},\left\{A_{d j} \mid j \in \bar{r}\right\}\right)$ invariant and it is contained in $\operatorname{Ker} C$. Being $A_{c 1^{-}}$ invariant, $S_{1}$ contains an eigenvector $v_{1}$ of $A_{c 1}$, corresponding to some eigenvalue $\lambda_{1}$. Then the set $S_{2}=$ $\left\{x \in \mathbf{C}^{n} \mid A_{c 1} x=\lambda_{1} x\right\}$ is a proper subspace of $\mathbf{C}^{n}$ as well as $S_{3}=S_{1} \cap S_{2}$. Following the ideas in the proof of Theorem 9 we can show that $S_{2}$ (and $S_{3}$ too) is $\left(\left\{A_{c i} \mid i \in \bar{q}\right\},\left\{A_{d j} \mid j \in \bar{r}\right\}\right)$-invariant. Then $S_{3}$ contains an eigenvector $v_{2}$ of $A_{c 2}$ corresponding to an eigenvalue $\lambda_{2}$, and since $S_{3} \subset S_{2}, v_{2}$ is an eigenvector of $A_{c 1}$ too. Now we consider the proper subspaces $S_{4}=\left\{x \in \mathbf{C}^{n} \mid A_{c 1} x=\lambda_{1} x, A_{c 2} x=\lambda_{2} x\right\}$ and $S_{5}=S_{3} \cap S_{4}$ and so on. Finally we obtain a sequence of proper subspaces of $\mathbf{C}^{n}, X_{u o}=S_{1} \supset$
$S_{3} \supset S_{5} \supset \ldots \supset S_{2 k-1} \supset \ldots \supset S_{2(q+r)-1}$ when $S_{2(q+r)-1}=\left\{x \in \mathbf{C}^{n} \mid A_{c i} x=\lambda_{i} x, i \in \bar{q}, A_{d j} x=\right.$ $\left.\mu_{j} x, j \in \bar{r}\right\} \neq\{0\}$ and $S_{2(q+r)-1} \subset X_{u o} \subset \operatorname{Ker} C$ hence there exists a common eigenvector $x$ of the drift matrices of $\Sigma$ which belongs to $\operatorname{Ker} C$.

We can derive from Theorem 14 a Belevitch-Hautus-Popov type criterion of observability

Theorem 15. The system $\Sigma$ is completely observable if and only if for any $\lambda_{i}, \mu_{j} \in \mathbf{C}, i \in \bar{q}, j \in \bar{r}$

$$
\begin{align*}
& \operatorname{rank}\left[C^{T} \lambda_{1} I-A_{c 1}^{T} \ldots \lambda_{q} I-A_{c q}^{T}\right. \\
& \left.\mu_{1} I-A_{d 1}^{T} \ldots \mu_{r} I-A_{d r}^{T}\right]^{T}=n \tag{30}
\end{align*}
$$

Proof. Let us denote the matrix in (30) by $M$. The existence of $\lambda_{i}, \mu_{j} \in \mathbf{C}, i \in \bar{q}, j \in \bar{r}$ such that $\operatorname{rank} M<$ $n$ is equivalent to the existence of $v \in \mathbf{R}^{n} \backslash\{0\}$ such that $M v=0$, i.e. such that $C v=0, A_{c i} v=\lambda_{i} v$, $i \in \bar{q}, A_{d j} v=\mu_{j} v, j \in \bar{r}$ which is equivalent by Theorem 14 to the fact that $\Sigma$ is not completely observable.

Corollary 4. The system $\Sigma$ is completely observable if and only if (30) holds for any $\lambda_{i} \in \sigma\left(A_{c i}\right)$, $i \in \bar{q}$ and $\mu_{j} \in \sigma\left(A_{d j}\right), j \in \bar{r}$.

Proof. This statement is an imediate consequence of Theorem 15, since for any square matrix $A, \lambda \in$ $\sigma(A)$ if and only if $\operatorname{det}(\lambda I-A)=0$. Then $\operatorname{rank}(\lambda I-$ $A)=n$ if and only if $\lambda \in \mathbf{C} \backslash \sigma(A)$, hence obviously (30) holds for $\lambda_{i} \in \mathbf{C} \backslash \sigma\left(A_{c i}\right)$ and $\mu_{1} \in \mathbf{C} \backslash \sigma\left(A_{d j}\right)$, $i \in \bar{q}, j \in \bar{r}$.

Conclusion. This paper studies a class of multidimensional hybrid linear systems from the point of view of observability. In the case of time-varying systems necessary and sufficient conditions are expressed by introducing a suitable observability Gramian. The connection between the concepts of observability and reachability is emphasized. For time-invariant systems ten criteria of observability are obtained. The geometric characterization of the subspace of unobservable states is given. This study can be continued for other concepts such as stability, stabilizability and detectability of multidimensional hybrid systems and it can be applied to the problems of minimal realizations.

Aknowledgement. This paper is partially supported by the Grant CNCSIS 86/ 2007 and by the 15-th Italian-Romanian Executive Programme of S\&T Co-operation for 2006-2008, University Politehnica of Bucharest.

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