A characterization of switched linear control systems with finite $L_2$-gain

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Abstract—Motivated by an open problem posed by J.P. Hespanha, we extend the notion of Barabanov norm and extremal trajectory to classes of switching signals that are not closed under concatenation. We use these tools to prove that the finiteness of the $L_2$-gain is equivalent, for a large set of switched linear control systems, to the condition that the generalized spectral radius associated with any minimal realization of the original switched system is smaller than one.

Index Terms—switched systems, $L_2$-gain

I. INTRODUCTION

Let $n, m, p$ be positive integers and $\tau$ be a positive real number. Consider the switched linear control system

$$\dot{x} = A_\sigma x + B_\sigma u, \quad y = C_\sigma x + D_\sigma u,$$  \hspace{1cm} (1)

where $x \in \mathbb{R}^n$, $u \in \mathbb{R}^m$, $y \in \mathbb{R}^p$ and $\sigma$ is in the class $\Sigma_\tau$ of piecewise constant signals with dwell time $\tau$ taking values in a fixed finite set $P$ of indices. Define the $L_2$-gain as

$$\gamma_2(\tau) = \sup \left\{ \|y_{u,\sigma}\|_2 \mid u \in L_2((0, \infty), \mathbb{R}^m), \{0\}, \sigma \in \Sigma_\tau \right\},$$

where $y_{u,\sigma}$ is the output corresponding to the trajectory of the system associated with $u$ and $\sigma$ starting at the origin at time $t = 0$. In [1, Problem 4.1], J.P. Hespanha asked the following questions: (i) under which conditions is the function $\tau \mapsto \gamma_2(\tau)$ bounded over $[0, \infty)$? (ii) when $\gamma_2$ is not a bounded function over $[0, \infty)$, how to compute $\tau_{\text{min}}$, the infimum of the dwell-times $\tau > 0$ for which $\gamma_2(\tau)$ is finite? (iii) how regular is $\gamma_2$?

In [2], Hespanha proved the noteworthy fact that, in general, $\lim_{\tau \to \infty} \gamma_2(\tau) = \max_{\sigma \in \Sigma_\tau} \gamma_2^P(\tau)$, where $\gamma_2^P(\tau)$ denotes the $L_2$-gain of the time-invariant control system where $\sigma(\cdot) \equiv p$. These results have been improved in [3] assuming that any operating mode satisfies a minimal realization assumption. In particular, $\gamma_2(\tau)$ is characterized in terms of a suitable switched Riccati equation, in the spirit of $H_\infty$-theory and worst-case switching laws (see also [4] for numerical investigations associated with this issue). Similar questions have been considered in [5] for other classes of switching signals (average dwell-time, persistent-dwell time...). In the latter work a restrictive hypothesis for a class of switching signals has been put forward, namely, that of closure under concatenation.

In the unswitched case $#P = 1$ the finiteness problem for the $L_2$-gain admits a classical solution, namely $\gamma_2$ is finite if and only if $A$ is Hurwitz, where $(A, B, C, D)$ is any minimal realization of $(A, B, C, D)$. A typical proof of the above fact involves a “worst” trajectory of the uncontrolled system $\dot{x} = \dot{\hat{x}}$, i.e., a trajectory whose norm behaves as $e^{\alpha t}$ where $\alpha$ is the maximum of the real parts of the eigenvalues of $\hat{A}$. As a consequence, one deduces that $\gamma_2$ is finite only if $\alpha < 0$. The argument can be adapted in the case of arbitrary switching and under a suitable observability assumption by replacing $\alpha$ with the maximal Lyapunov exponent of $\dot{x} = \dot{\hat{x}}$ and by choosing as “worst” trajectory an extremal trajectory for a Barabanov norm associated with $\dot{x} = \dot{\hat{x}}$ (cf. [6]).

In the dwell-time case the equivalence between finiteness of $\gamma_2$ and exponential stability of the associated uncontrolled linear switched system is commonly expected, at least under the minimal realization assumption for each operating mode (see, e.g., [3], [4]). In this paper, on the one hand, we confirm such an expectation under a slightly weaker hypothesis than minimal realization of each mode, but on the other hand we show by means of an example that the finiteness of $\gamma_2$ is not equivalent to the exponential stability of the uncontrolled linear switched system associated with any minimal realization of (1).

The argument for the arbitrary switched case sketched above fails in the dwell-time case simply because the existence of a Barabanov norm for a switched linear system requires the closedness under concatenation of the corresponding class of switching signals (cf. [7]), a condition which is not satisfied by $\Sigma_\tau$ for any $\tau > 0$. This is indeed the main obstacle to address Hespanha’s problem. The lack of concatenability is also the major issue addressed in [8], where a different problem is considered, namely, the continuity of the map $\tau \mapsto \rho(\tau)$. Here $\rho(\tau)$ denotes the generalized spectral radius associated with the switched linear system $\dot{x} = A_\sigma x$, $\sigma \in \Sigma_\tau$. In [8], the lack of concatenability is handled by covering $\Sigma_\tau$ with a family of subsets and by associating with each of them a suitable extremal norm. In this paper we follow a different path: we identify a single subset $\Sigma_\tau$ of $\Sigma_\tau$ that is large enough to encompass the asymptotic properties of the switched linear system $\dot{x} = A_\sigma x$, $\sigma \in \Sigma_\tau$, and well-behaved with respect to concatenation. The flows associated with $\Sigma_\tau$ define a semigroup of matrices, whose analysis allows us to describe the asymptotic worst-case behavior of the original switched system.

Several of our results hold true under a suitable observability
condition, that we call uniform observability, see Definition 22, given for classes $\Sigma$ of switching signals which are not closed under concatenation. In the context of Hespanha’s problem, i.e., for the class $\Sigma_\tau$ of switching signals with dwell-time larger than or equal to $\tau > 0$, the uniform observability assumption amounts to requiring every pair $(A_\sigma, C_\sigma), \sigma \in P$, to be observable and therefore this assumption is weaker than assuming that every $(A_\sigma, B_\sigma, C_\sigma, D_\sigma), \sigma \in P$, satisfies the minimal realization assumption.

Combining our approach with the characterization of controllability and observability for switched linear control systems given in [9] we prove, under the uniform observability condition, that $\gamma_2$ is finite if and only if $\rho_{\min} < 1$, where $\rho_{\min}$ is the generalized spectral radius associated with a minimal realization of the original switched control system (Theorem 27). This result allows one to reduce the specific questions posed by Hespanha in [1, Problem 4.1] about $L_2$-gain issues to the computation of the above mentioned generalized spectral radius. Such a computation can be successfully tackled from a numerical viewpoint thanks to the algorithms proposed in [10], [11].

When the uniform observability assumption does not hold true, the situation is more complicated as we illustrate by means of an example: we construct a 3-dimensional switched linear control system with piecewise constant switching laws, which is controllable, observable, has finite $L_2$-gain, and whose uncontrolled dynamics have generalized spectral radius equal to one (Example 31).

The last result of the paper establishes the right-continuity of the map $\tau \mapsto \gamma_2(\tau)$, providing a partial answer to Question (i) and to Question (iii) in [1, Problem 4.1].

II. PRELIMINARIES

A. Notations

If $n, m$ are positive integers, the set of $n \times m$ matrices with real entries is denoted $M_{n,m}(\mathbb{R})$ and simply $M_n(\mathbb{R})$ if $n = m$. We use $\mathbb{I}_n$ to denote the $n \times n$ identity matrix. A norm on $\mathbb{R}^n$ is denoted $\| \cdot \|$ and similarly for the induced operator norm on $M_n(\mathbb{R})$. A subset $\mathcal{M}$ of $M_n(\mathbb{R})$ is said to be irreducible if the only subspaces which are invariant for each element of $\mathcal{M}$ are $\{0\}$ or $\mathbb{R}^n$. For every $s, t \geq 0$ and $A \in L_\infty([s, s + t], M_n(\mathbb{R}))$, denote by $\Phi_A(s, t, s) \in M_n(\mathbb{R})$ the flow (or fundamental matrix) of $\dot{x}(\tau) = A(\tau)x(\tau)$ from time $s$ to time $s + t$. Given two signals $A_1 : [0, t_1] \rightarrow \mathcal{M}$, $A_2 : [t_1, t_2] \rightarrow \mathcal{M}$, we concatenate $A_1 \ast A_2 : [0, t_1 + t_2] \rightarrow \mathcal{M}$, i.e., the signal coinciding with $A_1(\cdot)$ on $[0, t_1]$ and with $A_2(\cdot - t_1)$ on $(t_1, t_1 + t_2]$. Similarly, if $\mathcal{A}$ and $\mathcal{B}$ are two subsets of signals, we use $\mathcal{A} \ast \mathcal{B}$ to denote the set of signals obtained by concatenation of a signal of $\mathcal{A}$ and a signal of $\mathcal{B}$.

Let $n, p$ and $m$ be positive integers. Consider a switched linear control system of the type

$$\dot{x}(t) = A(t)x(t) + B(t)u(t), \quad y(t) = C(t)x(t) + D(t)u(t),$$

(2)

where $x \in \mathbb{R}^n$, $u \in \mathbb{R}^m$, $y \in \mathbb{R}^p$, $(A, B, C, D)$ belongs to a class $\mathcal{T}$ of measurable switching laws taking values in a bounded set of quadruples of matrices $\mathcal{M} \subset M_n(\mathbb{R}) \times M_{n,m}(\mathbb{R}) \times M_{p,n}(\mathbb{R}) \times M_{p,m}(\mathbb{R})$.

For $t \geq 0$ and a switching law $(A, B, C, D) \in \mathcal{T}$, the controllability and observability Gramians in time $t$ are defined respectively as

$$\int_0^t \Phi_A(0, s)B(s)B(s)^T\Phi_A(0, s)^Tds,$$

$$\int_0^t \Phi_A(s, 0)^T C(s)^T C(s)\Phi_A(s, 0)ds.$$
\( S^{BV,T,ν}(\mathcal{M}) \) is the class of \((T, ν)\)-BV signals, i.e., the signals whose restriction to every interval of length \( T \) has total variation at most \( ν \), that is, \( M \in S^{BV,T,ν}(\mathcal{M}) \) if and only if

\[
\sup_{t-s_0, k \in \mathbb{N}} \left\| M(t_k) - M(t_{k-1}) \right\| \leq ν.
\]

Most of these classes have been already considered in [5]. For the class of persistently exciting signals, see for instance [12], [13], [14] and references therein. Notice that all the classes in the above list are shift-invariant.

Rather than addressing the issues at stake for each class of switching signal given above, we develop a unifying framework which can also be applied to other classes. For that purpose, we adopt an axiomatic approach which singles out some useful common properties satisfied by the classes above.

### III. ADAPTED NORMS FOR SWITCHED LINEAR SYSTEMS WITH CONCATENABLE SUBFAMILIES

We consider in this section a switched linear system

\[
\dot{x}(t) = A(t)x(t)
\]

where \( A \) belongs to a class \( \mathcal{S} \) of measurable switching laws taking values in a bounded nonempty set of matrices \( \mathcal{M} \subset \mathcal{M}_n(\mathbb{R}) \).

A useful assumption on the family \( \mathcal{S} \) that we are going to use in the following (which is satisfied by all the classes introduced in the previous section) concerns its invariance by time-shift.

**A0 (shift-invariance)** For every \( A(\cdot) \in \mathcal{S} \) and every \( t \geq 0 \), the signal \( A(t + \cdot) \) is in \( \mathcal{S} \).

Under Assumption **A0** a convenient measure of the asymptotic behavior of (3) is the **generalized spectral radius** (see, e.g., [7])

\[
\rho(\mathcal{S}) = \lim_{t \to +\infty} \sup_{A \in \mathcal{S}} \left\| \Phi_A(t, 0) \right\|^{1/t}.
\]

Notice that, since \( \mathcal{M} \) is bounded then \( \rho(\mathcal{S}) \) is finite.

As mentioned in introduction, our approach aims at extending the Barabanov norm construction (cf. [7]) beyond the class of signals with arbitrary switching. The main difficulty to do so lies in the fact that the set of all flows \( \Phi_A(s + t, s) \), for \( A \in \mathcal{S} \) and \( s, t \geq 0 \), does not form a semigroup, since in general signals in \( \mathcal{S} \) cannot be concatenated arbitrarily within \( \mathcal{S} \). A key object in what follows is then the identification of a subclass \( \mathcal{F}_x \) of \( \mathcal{S} \), constructed by concatenating in an arbitrary way some signals defined on finite intervals. We then attach to \( \mathcal{F}_x \) a semigroup of fundamental matrices that captures the asymptotic behaviour of \( \mathcal{S} \) if \( \mathcal{F}_x \) is large enough.

**A. Concatenable subfamilies**

Consider a set \( \mathcal{F} = \cup_{t \geq 0} \mathcal{F}_t \) with \( \mathcal{F}_t \subset L_\infty([0, t], \mathcal{M}) \), \( t \in [0, \infty) \). Define

\[
\mathcal{F}_x = \left\{ A_1 * A_2 * \cdots * A_k \cdots \mid A_k \in \mathcal{F}_{t_k} \text{ for } k \in \mathbb{N}, \sum_{k \in \mathbb{N}} t_k = \infty \right\}
\]

and \( \Phi(\mathcal{F}) = \cup_{t \geq 0} \Phi(\mathcal{F}_t) \), where, for every \( t \geq 0 \),

\[
\Phi(\mathcal{F}_t) = \{ A(t, 0) \mid A \in \mathcal{F}_{t} \}.
\]

Let moreover,

\[
\mu(\mathcal{F}) = \lim_{t \to +\infty} \sup \left\{ \left\| R_t \right\|^{1/t} \mid R_t \in \Phi(\mathcal{F}_t) \right\},
\]

with the convention that the quantity inside the parenthesis is equal to 0 if \( \mathcal{F}_t \) is empty. Notice that \( \mu(\mathcal{F}) \leq \rho(\mathcal{F}_{x}) \), but the converse is in general not guaranteed since the computation of \( \rho(\mathcal{F}_{x}) \) takes into account all intermediate instants between two concatenation times, unlike the one of \( \mu(\mathcal{F}) \).

We list below some useful assumptions on the pair \((\mathcal{S}, \mathcal{F})\) that will be exploited in the sequel.

**A1** (concatenability) \( \mathcal{F}_{x} \ast \mathcal{F}_t \subset \mathcal{F}_{x+t} \) for every \( s, t \geq 0 \);

**A2** (irreducibility) \( \Phi(\mathcal{F}) \) is irreducible;

**A3** (fatness) \( \mathcal{F}_x \subset \mathcal{S} \) and there exist two constants \( C, \Delta \geq 0 \) and a compact subset \( \mathcal{K} \) of \( \text{GL}(n) \) such that for every \( t \geq 0 \) and \( A \in \mathcal{S} \), there exist \( K \in \mathcal{K}, \delta \in \{ t, t + \Delta \} \), and \( R \in \Phi(\mathcal{F}_t) \) such that

\[
\left\| \Phi_A(t, 0) KR^{-1} \right\| \leq C.
\]

Moreover, if \( A \in \mathcal{F}_{x} \), one can take \( K = \text{Id}_n \) in (6).

**Remark 1.** As a consequence of the definition of \( \mathcal{F}_x \), if \( A \in \mathcal{F} \) and \( B \in \mathcal{F}_x \), then \( A \ast B \in \mathcal{F}_x \). Moreover, by Assumption **A1**, one has that \( \Phi(\mathcal{F}_x) \subset \Phi(\mathcal{F}_{x+1}) \) for every \( s, t \geq 0 \). Hence \( \Phi(\mathcal{F}) \) is a semigroup and Assumption **A2** above is equivalent to

\[
\forall x \in \mathbb{R}^n \setminus \{0\}, \text{ the linear span of } \Phi(\mathcal{F})x \text{ is equal to } \mathbb{R}^n.
\]

As in [7], one then says that \( \Phi(\mathcal{F}) \) is an irreducible semigroup. Note that [7] considers special classes of irreducible semigroups which verify the additional assumption

**decomposability** \( \mathcal{F}_{x+t} = \mathcal{F}_x \ast \mathcal{F}_t \) for every \( s, t \geq 0 \).

The decomposability assumption trivially implies that \( \Phi(\mathcal{F}_x) \Phi(\mathcal{F}_t) = \Phi(\mathcal{F}_{x+t}) \) for every \( s, t \geq 0 \).

**Remark 2.** Recall that the map \( L_\infty([0, t], \mathcal{M}) \ni A \mapsto \Phi_A(t, 0) \in \mathcal{M}_n(\mathbb{R}) \) is continuous with respect to the weak-* topology in \( L_\infty([0, t], \mathcal{M}) \) (see, for instance, [12, Proposition 21]). In particular if \( A3 \) holds true for \( \mathcal{S} \) then it also holds true for the weak-* closure of \( \mathcal{S} \).

**Lemma 3.** If \( A3 \) holds true then \( \rho(\mathcal{S}) \geq \rho(\mathcal{F}) \). If moreover \( A0 \) and **A1** hold true then the family \( \mathcal{F} \) verifies the following version of Fenichel’s uniformity lemma: assume that there exists \( m > 0 \) such that, for every sequence \( R_j \in \Phi(\mathcal{F}_{t_j}) \) with \( t_j \) tending to infinity, one has \( \lim_{j \to \infty} m^{-t_j} R_j = 0 \). Then \( \mu(\mathcal{F}) < m \).

**Proof:** The inequality \( \rho(\mathcal{S}) \geq \rho(\mathcal{F}) \) is immediate. The opposite one readily comes from Assumption **A3** and the definitions of \( \rho(\mathcal{S}) \) and \( \mu(\mathcal{F}) \). As regard the second part of the lemma, we can assume that \( m = 1 \) by replacing if necessary \( \mathcal{M} \) by the set \( \mathcal{M} - \log(m) \text{Id}_n \). Let \( \mathcal{S}^* \) be the closure of \( \mathcal{S} \) with respect to the weak-* topology induced by \( L_\infty([0, \infty), \mathcal{M}) \). We first show that every trajectory associated with a switching signal in \( \mathcal{S}^* \) tends to zero. Assume by contradiction that there
exist $A \in \mathcal{S}^*$, $x \in \mathbb{R}^n$, $\varepsilon > 0$ and a sequence $t_j$ tending to infinity such that
\[ \|\Phi_A(t_j, 0)x\| \geq \varepsilon. \]

By Remark 2, Assumption A3 actually extends to any switching signal in $\mathcal{S}^*$. Hence, for every $j \geq 0$ applying Assumption A3 to the switching signal $A$ and the time $t_j$ yields the inequality
\[ \|\hat{R}_j K_j^{-1}x\| \geq \varepsilon/C \]
for some $\hat{R}_j \in \Phi(F_{\hat{R}_j})$, $\hat{t}_j \in [t_j, t_j + \Delta]$, and $K_j$ belonging to a given compact of $\text{GL}(n)$. According to the hypotheses of the lemma, the left-hand side of the above inequality tends to $0$ as $j$ goes to infinity, which is a contradiction.

Since the switching laws of $\mathcal{S}$ take values in the bounded set $\mathcal{M}$, one has that the class $\mathcal{S}^*$ is weak-$*$ compact. Recalling that $\mathcal{S}$ is shift-invariant by A0, we notice that all the assumptions of Fenichel’s uniformity lemma are verified by the standard linear flow defined on $\mathbb{R}^n \times \mathcal{S}^*$ (cf. [15]). Therefore the convergence of trajectories of $\mathcal{S}$ to 0 is uniformly exponential, i.e., $\rho(\mathcal{S}) < 1$. By the first part of the lemma we deduce that $\mu(\mathcal{F}) < 1$.

We associate with each class of switching signals considered in the previous section a corresponding concatenable subfamily as listed below:
- $\mathcal{F}^{\text{arb}}(\mathcal{M})$: arbitrarily switching signals on finite intervals;
- $\mathcal{F}^{\text{pc}}(\mathcal{M})$: piecewise constant signals on finite intervals;
- $\mathcal{F}^{d,\tau}(\mathcal{M})$: piecewise constant signals on finite intervals with dwell-time $\tau$ and such that the first and last subintervals on which the signal is constant have length at least $\tau$ (notice that $\mathcal{F}^{d,\tau}(\mathcal{M}) = \emptyset$ for $t < \tau$);
- $\mathcal{F}^{\text{aw},d,\tau}(\mathcal{M})$: piecewise constant signals on finite intervals satisfying the $(\tau, N_0)$ average dwell-time condition and such that the first and last subintervals on which the signal is constant have length at least $N_0\tau$;
- $\mathcal{F}^{\text{pc},T,\mu}(\mathcal{M})$: $T$, $\mu$-persistently exciting signals on finite intervals $[0, t]$ for which $t \geq T - \mu$ and the signal is constantly equal to $M_1$ on $[t - T + \mu, t]$ (where we recall that $\mathcal{M} = \{M_\delta = (1 - \delta)M_0 + \delta M_1 | \delta \in [0, 1]\}$).

For the classes $\mathcal{S}^{\text{arb},L}(\mathcal{M})$ and $\mathcal{S}^{\text{BV},T,\mu}(\mathcal{M})$ we fix some $\tilde{M} \in \mathcal{M}$ and we define
- $\mathcal{F}^{\text{arb},L}(\mathcal{M})$: Lipschitz signals on finite interval starting and ending at $\tilde{M}$;
- $\mathcal{F}^{\text{BV},T,\mu}(\mathcal{M})$: $(T, \nu)$-BV signals on finite intervals $[0, t]$, $t \geq T$, starting and ending at $\tilde{M}$ and constant on $[t - T, t]$.

With these choices of $\mathcal{F}$ Assumption A1 is automatically satisfied. To address the validity of Assumption A2 for the previous classes of signals, we further introduce the following assumption on the set $\mathcal{F}$, which essentially says that the flow corresponding to any element in $\mathcal{F}^{\text{pc}}(\mathcal{M})$ can be approached in a suitable sense by an analytic deformation of flows corresponding to elements in $\mathcal{F}$.

A4 (Analytic propagation) For every $t, \varepsilon > 0$ and $A \in \mathcal{F}^{\text{pc}}(\mathcal{M})$, there exists a path $(0, 1] \ni \lambda \mapsto (t_\lambda, A_\lambda) \in (0, \infty) \times \mathcal{F}^{\text{arb}}(\mathcal{M})$ such that
- $\lambda \mapsto \Phi_{A_\lambda}(t_\lambda, 0)$ is analytic;
- $\|\Phi_{A_\lambda}(t_1, 0) - \Phi_A(t_0, 0)\| \leq \varepsilon$;
- the set $\{\lambda \in (0, 1] | A_\lambda \in \mathcal{F}\}$ has nonempty interior.

The relation between Assumptions A2 and A4 is clarified in the following proposition.

Proposition 4. Let $\mathcal{M}$ be irreducible. Then Assumption A4 implies Assumption A2.

Proof: Let $x, \xi \in \mathbb{R}^n \setminus \{0\}$. We have to prove that there exists $R \in \Phi(\mathcal{F})$ such that $z^T Rx \neq 0$. Since $\mathcal{M}$ is irreducible, it follows from [7, Lemma 3.1] that there exist $t > 0$ and $A$ in $\mathcal{F}^{\text{pc}}(\mathcal{M})$ such that
\[ z^T \Phi_A(t, 0)x \neq 0. \]

Consider the path $\lambda \mapsto A(\lambda, \lambda)$ provided by Assumption A4. The function
\[ \lambda \mapsto z^T \Phi_A(\lambda, 0)x \]
is analytic and not identically equal to zero. It therefore vanishes at isolated values of $\lambda$, whence the conclusion with $R$ of the type $\Phi_{A_\lambda}(t, 0)$. In the following proposition we establish the validity of Assumptions A0, A1, A3 and A4 for the classes of switching signals introduced in Section II-B and their corresponding concatenable subfamilies.

Proposition 5. Let $\mathcal{S}$ be one of the classes introduced in Section II-B with corresponding $\mathcal{F}$ as above. If $\mathcal{S} = \mathcal{S}^{\text{arb},L}(\mathcal{M})$ or $\mathcal{S} = \mathcal{S}^{\text{BV},T,\mu}(\mathcal{M})$, assume moreover that $\mathcal{M}$ is star-shaped, that is, there exists $\tilde{M} \in \mathcal{M}$ such that for any other $M \in \mathcal{M}$ the segment between $M$ and $\tilde{M}$ is contained in $\mathcal{M}$. Then Assumptions A0, A1, A3 and A4 hold true.

Proof: As already noticed, Assumptions A0 and A1 are satisfied.

Concerning Assumption A3, notice that every restriction $A_{[0,t]}$ of a signal in one of the classes $\mathcal{S}$ introduced in Section II-B can be extended to a signal $A_1 * A_{[0,t]} * A_2$ in the corresponding class $\mathcal{F}$, with $A_1 : [0, t_j] \to \mathcal{M}$, $j = 1, 2$, and $t_1, t_2 \leq t_4$ for some $t_4$ uniform with respect to $A \in \mathcal{S}$ and $t \geq 0$. Moreover, if $A \in \mathcal{F}_\infty$ then $t_1$ can be taken equal to zero.

Let us now prove Assumption A4. Take $t, \varepsilon > 0$ and $A \in \mathcal{F}^{\text{pc}}(\mathcal{M})$. For the cases $\mathcal{S} = \mathcal{S}^{\text{arb},L}(\mathcal{M})$ and $\mathcal{S} = \mathcal{S}^{\text{aw},d,\tau}(\mathcal{M})$ one can find the path $\lambda \mapsto A(\lambda, \lambda)$. For every $\delta > 0$, consider the time-reparameterized signal $\tilde{A}(\delta, \cdot) \in \mathcal{F}^{\text{pc}}(\mathcal{M})$. Then $A(\delta, \cdot) \in \mathcal{F}$ for $\delta$ small enough and the function $\delta \mapsto \Phi_{A(\delta)}(t, 0)$ is analytic. Indeed, the function $(0, \infty) \ni \delta \mapsto \Phi_{A(\delta)}(t_0, 0) = \Phi_{A(\delta)}(t_0, 0)$ is analytic, since the Volterra series associated with this flow define an analytic function of $\delta$.

We conclude by taking $A(\lambda, \cdot) = A(\lambda, \lambda)$. For the cases $\mathcal{S} = \mathcal{S}^{\text{arb},L}(\mathcal{M})$ and $\mathcal{S} = \mathcal{S}^{\text{BV},T,\mu}(\mathcal{M})$, the previous construction can be modified as follows. First notice that, up to adding some short intervals on which $A$ is constant we can assume that
- Case $\mathcal{S} = \mathcal{S}^{\text{arb},L}(\mathcal{M})$: the value of $A$ on the last interval on which it is constant is $M_1$;
- Case $\mathcal{S} = \mathcal{S}^{\text{BV},T,\mu}(\mathcal{M})$: the distance between two subsequent values of $A$ is smaller than $\nu$, and the first and last values of $A$ are both equal to $\tilde{M}$. 
The modification of $A$ can be taken so that the corresponding variation of $Φ_A(t,0)$ is smaller than $ε$. The argument now works as before in the case $S = S_{BW,T,ν}^{μ}(ℳ)$. In the case $S = S_{\text{lip},T}^{μ}(ℳ)$ one can take $A_{x}(s) = A(s/λ)$ on $[0,M_0]$ and $A_{x}(s) = M_1$ on $[λ,M_1]$, where $t_0$ is such that $A_{[t_0,t]} \equiv M_1$.

We are left to discuss the case $S = S_{\text{lip},L}^{μ}(ℳ)$. Similarly to what is done above, we can first assume that $A$ is equal to $M$ on the first and last interval on which it is constant. Then we modify $A$ by adding, at each switching time, a Lipschitz continuous arc defined on a small time interval and bridging the discontinuity (one may take a Lipschitz continuous arc whose graph is the union of two segments joining at $M$). These modifications can be done while keeping the corresponding variation of $Φ_A(t,0)$ smaller than $ε$. By a time-reparameterization $λ \mapsto A(λ⋅)$ we can lower the Lipschitz constant and complete the proof as above.

**Remark 6.** In Proposition 5, the hypothesis on the star-shapedness of $ℳ$ can be replaced by some weaker one. For instance we could assume:

- if $S = S_{\text{lip},L}^{μ}(ℳ)$ then there exists $C > 0$ such that every two distinct points of $ℳ$ can be connected by a Lipschitz-continuous curve lying in $ℳ$ of length smaller than $C$;
- if $S = S_{BW,T,ν}^{μ}(ℳ)$ then for every $M_0$, $M_1 ∈ ℳ$, there exists a finite sequence of points in $ℳ$ whose first and last elements are $M_0$ and $M_1$, respectively, and such that the distance between two subsequent elements is smaller than $ν$.

**B. Quasi-Barabanov semigroups**

The main goal of the section is to prove the result below.

**Theorem 7.** Let $(S,F)$ satisfy Assumptions A0–A3. Then there exists a constant $C \geq 1$ such that for any $x_0 ∈ ℜ^n \setminus \{0\}$ there exists a trajectory $t → Φ_A(t,0)x_0$ with $A$ belonging to the weak-∗ closure of $F$, such that, for every $t ≥ 0$,

$$\frac{1}{C} \rho(S)^t \|x_0\| ≤ \|x(t)\| ≤ C \rho(S)^t \|x_0\|. $$

For that purpose, we first need the following definitions.

**Definition 8.** Let $ℳ$, $F$, $F_{≥}^x$, and $Φ(F)$ be as in the previous section.

We say that $Φ(F)$ is a quasi-extremal semigroup if there exists $C_{qE} > 0$ such that, for every $t > 0$ and $R ∈ Φ(F)$, one has

$$\|R\| ≤ C_{qE} μ(F)^t. $$

Moreover, a quasi-extremal semigroup $Φ(F)$ is said to be extremal if there exists a norm $w$ on $ℜ^n$ such that the induced matrix norm $\|\cdot\|_w$ satisfies, for every $t > 0$ and $R ∈ Φ(F)$,

$$\|R\|_w ≤ μ(F)^t. $$

A norm $w$ satisfying Eq. (8) is said to be extremal for $Φ(F)$. A quasi-extremal semigroup $Φ(F)$ is said to be quasi-Barabanov if there exists $C_{qB} > 0$ such that for every $x ∈ ℜ^n$ and $t ≥ 0$ there exist $t' ≥ t$ and $R ∈ Φ(F_{t'})$ such that

$$\|Rx\| ≥ C_{qB} μ(F)^t' \|x\|. $$

Let $Φ(F)$ be a quasi-extremal semigroup. A trajectory $x : t → Φ_A(t,0)x_0$ with $x_0 ≠ 0$ and $A$ belonging to the weak-∗ closure of $F$, is said to be quasi-extremal with constant $C_{qE} ≥ 1$ if for every $t > 0$,

$$\frac{1}{C_{qE}} μ(F)^t \|x_0\| ≤ \|x(t)\| ≤ C_{qE} μ(F)^t \|x_0\|. $$

The notion of generalized spectral radius for a quasi-Barabanov semigroup $Φ(F)$ is actually equivalent to the following adaptation of the definition of maximal Lyapunov exponent

$$\lambda(F) = \sup \left\{ \limsup_{k→∞} \frac{\log \|R_{t_k} \cdots R_{t_1}\|}{t_1 + \cdots + t_k} | \right\}
\begin{align*}
R_{t_k} ∈ Φ(F_{t_k}) & \text{ for } k ∈ ℤ, \sum_{k∈ℤ} t_k = \infty,
\end{align*}
$$
as stated below.

**Proposition 9.** Let $F$ be a family of switching laws satisfying the concatenability condition A1 and assume that $Φ(F)$ is a quasi-Barabanov semigroup. Then $μ(F) = e^{λ(F)}$.

**Proof:** The inequality $μ(F) ≥ e^{λ(F)}$ easily comes from (7). In order to show the opposite inequality one observes that, by (9), for any $x_0 ∈ ℜ^n$ and $i ∈ ℤ$, there exist $t_i ≥ i$ and $R_{t_i} ∈ Φ(F_{t_i})$ such that

$$\|R_{t_k} \cdots R_{t_1} x_0\| ≥ (C_{qB})^k \mu(F)^{t_1 + \cdots + t_k} \|x_0\|. $$

In particular

$$\limsup_{k→∞} \frac{\log \|R_{t_k} \cdots R_{t_1}\|}{t_1 + \cdots + t_k} ≥ - \lim_{k→∞} \frac{2\log(C_{qB})}{k + 1} + \log(μ(F))$$

which concludes the proof.

**Lemma 10.** Let Assumptions A0, A1 and A3 be satisfied. If $Φ(F)$ is a quasi-Barabanov semigroup then there exists $C_{qE} ≥ 1$ such that any nonzero point of $ℜ^n$ is the initial condition of a quasi-extremal trajectory with constant $C_{qE}$.

**Proof:** Let $C_{qE}$, $C_{qB}$ be as in (7), (9) and $C$, $X$ as in A3. Let $κ ≥ 1$ verify $|K|, |K^{-1}| ≤ κ$ for every $K ∈ X$.

Fix $x_0 ≠ 0$. There exists an increasing sequence of times $t_k$ going to infinity and signals $A_k ∈ F_{t_k}$ such that $R_{t_k} x_0 ≥ C_{qB} μ(F)^{t_k} \|x_0\|$, where $R_{t_k} = Φ_A(t_k,0)$.

We now claim that, for every $k ∈ ℤ$ and every $s ∈ [0, t_k]$, one has

$$\|Φ_A(s,0)x_0\| ≥ \frac{1}{C_0} μ(F)^s \|x_0\| $$

for some constant $C_0 > 0$ independent of $x_0$ and $s$. Indeed, because of Assumption A0 and applying (6) one gets that $Φ_{A_k}(t_k,s) = MRK^{-1}$ with $\|M\| ≤ C$, $K ∈ X$, and $R$ in $Φ(F_{t_k-s-δ})$, for some $δ ∈ [0, Δ]$. One can therefore write $R_k = MRK^{-1} Φ_A(s,0)$. It follows that

$$\|R_k x_0\| ≤ \|M\| \|RK^{-1} Φ_A(s,0)x_0\| ≤ CC_{qE} μ(F)^{t_k-s-δ} |K^{-1} Φ_A(s,0)x_0| ≤ κCC_{qE} μ(F)^{t_k-s+δ} \|Φ_A(s,0)x_0\|. $$
On the other hand \(|R_k x_0| \geq C_{q_0} \mu(F)^{t_k} \|x_0\|\), which proves (10) with
\[
C_0 = \frac{\kappa C_{q_0} \max(1, \mu(F)^{t_k})}{C_{q_0}}.
\]

Notice that each \(A_k\) is the restriction on \([0, t_k]\) of a signal \(B_k\) in \(F_k\). Up to a subsequence we can assume that \(B_k\) weak-* converges to some \(B_x\) in the weak-* closure of \(F_x\). Passing to the limit in (10) we deduce that for every \(s \geq 0\),
\[
\|\hat{\Phi}_s(0, 0)x_0\| \geq \frac{1}{C_0} \mu(F)^s \|x_0\|. \tag{11}
\]

We next prove that there exists \(C_1 > 0\) such that for every \(s \geq 0\), \(x_0 \in \mathbb{R}^n\) and \(B_x\) in the weak-* closure of \(F_x\), it holds
\[
\|\hat{\Phi}_s(0, 0)x_0\| \leq C_1 \mu(F)^s \|x_0\|. \tag{12}
\]

For that purpose, consider a sequence \(B_k\) in \(F_x\) weak-* converging to \(B\). Applying (6), Remark 2 and a compactness argument, we get that \(\|\hat{\Phi}_s(0, 0)x_0\| = M R K^{-1}\) with \(M \leq C\), \(K \in \mathscr{K}\) and \(R\) in the closure of \(\cup_{\delta \in [0, \Delta]} \Phi(F_{s+\delta})\). It follows that
\[
\|\hat{\Phi}_s(0, 0)x_0\| \leq \|M\| R K^{-1} \|x_0\| \leq C_1 \mu(F)^s \|x_0\|,
\]
where \(C_1 = \kappa C_{q_0} \max(1, \mu(F)^{t_k})\), proving (12). Together with (11), this concludes the proof of the lemma with \(C_{q_0} = \max(C_0, C_1)\).

Set
\[
\mathscr{R}_x = \{R \mid 3t_k \to \infty, \ R_k \in \Phi(F_{t_k}) \text{ such that } \mu(F)^{-t_k} R_k \to R\}.
\]

The resulting result is the counterpart of [7, Proposition 3.2] in our setting.

**Proposition 11.** Let \((S, F)\) satisfy Assumptions A0–A3 and define \(\mathscr{R}_x\) as above. Then

(i) \(\mathscr{R}_x\) is compact and nonempty, \(\mathscr{R}_x \neq \{0\}\;
(ii) \(\mathscr{R}_x\) is a semigroup;
(iii) for every \(t \geq 0\), \(T \in \Phi(F_t)\) and \(S \in \mathscr{R}_x\), both \(\mu(F)^{-t} TS\) and \(\mu(F)^{-t} ST\) belong to \(\mathscr{R}_x\);
(iv) \(\mathscr{R}_x\) is irreducible.

**Proof:** To prove the proposition one follows exactly the arguments provided in the proof of [7, Proposition 3.2] except for the fact that \(\mathscr{R}_x \neq \{0\}\). In our setting this result is easily proved by using Lemma 3 and the definition of \(\mathscr{R}_x\).

The following result can be proven as in [7, Lemma 3.4].

**Proposition 12.** Let \((S, F)\) satisfy Assumptions A0–A3 and define \(\mathscr{R}_x\) as above. Let \(\hat{v} : \mathbb{R}^n \to (0, \infty)\) be defined as
\[
\hat{v}(x) = \max_{R \in \mathscr{R}_x} \|Rx\|.
\]

Then \(\hat{v}\) is an extremal norm for \(\Phi(F)\).

**Remark 13.** If Assumption A1 is replaced by the stronger decomposability assumption (see Remark 1), then \(\hat{v}\) is a Barabanov norm (see e.g. [7]).

We have the following result.

**Proposition 14.** Let \((S, F)\) satisfy Assumptions A0–A3. Then \(\Phi(F)\) is an extremal and quasi-Barabanov semigroup.

**Proof:** The fact that \(\Phi(F)\) is an extremal semigroup readily comes from Proposition 12. In order to show the second part of the statement, let us consider \(\kappa \geq 1\) as in the proof of Lemma 10. Without loss of generality, let us also assume that \(\kappa\) satisfies
\[
\frac{\lambda}{\kappa \|x\|} \leq \hat{v}(x) \leq \lambda \kappa\|x\| \text{ for all } x \in \mathbb{R}^n
\]

and that \(\|M\| \leq C\) implies that \(\|M\| \leq \kappa\), where \(C\) is as in A3.

For every \(x_0 \in \mathbb{R}^n\), there exists a sequence \((R_k)_{k \in \mathbb{N}}\) so that \(R_k \in \Phi(F_{t_k})\) with \(3t_k \to \infty\) and \(\hat{v}(x_0) = \lim_{k \to \infty} \mu(F)^{-t_k} \|R_k x_0\|\). For every \(k \in \mathbb{N}\), let \(A_k \in F_{t_k}\) be such that \(R_k = A_k(t_k, 0)\).

Fix now \(s \geq 0\). For every \(k \in \mathbb{N}\) such that \(t_k \geq s\), we deduce from Assumption A3 that \(A_k(0, 0) = M_k Q_1^k(k)\) and \(A_k(s, 0) = M_2^k Q_2^k(k) (K_2^k)\), with \(\|M_2^k\| \leq C\), \(K_2^k \in \mathscr{K}\), \(Q_1^k \in \Phi(F_{s+\delta_1^k})\), and \(Q_2^k \in \Phi(F_{t_k-s+\delta_2^k})\), where \(\delta_i^k \in [0, \Delta]\), for \(i = 1, 2\). One can therefore write \(R_k = M_2^k Q_2^k(k) (K_2^k)^{-1} M_1^k Q_1^k(k)\). Using the extremality of \(\hat{v}\) and the requirements imposed on \(\kappa\), it follows that
\[
\mu(F)^{-t_k} \|R_k x_0\| \leq \kappa \mu(F)^{-t_k} \hat{v}(R_k x_0)
\]
\[
\leq \kappa^6 \mu(F)^{-s+\delta_2^k} \hat{v}(Q_1^k(k) x_0).
\]

Taking limits as \(k\) tends to infinity and up to subsequences, one gets
\[
\hat{v}(x_0) \leq \kappa^6 \mu(F)^{-s+\delta_2} \hat{q}(Q_1 x_0),
\]
where \(Q_1\) is in the closure of \(\cup_{\delta \in [0, \Delta]} \Phi(F_{s+\delta})\), and \(\delta_2\) belongs to \([0, \Delta]\). In particular, there exist \(Q \in \Phi(F_{s+\delta})\) for some \(\delta \in [0, \Delta]\) such that
\[
\hat{v}(x_0) \leq 2 \kappa^6 \mu(F)^{-s+\delta_2} \hat{q}(Q x_0).
\]

One deduces that \(\hat{v}(Q x_0) \geq C_0 \mu(F)^{s+\delta} \hat{v}(x_0)\) with \(C_0 = \min_{1 \leq s \leq 2} \lambda^{\Delta} (S)\).

As a consequence of Lemma 10, we have the following corollary.

**Corollary 15.** Let \((S, F)\) satisfy Assumptions A0–A3. Then there exists \(C_{q_0} \geq 1\) such that any nonzero point of \(\mathbb{R}^n\) is the initial condition of a quasi-extremal trajectory with constant \(C_{q_0}\).

Theorem 7 follows directly from Corollary 15 and Lemma 3.

Note that, as a consequence of A3 and the above results, the right-hand side inequality in the statement of Theorem 7 holds true for any trajectory associated with a signal in \(S\), up to adapting the constant \(C\).

**Remark 16.** As a consequence of Proposition 9, Proposition 14 and Lemma 3 one easily deduces that \(\rho(S) = e^{\lambda(S)}\), where
\[
\lambda(S) = \sup_{A \in S} \limsup_{t \to +\infty} \frac{\log \|\Phi_A(t, 0)\|}{t}
\]
is the maximal Lyapunov exponent associated with the family \(S\). This result was already obtained in [8] if the class \(S\) is assumed to be weak-* closed.
IV. $L_2$-GAIN FOR SWITCHED LINEAR CONTROL SYSTEMS

Consider a switched linear control system of the type

$$
\dot{x}(t) = A(t)x(t) + B(t)u(t), \quad y(t) = C(t)x(t),
$$

(13)

where $(x, u, y) \in \mathbb{R}^n \times \mathbb{R}^m \times \mathbb{R}^p$ and $(A, B, C)$ belongs to a class $\mathcal{T}$ of measurable switching laws taking values in a bounded set of triplets of matrices $\mathcal{M} = M_n(\mathbb{R}) \times M_m(\mathbb{R}) \times M_p(\mathbb{R})$. We denote by $\pi_A$ and $\pi_{AXB}$ the projections from $M_n(\mathbb{R}) \times M_m(\mathbb{R}) \times M_p(\mathbb{R})$ to its first and first two factors, respectively. We set

$$
\mathcal{M}_A = \pi_A(\mathcal{M}) = \{ A \mid \exists (A, B, C) \in \mathcal{M} \},
$$

$$
\mathcal{T}_A = \pi_A(\mathcal{T}) = \{ A \mid (A, B, C) \in \mathcal{T} \}
$$

and we define similarly $\mathcal{M}_B$, $\mathcal{M}_C$, $\mathcal{T}_B$ and $\mathcal{T}_C$.

In the sequel, we also assume that $\mathcal{T}$ contains a subset $\mathcal{G}$ of the family $\mathcal{G} = \bigcup_{i \geq 0} \mathcal{G}_i$ as in Eq. (5). Useful properties on $\mathcal{T}$ and $\mathcal{G}$ are the following:

- **T1** (T$_A$, G$_A$) satisfies Assumptions A0 and A1, where G$_A = \pi_A(\mathcal{G})$.
- **T2** (T$_A$, G$_A$) satisfies Assumption A4. Moreover, for every $t^*, \varepsilon > 0$ and $(A, B, C) \in \mathcal{F}^p_{t^*}(\mathcal{M})$, there exists a path $[0, 1] \ni \lambda \mapsto (t_A, B_A, C_A) \in (0, \infty) \times \mathcal{F}^{tb}_t(\mathcal{M})$ such that:
  - $\lambda \rightarrow (W^x(t_A), W^y(t_A))$ is analytic, where $W^x(t_A)$ and $W^y(t_A)$ denote, respectively, the controllability and observability Gramians in time $t_A$ associated with $\dot{x}(t) = A_1(t)x(t) + B_1(t)u(t)$, $y(t) = C(t)x(t)$;
  - $|W^x(t_0) - W^x(t_1)| - \varepsilon \leq \|W^x(t_1) - W^x(t_0)\| \leq \varepsilon$, and $\|W^y(t_0) - W^y(t_1)\| \leq \varepsilon$, where $W^c(t)$ and $W^o(t)$ denote, respectively, the controllability and observability Gramians in time $t$ associated with $\dot{x}(t) = A(t)x(t) + B(t)u(t)$, $y(t) = C(t)x(t)$;
  - the set $\{ \lambda \in [0, 1] \mid (A_1, B_1, C_1) \in \mathcal{G}\}$ has nonempty interior.

A trivial adaptation of the proof of Proposition 5 yields the following result.

**Lemma 17.** Let S be one of the classes introduced in Section II-B with corresponding $\mathcal{F}$. Assume moreover that if $\mathcal{S} = S^{\text{lip},L}(\mathcal{M})$ or $\mathcal{S} = S^{\text{BV},T',\nu}(\mathcal{M})$ then $\mathcal{M}$ is star-shaped. Then $(\mathcal{S}, \mathcal{F})$ satisfies Assumptions T1 and T2.

A. Minimal realization for switched linear control systems

We start by giving the following definitions.

**Definition 18.**

1. A point $x \in \mathbb{R}^n$ is $G$-reachable for the switched linear control system (13) if there exist $t \geq 0$, a switching law $(A, B, C) \in \mathcal{G}_t$ and an input $u \in L_2$ such that the corresponding trajectory $x_u$ starting at 0 reaches $x$ in time $t$. The reachable set $\mathcal{R}(G)$ is the set of all $G$-reachable points. System (13) is said to be $G$-controllable if $\mathcal{R}(G) = \mathbb{R}^n$.

2. A point $x \in \mathbb{R}^n$ is $G$-observable for the switched linear control system (13) if there exist $t \geq 0$ and a switching law $(A, B, C) \in \mathcal{G}_t$ such that the trajectory $x_0$ associated with the zero input and starting at $x$ gives rise to an output $y$ verifying $y(t) \neq 0$. The observable set $\mathcal{O}(G)$ is the set of all $G$-observable points. System (13) is said to be $G$-observable if $\mathcal{O}(G) = \mathbb{R}^n$.

It is not immediate from their definitions that the reachable and observable sets $\mathcal{R}(G)$ and $\mathcal{O}(G)$ are linear subspaces. It has been shown in [9] that this is the case if $T = \mathcal{G}_x = S^{\text{arb}}(\mathcal{M})$, where $\mathcal{M} = \{(A_1, B_1, C_1), \ldots, (A_k, B_k, C_k)\}$ with $k$ a positive integer. In addition, it is proved in the same reference that the state space admits a direct sum decomposition into a controllable and an uncontrollable part for the switched linear control system exactly as in the unswitched situation. More precisely, there exists a direct sum decomposition of the state space $\mathbb{R}^n = \mathcal{R}(G) \oplus E$ and an invertible $n \times n$ matrix $P$ such that, if $r = \dim \mathcal{R}(G)$ and $P^{-1}x = (x^c, x^u)$ with $x^c \in \mathcal{R}(G)$ and $x^u \in E$, one has for $1 \leq i \leq k$,

$$
P^{-1}A_iP = \begin{pmatrix} A_i^c & 0 \\ 0 & A_i^u \end{pmatrix}, \quad P^{-1}B_i = \begin{pmatrix} B_i^c \\ 0 \end{pmatrix}, \quad C_iP = \begin{pmatrix} C_i^c \\ 0 \end{pmatrix},
$$

where $A_i^c$ and $B_i^c$ belong to $M_r(\mathbb{R})$ and $M_{r,m}(\mathbb{R})$ respectively. Moreover, the switched linear control system defined on $\mathbb{R}^r$ associated with $S^{\text{arb}}(\mathcal{M})$, where

$$
\mathcal{M}^c = \{ (A_1^c, B_1^c, C_1^c), \ldots, (A_k^c, B_k^c, C_k^c) \},
$$

is $\mathcal{F}^{\text{arb}}(\mathcal{M}^c)$-controllable. We refer to $\mathcal{M}^c$ as the controllable part of $\mathcal{M}$. Notice that $\mathcal{M}^c = \Pi^c(\mathcal{M})$ where

$$
\Pi^c(A, B, C) = (UP^{-1}APUT, UP^{-1}B, CPUT),
$$

with $U = (Id_{r,n}, 0_{r,n-r})$. Also notice that the output $y$ corresponding to the original system is equal to $y = C^c x^c$ and thus the original switched linear control system has the same $L_2$-gain as the one reduced to the reachable space.

Similarly, there exists a direct sum decomposition of the state space $\mathbb{R}^n = \mathcal{O}(G) \oplus F$ and an invertible $n \times n$ matrix $Q$ such that, if $s = \dim \mathcal{O}(G)$ and $Q^{-1}x = (x^c, x^u)$ with $x^c \in \mathcal{O}(G)$ and $x^u \in F$, one has for $1 \leq i \leq k$,

$$
Q^{-1}A_iQ = \begin{pmatrix} A_i^o & 0 \\ 0 & A_i^o \end{pmatrix}, \quad Q^{-1}B_i = \begin{pmatrix} B_i^o \\ 0 \end{pmatrix}, \quad C_iQ = \begin{pmatrix} C_i^o \\ 0 \end{pmatrix},
$$

where $A_i^o$ and $C_i^o$ belong to $M_s(\mathbb{R})$ and $M_{p,s}(\mathbb{R})$ respectively. Moreover, the switched linear control system defined on $\mathbb{R}^s$ associated with $S^{\text{arb}}(\mathcal{M}^o)$, where

$$
\mathcal{M}^o = \{ (A_1^o, B_1^o, C_1^o), \ldots, (A_k^o, B_k^o, C_k^o) \},
$$

is $\mathcal{F}^{\text{arb}}(\mathcal{M}^o)$-observable. We refer to $\mathcal{M}^o$ as the observable part of $\mathcal{M}$. Notice that $\mathcal{M}^o = \Pi^o(\mathcal{M})$ where

$$
\Pi^o(A, B, C) = (VQ^{-1}AQVT, VQ^{-1}B, CVT),
$$

and $V = (Id_s, 0_{s,n-s})$. The corresponding output $y$ being equal to $y = C^o x^o$, one deduces the equality of the $L_2$-gains of the original switched linear control system and of the one reduced to the observable space.

From the above, one can proceed as follows. Consider a switched linear control system (13) associated with $S^{\text{arb}}(\mathcal{M})$, where $\mathcal{M} = \{(A_1, B_1, C_1), \ldots, (A_k, B_k, C_k)\}$. One first reduces it to its reachable space $\mathcal{R}(\mathcal{F}^{\text{arb}}(\mathcal{M}))$. We thus
get a $F_{\text{arb}}(\Pi^*(M))$-controllable switched linear control system with same $L_2$-gain as that of the original system. Then, one reduces the latter system to its observable space $\mathcal{O}(F_{\text{arb}}(\Pi^*(M)))$ to finally obtain a switched linear control system associated with $\mathcal{S}_{\text{arb}}(\mathcal{M}_{\text{min}})$ where $\mathcal{M}_{\text{min}} = \Pi^0(\Pi^*(M))$. The latter system is finally $F_{\text{arb}}(\mathcal{M}_{\text{min}})$-controllable and $F_{\text{arb}}(\mathcal{M}_{\text{min}})$-observable and

$$\gamma_2(S_{\text{arb}}(\mathcal{M}_{\text{min}})) = \gamma_2(S_{\text{arb}}(\mathcal{M})).$$

(14)

We refer to $S_{\text{arb}}(\mathcal{M}_{\text{min}})$ and the corresponding switched linear control system as a minimal realization for the linear switched linear control system associated with $S_{\text{arb}}(\mathcal{M})$ and (13). We also say that $n'$ is the dimension of such a minimal realization.

**Remark 19.** Note that even though the dimension $n'$ of $\mathcal{O}(F_{\text{arb}}(\Pi^*(M)))$ is uniquely defined by the original switched linear control system, the minimal realization is not unique since it depends on the choice of supplementary subfamilies of switching signals introduced in Section III-A. Section II-B together with its corresponding concatenation of realization associated with (13) and any class $\mathcal{M}$ of realization theory of switched linear control systems and we get a $\mathcal{M}_{\text{min}}$-controllable switched linear control system associated with $\mathcal{S}_{\text{arb}}(\mathcal{M})$ and (13). All the results presented in this paragraph belong to the theme of realization theory of switched linear control systems and we refer to [16] for a thorough presentation of such a theory.

Finally, it must be recalled that [9] also provides a nice and explicit geometric description of $\mathcal{R}(F_{\text{arb}}(\mathcal{M}))$ and $\mathcal{O}(F_{\text{arb}}(\mathcal{M}))$ in terms of the data of the problem. Let us recall here the details of such a geometric description for $\mathcal{R}(F_{\text{arb}}(\mathcal{M}))$ (the corresponding results for $\mathcal{O}(F_{\text{arb}}(\mathcal{M}))$ being standardly derived by duality).

We first need the following notation. If $A$ is an $n \times n$ matrix and $B$ is a subspace of $\mathbb{R}^n$, let $\Gamma_A B$ be the subspace of $\mathbb{R}^n$ given by

$$\Gamma_A B = B + AB + \cdots + A^{n-1} B.$$  

For $1 \leq i \leq k$, let $D_i = \text{Im}[B_i A_1 B_i \cdots A_i^{n-1} B_i]$. Moreover (see [9, Section 3.1]), define recursively the sequence of subspaces of $\mathbb{R}^n$ denoted $V_j, j \geq 1$, by

$$V_1 = D_1 + \cdots + D_k,$$

$$V_{j+1} = \Gamma_A V_j + \cdots + \Gamma_A V_j,$$

and finally set $\mathcal{V}(\mathcal{M}) = \bigcup_{j \geq 1} V_j$. From the variation of constants formula, it is not difficult to see that $\mathcal{R}(F_{\text{arb}}(\mathcal{M}))$ is included in $\mathcal{V}(\mathcal{M})$. The converse inclusion is also true but more delicate to establish, cf. [9, Theorem 1] and the proof of Proposition 20 below. In the subsequent paragraphs, we will use these results to derive the existence of a minimal realization associated with (13) and any class $\mathcal{T}$ considered in Section II-B together with its corresponding concatenation of subfamilies of switching signals introduced in Section III-A.

We now generalize the above construction to a bounded set $\mathcal{M} \subset M_n(\mathbb{R}) \times M_{n,m}(\mathbb{R}) \times M_{p,n}(\mathbb{R})$. We associate a subspace $\mathcal{V}(\mathcal{M})$ of $\mathbb{R}^n$ as follows. First, consider

$$\mathcal{V}_1(\mathcal{M}) = \text{Span}\{A_1 b_l \mid 0 \leq j \leq n-1, 1 \leq l \leq m, (A, [b_1 \ldots b_m]) \in \pi_{A \times B}(\mathcal{M})\}.$$  

Then, define recursively for $j \geq 1$

$$\mathcal{V}_{j+1}(\mathcal{M}) = \text{Span}\{A_j^i v \mid 0 \leq j \leq n-1, A \in \pi_{A}(\mathcal{M}), v \in \mathcal{V}_j(\mathcal{M})\},$$

and finally set $\mathcal{V}(\mathcal{M}) = \sum_{j \geq 1} \mathcal{V}_j(\mathcal{M})$. Taking a (finite) generating family of $\mathcal{V}(\mathcal{M})$, one can extract a finite subset $\mathcal{M}_{\text{finite}}$ of $\mathcal{M}$ such that $\mathcal{V}(\mathcal{M}_{\text{finite}}) = \mathcal{V}(\mathcal{M})$. Hence

$$\mathcal{V}(\mathcal{M}_{\text{finite}}) = \mathcal{V}(\mathcal{M}),$$

(16)

where the first inclusion is deduced from the variation of constants formula. Therefore, all the sets appearing in (16) coincide.

We thus prove the following proposition.

**Proposition 20.** Consider a switched linear control system of the type (13) associated with a class $\mathcal{T}$ of measurable switching laws taking values in a bounded set $\mathcal{M} \subset M_n(\mathbb{R}) \times M_{n,m}(\mathbb{R}) \times M_{p,n}(\mathbb{R})$. Let $(T, G)$ satisfy Assumption T2. Then $\mathcal{R}(G) = \mathcal{R}(F_{\text{arb}}(\mathcal{M})) = \mathcal{V}(\mathcal{M})$ and there exist $t^* > 0$ and a switching law $(A, B, C) \in \mathcal{G}_{t^*}$ such that the range of $W^c(t^*)$, the controllability Gramian in time $t^*$ associated with $\hat{x}(t) = A(t)x(t) + B(t)u(t)$, is equal to $\mathcal{V}(\mathcal{M})$.

**Proof:** First notice that the equality $\mathcal{R}(F_{\text{arb}}(\mathcal{M})) = \mathcal{V}(\mathcal{M})$ is contained in Eq.(16) and one has the trivial inclusion $\mathcal{R}(G) \subset \mathcal{V}(\mathcal{M})$. Let us prove the opposite inclusion.

As done above there exists a finite subset $\mathcal{M}_{\text{finite}}$ of $\mathcal{M}$ such that $\mathcal{V}(\mathcal{M}_{\text{finite}}) = \mathcal{V}(\mathcal{M})$. It is proved in [9, Theorem 1] that there exists a piecewise-constant periodic switching law $(A, B, C)$ taking values in $\mathcal{M}_{\text{finite}}$ and a time $t^* > 0$ such that the range of $W^c(t^*)$, the controllability Gramian in time $t^* > 0$ associated with $\hat{x}(t) = A(t)x(t) + B(t)u(t)$, is equal to $\mathcal{V}(\mathcal{M})$. Fix $\varepsilon > 0$ such that if an $n \times n$ matrix $W$ satisfies $\|W - W^c(t^*)\| < \varepsilon$ then the rank of $W$ is larger than or equal to the rank of $W^c(t^*)$. Let $\lambda \mapsto (A_\lambda, B_\lambda, C_\lambda)$ be the path provided by Assumption T2. Then there exists $\lambda^* \in [0, 1]$ such that $(A_{\lambda^*}, B_{\lambda^*}, C_{\lambda^*}) \in \mathcal{G}$ and the rank of the corresponding controllability Gramian $W_{\lambda^*}^c(t_{\lambda^*})$ is larger than or equal to the rank of $W^c(t^*)$, itself equal to dim $\mathcal{V}(\mathcal{M})$. Since the range of $W_{\lambda^*}^c(t_{\lambda^*})$ is included in $\mathcal{R}(G)$ this concludes the proof of the proposition.

An analogous statement is obtained regarding observability spaces and, following the first part of the section, one finally derives the subsequent result.

**Proposition 21.** Consider a switched linear control system of the type (13) associated with a class $\mathcal{T}$ of measurable switching laws taking values in a bounded set $\mathcal{M} \subset M_n(\mathbb{R}) \times M_{n,m}(\mathbb{R}) \times M_{p,n}(\mathbb{R})$. Let $(T, G)$ satisfy Assumption T2. Let $n'$ be the dimension of any minimal realization of the switched linear control system associated with $S_{\text{arb}}(\mathcal{M})$. Pick one such
minimal realization and consider the corresponding surjective linear mapping \( \Pi \) from \( M_n(\mathbb{R}) \times M_{n,m}(\mathbb{R}) \times M_{p,n}(\mathbb{R}) \) to \( M_{n'}(\mathbb{R}) \times M_{n',m}(\mathbb{R}) \times M_{p,n'}(\mathbb{R}) \) describing its matrix representation. Denote by \( T^{\text{min}} \) the class \( T^{\text{min}} = \{ t \mapsto \Pi(A(t), B(t), C(t)) \mid (A, B, C) \in T \} \) taking values in \( \mathcal{M}^{\text{min}} = \Pi(\mathcal{M}) \). Then the switched linear control system corresponding to \( T^{\text{min}} \) is \( G^{\text{min}} \)-controllable and \( G^{\text{min}} \)-observable in the sense of Definition 18 with \( G^{\text{min}} = \Pi(G) \). Moreover, \( (T^{\text{min}}, G^{\text{min}}) \) satisfies Assumption T2 and \( \gamma_2(T^{\text{min}}) \). Finally, if \( (T, G) \) satisfies Assumption T1 then \( (T^{\text{min}}, G^{\text{min}}) \) does.

Proof: This proposition is a simple consequence of the construction of minimal realizations in the case of arbitrary switching, as described in the first part of the section, and of Proposition 20.

For \( T \) and \( T^{\text{min}} \) as in the statement of Proposition 21, we say that the switched linear control system corresponding to \( T^{\text{min}} \) is a minimal realization of the switched linear control system associated with \( T \). It follows from Remark 19 that any two such minimal realizations are algebraically similar.

### B. Finiteness of the \( L_2 \)-gain

Consider the switched linear control system (13) and the corresponding class \( T \) of switching laws with values in the bounded set \( \mathcal{M} \) of triples of matrices. Let us introduce the following definition.

**Definition 22.** We say that (13) is uniformly observable (or, equivalently, that \( T \) is uniformly observable) if there exist \( T, \gamma > 0 \) such that, for every \( (A, B, C) \in T \) and every \( t \geq 0 \), one has \( W^\alpha(t, t + T) \geq \gamma 1_{\mathbb{R}_+} \), where \( W^\alpha(t, t + T) \) is the observability Gramian in time \( T \) associated with \( (A(t + \cdot), B(t + \cdot), C(t + \cdot)) \).

The following remark will be used in the sequel.

**Remark 23.** If (13) is uniformly observable and \( T, \gamma \) are as in Definition 22, the observability Gramian \( W^\alpha(t, t + T) \) in time \( T \) associated with a switching law belonging to the weak-* closure of \( T \) still satisfies \( W^\alpha(t, t + T) \geq \gamma 1_{\mathbb{R}_+} \).

**Remark 24.** Consider the case where \( T \) contains all the constant \( \mathcal{M} \)-valued switching signals. (Notice that this is the case for all classes of switching signals introduced in Section II-B except that of persistently exciting signals.) It is then easy to see that uniform observability implies that \( (A, C) \) is observable for every \( (A, B, C) \in \mathcal{M} \). In the case where \( \mathcal{M} \) is compact and the class of signals under consideration is \( S^A, \tau(\mathcal{M}) \) for some \( \tau > 0 \), one easily shows that the converse is also true, namely, the observability of each pair \( (A, C) \) implies uniform observability. Let us stress that the uniform observability assumption is weaker than the minimal realization assumption (every triple \( (A, B, C) \in \mathcal{M} \) is a minimal realization, i.e., \( (A, B) \) is controllable and \( (A, C) \) is observable) needed in [2], [3].

We can now state the main result of this section.

**Theorem 25.** Assume that (13) admits a minimal realization defined on \( \mathbb{R}^n, n' \leq n \), which is given by \( \dot{x}^\text{min}(t) = A^\text{min}(t)x(t) + B^\text{min}(t)u(t) \) with output \( y^\text{min}(t) = C^\text{min}(t)x(t) \) and associated with a class \( T^{\text{min}} \) of switching signals taking values in \( \mathcal{M}^{\text{min}} \), and a family \( G^{\text{min}} \) satisfying the following assumptions:

(a) \( (T^{\text{min}}, G^{\text{min}}) \) satisfies Assumptions T1 and T2;

(b) for every subspace \( V \) of \( \mathbb{R}^{n'} \) and every linear projection \( P_V : \mathbb{R}^{n'} \to V \), the class of signals \( P_V G^{\text{min}} \) satisfies Assumption A3, where we use \( P_V^\# \) to denote the dual map from \( V \) to \( \mathbb{R}^{n'} \) defined by \( x^TP_V y = y^TP_V^\# x \) for every \( x \in V \) and \( y \in \mathbb{R}^n \).

Then \( \gamma_2(T) \) is finite if \( \rho(T^{\text{min}}) < 1 \) and infinite if either \( \rho(T^{\text{min}}) > 1 \) or \( \rho(T^{\text{min}}) = 1 \) and \( T^{\text{min}} \) is uniformly observable.

Proof: Thanks to Proposition 21, it is enough to treat the case \( T = T^{\text{min}} \).

Assume first that \( \rho(T_A) < 1 \). Taking into account the definition of \( \rho(T_A) \) (see (4)) and the boundedness of \( \pi_A(\mathcal{M}) \), one gets the following exponential decay estimate: there exist \( K_1, \lambda > 0 \) such that, for every \( A \in T_A \) and every \( 0 \leq s \leq t \), one has

\[
\|\Phi_A(t, s)\| \leq K_1 e^{-\lambda(t-s)}.
\]

As a consequence of the above and the boundedness of \( \mathcal{M} \), one deduces that there exists \( K_2 > 0 \) such that, for every \( u \in L_2, (A, B, C) \in T \) and \( t \geq 0 \), one has

\[
\|y_u(t)\| \leq K_2 \int_0^t e^{-\lambda(t-s)}\|u(s)\| ds.
\]

If \( \chi_{[0, +\infty)} \) denotes the characteristic function of \( [0, +\infty) \), the integral function on the right-hand side of (17) can be interpreted as the convolution of

\[
f_1(t) = \chi_{[0, +\infty)}(t)e^{-\lambda t}, \quad f_2(t) = \chi_{[0, +\infty)}(t)\|u(t)\|.
\]

That yields at once that \( \|y_u\|_2 \leq \frac{K_2}{\lambda} \|u\|_2 \), hence the conclusion.

Assume now that \( \rho(T_A) \geq 1 \). It is well-known (see, e.g., [17, Proposition 2] for details), that, up to a common linear change of coordinates, every matrix in \( A \in \mathcal{M} \) has the upper triangular block form

\[
A = \begin{pmatrix}
A_{11} & A_{12} & \cdots & A_{1q} \\
0 & A_{22} & A_{23} & \cdots \\
0 & 0 & A_{33} & A_{34} & \cdots \\
\vdots & \vdots & \ddots & \ddots & \ddots \\
0 & \cdots & \cdots & \cdots & A_{qq}
\end{pmatrix},
\]

where, for \( i = 1, \ldots, q \), each \( A_{ii} \in M_{n_i-n_{i-1}}(\mathbb{R}) \), \( n_i \in \mathbb{N} \) and the set \( \mathcal{A}_i := \{ A_{ii} \mid A \in \mathcal{M}_A \} \) is irreducible (whenever \( \mathcal{M}_A \) is irreducible one has \( \mathcal{A}_i \). Define the subsystems of \( \mathcal{M}_A \) as the switched systems corresponding to the sets \( \mathcal{A}_i \) and the class of switching signals \( T_{A,i} := \{ A_{ii} \mid A \in T_A \} \) and \( G_{A,i} := \{ A_{ii} \mid A \in G_A \} \). One can then show that \( (T_{A,1}, G_{A,1}), \ldots, (T_{A,q}, G_{A,q}) \) satisfy Assumptions A0–A3. Indeed, Assumptions A0 and A1 follow from Assumption T1; Assumption A2 is a consequence of Assumption T2 and Proposition 4 while Item (b) yields Assumption A3.
Moreover, an induction argument on the number of subsystems and a standard use of the variation of constants formula yields

$$\rho(T_A) = \max_{1 \leq i \leq q} \rho(T_{A,i}).$$

Let $i \leq q$ be the largest index such that $\rho(T_{A,i}) = \rho(T_A)$. Since (13) is $G$-controllable, there exists a time $t > 0$, a switching law $(A^0, B^0, C^0) \in G$ and a measurable control $\bar{u}$ defined on $[0, \bar{t}]$ so that the corresponding trajectory $\bar{x}_\bar{u}$ starting at $0$ reaches some point $\bar{x} = (0, \ldots, 0, \bar{x}, 0, \ldots, 0)^T \neq 0$ in time $\bar{t}$.

Since $(T_{A,i}, G_{A,i})$ satisfies the hypotheses of Theorem 7, the corresponding semigroup $\Phi(G_{A,i})$ is quasi-Baranov and admits a quasi-extremal trajectory $x^{\bar{u}, \text{ext}}_\bar{u}$ starting at $x$ and corresponding to a signal $A_i$ belonging to the weak-* closure of $(G_{A,i})_{<\infty}$. Let $\bar{x}$ be a trajectory of $x = A(t)x(t)$ starting at $\bar{x}$ and corresponding to some signal $\bar{A}$ in the closure of $(G_{A,i})_\infty$ so that $\bar{A}_{\bar{t}} = A_{\bar{t}}$. Notice that $\bar{x} = x^{\bar{u}, \text{ext}}_\bar{u}$.

Let $\bar{E}$ be the signal defined as the concatenation of $A^0$ and $\bar{A}$ and let $\bar{A}^0_{\bar{t}}$ and $\bar{E}_{\bar{t}}$ be the $(\bar{t}, \bar{t})$-blocks corresponding to $A^0$ and $\bar{E}$ respectively. Let $(E^0_{\bar{t}})_{\bar{t} > 0}$ and $(E_{\bar{t}})_{\bar{t} > 0}$ be a sequence in $(G_{A,i})_{\infty} (G_{A,i})_{\infty}$ respectively weak-* converging to $\bar{A}$ ($\bar{E}_{\bar{t}}$ respectively). Let $(F^i, G^i)_{\bar{t} > 0}$ be such that $(E^0_{\bar{t}}E^i_{\bar{t}})_{\bar{t} > 0} \in (G_{A,i})_{\infty}$ and weak-* converges to $\bar{E}$ ($\bar{E}_{\bar{t}}$ respectively). For every positive time $t \geq \bar{t}$, the convergence of $s \to \Phi_{E^i}(s-\bar{t}, 0)$ to $s \to \Phi_{\bar{E}}(s, \bar{t})$ is uniform with respect to $s \in [\bar{t}, \bar{t}]$ and similarly for the corresponding $(\bar{t}, \bar{t})$-blocks.

Denote by $x^i$ the trajectory of the original system and associated with the control equal to $\bar{u}$ on $[0, \bar{t}]$ and zero on $[\bar{t}, +\infty)$, and with the switching signal $(A^0 \ast E^i, B^0 \ast F^i, C^0 \ast G^i)$. Let $y^i$ be the corresponding output. Notice that $x^i(s) = \Phi_{E^i}(s-\bar{t}, 0)\bar{x}$ for every $s \geq \bar{t}$. For $l$ large enough and for every $s \in [\bar{t}, \bar{t}]$, one has

$$\|x^i(s)\| \geq \|\Phi_{E^i}(s-\bar{t}, 0)\bar{x}\| \geq \frac{1}{2}\|\Phi_{E_{\bar{t}}}(\bar{t}, \bar{t})\bar{x}\| \geq \frac{1}{2C_{\text{qu}}} \rho(T_A)^{s-\bar{t}} \|\bar{x}\|.$$

Hence, there exists a positive constant $K_3$ independent of $t \geq \bar{t}$ such that, for every $l$ large enough and for every $s \in [\bar{t}, \bar{t}]$, one has

$$\|x^i(s)\| \geq K_3 \rho(T_A)^{s-\bar{t}}. \quad (18)$$

Assume now that $\rho(T_A) = 1$ and $T$ is uniformly observable, and let $T, \gamma$ be as in Definition 22. Given two positive integers $l, j$, let $W_{l,j}$ be the observability Gramian in time $T$ associated with $(E^l(jT + \cdot), F^l(jT + \cdot), G^l(jT + \cdot))$. Applying the uniform observability assumption to each $(E^i(jT + \cdot), F^l(jT + \cdot), G^l(jT + \cdot))$ and taking into account Eq. (18), we get for $J \in \mathbb{N}$ and $l$ large enough,

$$\gamma_2(T) \geq \frac{\|y^i\|_{T_j}^2}{\|u\|_2^2} \geq \frac{\int_0^T \|G^i(t)x^i(t)\|^2 dt}{\|u\|_2^2} = \sum_{j=0}^{J-1} \|G^i(t)x^i(t+jT)\|^2 dt$$

$$\geq \sum_{j=0}^{J-1} \|x^i(t+jT)\|^2 W_{l,j} \|x^i(t+jT)\|^2 dt$$

$$\geq \gamma \sum_{j=0}^{J-1} \|x^i(t+jT)\|^2 \|u\|_2^2 \geq \gamma J K_3^2 \rho(T_A)^{2l}.$$ 

This implies that the $L_2$-gain $\gamma_2(T)$ is infinite.

Let now $\rho(T_A) > 1$. Let $\bar{t}, \bar{x}$ and $\bar{u}$ be as above. Since the semigroup $\Phi(G_{A,i})$ is quasi-Barabanov, we can find for every $l \in \mathbb{N}$ a time $t_l \geq l$ and a switching law $(\bar{E}^l, \bar{F}^l, \bar{G}^l) \in G_{t_l}$ such that

$$\|\Phi_{\bar{E}^l}(t_l, 0)\bar{x}\| \geq C_{\text{qu}} \rho(T_A)^{t_l} \|\bar{x}\| \geq C_{\text{qu}} \rho(T_A)^l \|\bar{x}\|.$$ 

(19)

According to Theorem 2 and because of the observability counterpart of Proposition 20, there exist $\tilde{s} > 0$ and a switching law $(A_\tilde{s}, B_\tilde{s}, C_\tilde{s}) \in G_{\tilde{s}}$ such that the observability Gramian $W_\tilde{s}(\tilde{s})$ in time $\tilde{s}$ associated with $x(t) = A_\tilde{s}(t)x(t) + B_\tilde{s}(t)\bar{u}(t)$, $y(t) = C_\tilde{s}(t)x(t)$ is invertible. Up to a suitable extension on $(\tilde{s}, \infty)$, we can assume that $(A_\tilde{s}, B_\tilde{s}, C_\tilde{s})$ belongs to $G_{\tilde{s}}$. For $l \geq 0$, consider the sequence of switching signals

$$S^l = (A^0 \ast \bar{E}^l \ast A_\tilde{s}, B^0 \ast \bar{F}^l \ast B_\tilde{s}, C^0 \ast \bar{G}^l \ast C_\tilde{s}) \in G_{\tilde{s}}.$$ 

Denote by $x^l$ the corresponding trajectory of Eq. (13) starting at the origin and associated with the control equal to $\bar{u}$ on $[0, \bar{t}]$ and zero for $t > \bar{t}$ and let $y^l$ be the corresponding output. Note that $x^l(t + t_l) = \Phi_{E^l}(t_l, 0)\bar{x}$ in $[\bar{t}, \bar{t}]$. It then follows from (19) that there exists a positive constant $K_4$ independent of $l$ so that

$$\gamma_2(T) \geq \|y^l\|_{T_j}^2 \geq \frac{\int_{\bar{t} + t_l}^{\bar{t} + t_l + \tilde{s}} \|C_\tilde{s}(s)x^l(s)\|^2 ds}{\|u\|_2^2} = \frac{\int_{\bar{t} + t_l}^{\bar{t} + t_l + \tilde{s}} W_\tilde{s}(\tilde{s})x^l(t + t_l)}{\|u\|_2^2} \geq K_4 \rho(T_A)^{2l},$$

and the right-hand side clearly tends to infinity as $l$ tends to infinity.

Remark 26. Note that the value of $\rho(T_A^\text{min})$ does not depend on the particular choice of the minimal realization thanks to Eq. (15).

Under the hypotheses of Proposition 5, the classes of switching signals considered in Section II-B together with the corresponding families of Section III-A satisfy all the hypotheses of Theorem 25. As a consequence, we can now answer some of the questions raised by Hespanha in [1].

Theorem 27. Let $\mathcal{M}$ be a bounded subset of $M_{n,m}(\mathbb{R}) \times M_{m,n}(\mathbb{R}) \times M_{p,n}(\mathbb{R})$ with $n, m, p$ positive integers and $\tau \geq 0$. Consider the switched linear control system $\dot{x}(t) = A(t)x(t) + B(t)\bar{u}(t)$, $y(t) = C(t)x(t)$, where the switching signal $(A, B, C)$ belongs to the class $S^{d,\tau}(\mathcal{M})$ of piecewise constant signals with dwell-time $\tau$. Let $\gamma_2(\tau)$ be the $L_2$-gain associated with $S^{d,\tau}(\mathcal{M})$. Consider a minimal realization defined on
\[ \mathbb{R}^n, \ k' \leq n, \text{ given by } \dot{x}^{\min}(t) = A^{\min}(t)x(t) \]
\[ + B^{\min}(t)u(t) \text{ with output } y^{\min}(t) = C^{\min}(t)x^{\min}(t) \text{ and associated with a class } S^{d, r}(\mathcal{M}^{\min}) \text{ where } \mathcal{M}^{\min} \text{ is a bounded subset of } M_n(\mathbb{R}) \times M_{n', m}(\mathbb{R}) \times M_{p, n'}(\mathbb{R}). \text{ Assume furthermore that this minimal realization is uniformly observable.} \]

Then, \( \gamma_2(\tau) \) is finite if and only if \( \rho(S^{d, r}(\mathcal{M}^{\min})) < 1 \) and, if \( \tau_{\min} \) is defined as
\[ \tau_{\min} = \inf\{ \tau > 0 \mid \text{ only if } \rho(S^{d, r}(\mathcal{M}^{A^{\min}})) < 1 \}. \]

one has the following characterization: \( \tau_{\min} = \inf\{ \tau > 0 \mid \rho(S^{d, r}(\mathcal{M}^{A^{\min}})) < 1 \} \).

**Remark 28.** The theorem still holds true if we replace the uniformly observable assumption by the hypothesis that there exists at most one \( \tau > 0 \) such that \( \rho(S^{d, r}(\mathcal{M}^{A^{\min}})) = 1 \).

**Remark 29.** One can derive results similar to the previous theorem when one considers variations of certain parameters used in the definition of classes other than \( S^{d, r}(\mathcal{M}) \). For instance, one can characterize the set of \( \mu \in (0, T) \) such that \( \gamma_2(\nu_S^{\text{spe}, T, \mu}(\mathcal{M})) \) is finite in terms of the value of \( \rho(S^{\text{spe}, T, \mu}(\mathcal{M}^{A^{\min}})) \).

**Remark 30.** Recall that the computation of \( t_{\min} = \inf\{ \tau > 0 \mid \rho(S^{d, r}(\mathcal{M}^{A^{\min}})) < 1 \} \) turns out to be a numerically tractable task. Indeed, [10] proposes an LMI procedure providing a sequence of upper bounds of \( t_{\min} \) approximating it arbitrarily well. Therefore, combining Theorem 27 and the LMI-based algorithm of [10] yields a numerical procedure for estimating \( t_{\min} \).

Theorem 25 shows that, under the assumption of uniform observability of a minimal realization, the necessary and sufficient condition for finiteness of the \( L_2 \)-gain (i.e., generalized spectral radius smaller than one) is exactly the same as in the switched framework. We prove below by means of an example that this is no more the case when the assumption of uniform observability does not hold. For this purpose, we next define a switched linear control system satisfying all the assumptions of Theorem 25 (with \( T_A = T_A^{\min} = (G^{\min})_{A^{\bullet}} \) and for which uniform observability does not hold. We have the following example.

**Example 31.** Assume that \( \mathcal{M}(\alpha) = S^{\text{unb}}(\mathcal{M}(\alpha)) \) where \( \mathcal{M}(\alpha) = \{ \mathcal{M}(i, b, c, d) \}_{i, b, c, d} \in M_3(\mathbb{R}) \times \mathbb{R}^3 \times \mathbb{R}^3 \), with
\[
A_1 = \begin{pmatrix}
-1 & -\alpha & 0 \\
\alpha & -1 & 0 \\
0 & 0 & -1
\end{pmatrix}, \quad A_2 = \begin{pmatrix}
-1 & -\alpha & 0 \\
1/\alpha & -1 & 0 \\
0 & 0 & -1
\end{pmatrix},
\]
\[
A_3 = \begin{pmatrix}
-4 & 0 & 1 \\
-4 & 0 & -1 \\
0 & 1 & -1
\end{pmatrix},
\]
\[
b_1 = b_2 = 0, \quad c_1 = c_2 = c_3 = (0, 0, 1)^T,
\]
for \( \alpha > 0 \). We use \( \gamma_2(\alpha) \) to denote the \( L_2 \)-gain induced by the switched linear control system given by \( \dot{x} = A(t)x + b(t)u(t), \ y(t) = c^T(t)x(t) \) and \( (A, b, c) \in \mathcal{M}(\alpha) \). Then we claim that \( \mathcal{T}(\alpha) = \mathcal{T}_A^{\min}(\alpha) \) for every \( \alpha > 0 \) and there exists \( \alpha_* \) (approximately equal to 4.5047) such that \( \rho(T_A(\alpha_*)) = 1 \) and \( \gamma_2(\alpha_*) \leq 4 \).

Before providing a proof, let us note that the assumption of uniform observability does not hold since the observability Gramian in any positive time associated with a switching signal only activating the first two modes contains \( b_3 \) in its kernel.

Let us now prove the claim. Using the results in [18], one determines the value \( \alpha = \alpha_* \approx 4.5047 \) for which the switched system associated with \( \mathcal{M}(\alpha) = \{ A_1, A_2 \} \) is marginally stable (and reducible). In this case, starting from every point \( (x_1, x_2, 0) \), there exists a closed (periodic) \( C^1 \) trajectory \( \Gamma_{x_1, x_2} \) of the switched system lying on the plane \( x_3 = 0 \), which can be completely determined by explicit computations. In particular, we can pick such a trajectory so that its support \( \Gamma \) is contained in the set \( \{ (x_1, x_2, 0) \in \mathbb{R}^3 \mid 1 \leq x_1^2 + x_2^2 \leq 3 \} \). We define the norm \( v(x_1, x_2, 0) \) on the plane \( x_3 = 0 \) by setting \( v^{-1}(1) = \Gamma \). Then \( v \) is a Barabanov norm for the restriction of \( \mathcal{M}_A(\alpha_*) \) on the plane \( x_3 = 0 \). We extend \( v \) to a function on \( \mathbb{R}^3 \), still denoted by \( v \), by setting \( v(x_1, x_2, x_3) = v(x_1, x_2, 0) \) and it follows by explicit computations that \( \| Vv(x) \| \leq \sqrt{3} \) and, by homogeneity, that \( v(x) = \sqrt{v(x^T(x_1, x_2, 0))^T} \) for every \( x \in \mathbb{R}^3 \). Notice, moreover, that
\[ v(x) \geq \sqrt{\frac{x_1^2 + x_2^2}{3}}, \forall x \in \mathbb{R}^3. \]

Let us consider the positive definite function \( V(x) = \frac{1}{2}v(x)^2 + \frac{1}{2}x_3^2 \) and observe that \( \frac{d}{dt} V(x(t)) \leq -x_3^2(t) \) whenever \( A(t) = A_1 \) or \( A(t) = A_2 \). If \( A(t) = A_3 \), one deduces from the above properties of \( v \) that
\[ \frac{d}{dt} V(x(t)) = -x_3^2(t)/4 + |u(t)x_3(t)| \text{ along any trajectory of the switched linear control system. By integrating the previous inequality, using the fact that } x(0) = 0 \text{ and applying Cauchy–Schwarz inequality, we get}
\]
y(t) = C(t)x(t), where the switching law (A, B, C) belongs to the class $\mathcal{S}^{d}$. of piecwise constant signals with dwell-time $\tau \geq 0$. Let $\gamma_2(\tau)$ be the $L_2$-gain associated with $\mathcal{S}^{d}$. We can now state our result on the right-continuity of $\tau \mapsto \gamma_2(\tau)$, which answers Questions (i) and—partially—(iii) of [1].

**Proposition 32.** The function $\gamma_2 : [0, \infty) \to [0, \infty]$ is right-continuous, i.e., for every $\bar{\tau} \in [0, \infty)$, $\lim_{\tau \to \bar{\tau}} \gamma_2(\tau) = \gamma_2(\bar{\tau})$.

**Proof:** For every $T > 0$ and every $\tau \in [0, \infty)$, define the $L_2$-gain in time $T$ as

$$\gamma_2(\tau, T) := \sup \left\{ \frac{\|y_{u, \sigma}2,T\|}{\|u\|_{2,T}} \mid u \in L_2([0, T], \mathbb{R}^m), \{0\}, \sigma \in \mathcal{S}^{d} \right\}.$$ 

It is immediate to see that $\gamma_2(\tau, T)$ is finite for every $(\tau, T) \in [0, \infty) \times (0, \infty)$ and that the maps $\tau \mapsto \gamma_2(\tau, T)$ (for fixed $T > 0$) and $T \mapsto \gamma_2(\tau, T)$ (for fixed $\tau \geq 0$) are non-increasing and non-decreasing respectively. Also notice that $\tau \mapsto \gamma_2(\tau)$ is non-increasing.

We claim that

(i) $\lim_{T \to \infty} \gamma_2(\tau, T) = \gamma_2(\tau)$ for every $\tau \geq 0$;
(ii) the map $\tau \mapsto \gamma_2(\tau, T)$ is right-continuous for every $T > 0$.

In order to prove property (i) of the claim, notice that, given $\tau \geq 0$, $0 < T \leq \infty$, any switching signal $\sigma$ and nonzero control $u \in L_2$, one has

$$\lim_{T \to \infty} \frac{\|y_{u, \sigma}2,T\|}{\|u\|_{2,T}} = \frac{\|y_{u, \sigma}2,T\|}{\|u\|_{2,T}},$$

since $\frac{\|y_{u, \sigma}2,T\|}{\|u\|_{2,T}}$ and $\|u\|_{2,T}$ converge to $\frac{\|y_{u, \sigma}2,T\|}{\|u\|_{2,T}}$ and $\|u\|_{2,T}$ respectively. Property (ii) then follows from the definition of $\gamma_2(\tau)$ and the monotonicity of the $T \mapsto \gamma_2(\tau, T)$.

Let us now prove point (ii) of the claim. With $0 \leq \tau < \tau'$, $T > 0$, and $\sigma \in \mathcal{S}^{d}$, we associate $\sigma' \in \mathcal{L}_\tau([0, \infty), \mathcal{M})$ as follows: $\sigma'(\cdot) := \sigma(\xi)$ where $\xi$ is the largest number in $[0, 1]$ such that the restriction of $\sigma''$ to $[0, T]$ has dwell-time $\tau'$. Notice that $\xi$ is larger than or equal to $\min(\tau, \tau')/\tau'$, where $\tau$ is the largest number such that the restriction of $\sigma$ to $[0, T]$ has dwell-time $\tau$. In particular, $\xi$ is positive and converges to 1 as $\tau' \searrow \tau$. Note that $\sigma$ and $\sigma'$ are equal except on a set of measure upper bounded by $C(T, \tau)(\tau' - \tau)$, where $C(T, \tau)$ denotes some positive constant only depending on $T$ and $\tau$. As a consequence, for any $u \in L_2(0, T)$,

$$\lim_{\tau' \searrow \tau} \frac{\|y_{u, \sigma'}2,T\|}{\|u\|_{2,T}} = \frac{\|y_{u, \sigma}2,T\|}{\|u\|_{2,T}}.$$ 

One immediately deduces property (ii) from the definition of $\gamma_2(\tau, T)$ and the monotonicity of $\tau' \mapsto \gamma_2(\tau', T)$.

We can then conclude the proof of the proposition as follows: given $\tau \geq 0$ and $\kappa < \gamma_2(\tau)$, let $T > 0$ be such that $\kappa < \gamma_2(\tau, T)$ ($T$ exists by property (i)). Take then a right-neighborhood $[\tau, \tau + \eta)$ of $\tau$ such that $\kappa < \gamma_2(\tau', T)$ for every $\tau' \in [\tau, \tau + \eta)$ (the existence of $\eta > 0$ follows by property (ii)). We deduce from the monotonicity of each function $T \mapsto \gamma_2(\tau', T)$ that $\kappa < \gamma_2(\tau')$ for every $\tau' \in [\tau, \tau + \eta)$.
the closed interval $I = \{ \tau \geq 0 | \rho(\tau) = 1 \}$. Notice, however, that $I$ may be nontrivial, as it is the case for the system considered in Example 31, for which $\tau_{\text{min}}$ coincides with the left endpoint of $I$, contrarily to the uniformly observable case, where it is located at its right endpoint. It would be interesting to understand the exact location of $\tau_{\text{min}}$ within $I$ in the general not uniformly observable case.

Another challenging open problem consists in extending the results of [3], [4], which provide an algorithmic approach based on optimal control for the computation of the $L_2$-gain, in the abstract framework introduced here.

REFERENCES


