# Switched nonlinear differential algebraic equations: Solution theory, Lyapunov functions, and stability ${ }^{\text {an }}$ 

Daniel Liberzon ${ }^{\text {a }}$, Stephan Trenn ${ }^{\text {b,* }}$<br>${ }^{a}$ Coordinated Science Laboratory, University of Illinois at Urbana-Champaign, Urbana-Champaign, IL, USA<br>${ }^{b}$ Institute of Mathematics, University of Würzburg, Würzburg, Germany


#### Abstract

We study switched nonlinear differential algebraic equations (DAEs) with respect to existence and nature of solutions as well as stability. We utilize piecewise-smooth distributions introduced in earlier work for linear switched DAEs to establish a solution framework for switched nonlinear DAEs. In particular, we allow induced jumps in the solutions. To study stability, we first generalize Lyapunov's direct method to non-switched DAEs and afterwards obtain Lyapunov criteria for asymptotic stability of switched DAEs. Developing appropriate generalizations of the concepts of a common Lyapunov function and multiple Lyapunov functions for DAEs, we derive sufficient conditions for asymptotic stability under arbitrary switching and under sufficiently slow average dwell-time switching, respectively.


Keywords: Nonlinear differential algebraic equations, piecewise-smooth distributions, Lyapunov functions, asymptotic stability

## 1. Introduction

We consider switched nonlinear differential algebraic equations (DAEs) of the form

$$
\begin{equation*}
E_{\sigma(t)}(x(t)) \dot{x}(t)=f_{\sigma(t)}(x(t)) \tag{1}
\end{equation*}
$$

where $\sigma: \mathbb{R} \rightarrow\{1, \ldots, \mathrm{p}\}, \mathrm{p} \in \mathbb{N}$, is the switching signal and $E_{p}: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n \times n}, f_{p}: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}, p \in\{1, \ldots, \mathrm{p}\}$, are smooth functions. In particular, we assume that each subsystem is a DAE in quasi-linear form (Reich, 1990)

$$
\begin{equation*}
E(x) \dot{x}=f(x) \tag{2}
\end{equation*}
$$

Equations of this kind occur for example when modeling (nonlinear) electrical circuits (Chua and Rohrer, 1965) or coupled mechanical systems (Schiehlen, 1990). Classical linear DAEs (i.e. without switching) of the form $E \dot{x}=A x$, with matrices $E, A \in \mathbb{R}^{n \times n}$, which are also known as singular systems or descriptor systems, naturally appear when modeling electrical circuits because Kirchhoff's circuit laws add algebraic equations to the differential equations stemming from capacitors and inductances. For more details and further motivation for studying (non-switched) DAEs the reader is referred to Kunkel and Mehrmann (2006). Adding, for example, (ideal) switches to an electrical circuit or allowing for sudden structural changes in mechanical systems yield a switched DAE as in (1). When studying the zero dynamics of an ordinary differential equation

[^0](ODE) one arrives at a DAE because of the additional algebraic constraint $0=y=h(x)$, where $h: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ is the output function. In particular, using a switched controller to stabilize the zero dynamics (as was done in Nešić et al., 1999) yields a switched DAE (1) even if one starts with an ODE.

The switching signal in (1) is time-dependent and not state-dependent. Although state-dependent switching has high relevance for applications, we focus our attention in this paper only on time-dependent switching. Some reasons for this are the following: 1) We view the switching signal as an exogenous signal, which is a natural approach for studying electrical circuits with (physical) switches or sudden component faults in electrical and mechanical systems, 2) The distributional solution framework utilized in this paper does not allow for accumulation of switching times (Zeno behavior) which in general can occur for statedependent switching.

The aim of this paper is a stability analysis of (1) with the help of Lyapunov functions. For this we first need to establish a Lyapunov theory for non-switched DAEs in quasi-linear form (2) and secondly we need to define a suitable solution framework for the switched DAE (1).

The use of Lyapunov functions is a powerful tool to study stability of nonlinear differential equations. However, it is not immediately clear how Lyapunov functions can be defined for implicit differential equation such as (2). Of course, it is possible to define a Lyapunov function in a very general setting just by the property that it decreases along solutions; but we believe that only a definition for a Lyapunov function which does not refer to the individual solutions makes Lyapunov functions so useful. The main
problem is that, given a function $x \mapsto V(x)$, its derivative along solutions $\dot{V}(x)=\nabla V(x) \dot{x}$ can not be expressed directly in terms of the right-hand side $f(x)$, because $\dot{x}$ is not explicitly given. We resolve this problem and generalize the well known Lyapunov's Direct Method to implicit differential equations of the form (2). In the linear case $E \dot{x}=A x$ there have been generalizations of Lyapunov's Direct Method (e.g. in Owens and Debeljkovic, 1985; Takaba et al., 1995) but no general definition of a Lyapunov function was given.

One major problem of studying switched DAEs of the form (1) is the presence of jumps in the solutions induced by the presence of so-called consistency spaces. A special case is the problem of inconsistent initial values which has been studied extensively (see e.g. Verghese et al., 1981; Cobb, 1982; Liu et al., 1995; Frasca et al., 2010) and the references in the latter. We are using the piecewise-smooth distributional framework from (Trenn, 2009a,b) to define solutions of the switched DAE (1). In this framework $\dot{x}$ is well defined even when $x$ contains jumps, in which case $\dot{x}$ contains Dirac impulses. It should be noted that a general distributional solution framework (i.e. not considering the smaller space of piecewise-smooth distributions) will not work, because 1) the nonlinear function evaluations $E(x)$ and $f(x)$ are not defined for distributions and 2) the product $E(x) \dot{x}$ is not defined even when $E(x)$ is a piecewisesmooth function.

All results presented here apply of course also to the linear switched DAE

$$
\begin{equation*}
E_{\sigma} \dot{x}=A_{\sigma} x \tag{3}
\end{equation*}
$$

where $E_{p}, A_{p} \in \mathbb{R}^{n \times n}$ for $p \in\{1, \ldots, \mathrm{p}\}$. In this case some of the results simplify significantly and we will formulate corollaries to highlight the linear case. We have studied stability of the linear switched DAE (3) already in Liberzon and Trenn (2009). However, our nonlinear results presented here applied to the linear switched DAE (3) still generalize these results. In particular, the notion of a Lyapunov function as well as the dwell-time stability results are significantly generalized.

Although the two research fields "DAEs" and "switched systems" are now relatively mature (see e.g. the textbooks Kunkel and Mehrmann, 2006; Liberzon, 2003) the combination of both has not been studied much even in the linear case. The existing literature available on switched DAEs (Geerts and Schumacher, 1996a,b; Meng, 2006; Meng and Zhang, 2006; Wunderlich, 2008; Zhai et al., 2006; Raouf and Michalska, 2010) does not consider stability problems in a nonlinear setup. Furthermore, the fundamental problem that one needs distributional solutions for a switched linear DAE (3) and at the same time the equation (3) cannot be evaluated for distributional $x$ is not resolved there.

It might be possible to reformulate the switched DAE (1) as a hybrid system in the framework of (Goebel et al., 2009) by writing (1) as $\dot{x} \in E_{\sigma}(x)^{-1}\left(f_{\sigma}(x)\right)$; however, by doing so, we lose the special structure of (1). In particular,
the jumps of the states are implicitly given by (1) and no additional jump map needs to be considered. This is a major difference between switched DAEs and switched ODEs with reset maps.

A system class which has some similarities with switched DAEs (1) is that of complementarity systems (see e.g. van der Schaft and Schumacher, 1996; Heemels et al., 2000; Çamlıbel et al., 2003; Acary et al., 2008). The main similarity is the existence of different modes which are described by differential-algebraic equations. Roughly speaking, the different modes in the complementarity framework stem directly from the complementarity condition (certain variables must be zero) and a mode change is triggered by violation of positivity of certain variables. In particular, the switches between the different modes are statedependent; hence the solution theory is rather different. Another difference of the complementarity framework is the existence of two different types of variables: the state variable (whose derivative appears explicitly in the system description) and complementarity variables which have to fulfill the complementarity conditions. This distinction is not made in our approach: In one mode a certain statevariable could be governed by a differential equation, in another mode this variable could be governed by a simple algebraic equation. A further comparison of the linear switched DAE (3) with the linear complementarity framework from Heemels et al. (2000) reveals that the consistency projectors are (modulo a restriction to the state variable) identical in both frameworks but the different modes in Heemels et al. (2000) have the same $E$-matrix which simplifies the analysis significantly.

The structure of the paper is as follows. In Section 2 we study the non-switched DAE (2) and generalize Lyapunov's Direct Method to the DAE case in Theorem 2.7. This result is based on a presumably new definition of a Lyapunov function for the DAE (2) as formulated in Definition 2.5. In Section 3 the distributional solution framework for switched DAEs of the form (1) is introduced. We formulate Assumption A4 which under certain regularity assumptions on the subsystems guarantees existence and uniqueness of solutions of the switched DAE (1), see Theorem 3.3. In Section 3.2 we consider the linear case and observe with Corollary 3.9 that the linear equivalent of Assumption A4 ensures the existence of impulse-free solutions of the linear switched DAE (3). Finally, in Section 4 we generalize the well-known results that the existence of a "common Lyapunov function" implies asymptotic stability under arbitrary switching; the novel element is that this Lyapunov function must take into account the consistency projectors as formulated in Theorem 4.1. We also prove a result on stability under average dwell time in the spirit of Hespanha and Morse (1999) for switched nonlinear DAEs (1) in Theorem 4.2.

The following notation is used throughout the paper. $\mathbb{N}, \mathbb{Z}, \mathbb{R}, \mathbb{C}$ are the natural numbers, integers, real and complex numbers, respectively. For a matrix $M \in \mathbb{R}^{n \times m}$, $n, m \in \mathbb{N}$, the kernel (null space) of $M$ is ker $M$, the im-
age (range, column space) of $M$ is $\operatorname{im} M$, and the transpose of $M$ is $M^{\top} \in \mathbb{R}^{m \times n}$. For a matrix $M \in \mathbb{R}^{n \times n}$ and a set $\mathcal{S} \subset \mathbb{R}^{n}$, the image of $\mathcal{S}$ under $M$ is $M \mathcal{S}:=$ $\left\{M x \in \mathbb{R}^{n} \mid x \in \mathcal{S}\right\}$ and the pre-image of $\mathcal{S}$ under $M$ is $M^{-1} \mathcal{S}:=\left\{x \in \mathbb{R}^{n} \mid \exists y \in \mathcal{S}: M x=y\right\}$. The identity matrix is denoted by $I$. For a piecewise-continuous function $f: \mathbb{R} \rightarrow \mathbb{R}$ the left-sided evaluation $\lim _{\varepsilon \searrow 0} f(t-\varepsilon)$ at $t \in \mathbb{R}$ is denoted by $f(t-)$. The space of differentiable functions $f: \mathbb{R} \rightarrow \mathbb{R}$ is denoted by $\mathcal{C}^{1}$, the space of piecewisesmooth functions is denoted by $\mathcal{C}_{\mathrm{pw}}^{\infty}$, the space of distributions is denoted by $\mathbb{D}$ and the space of piecewise-smooth distributions is denoted by $\mathbb{D}_{\mathrm{pw} \mathcal{C}} \infty$; for detailed definitions of these spaces see the Appendix. The set of switching signals considered here is

$$
\Sigma:=\left\{\begin{array}{l|l}
\sigma: \mathbb{R} \rightarrow\{1, \ldots, \mathrm{p}\} & \begin{array}{l}
\sigma \text { is right continuous with a } \\
\text { locally finite number of jumps }
\end{array}
\end{array}\right\}
$$

where $p \in \mathbb{N}$ is the number of subsystems.

## 2. Non-switched DAEs and Lyapunov functions

### 2.1. Classical solutions and consistency spaces

Consider for now the (non-switched) nonlinear DAE (2). A (classical, local) solution of (2) is any differentiable function $x: \mathfrak{I} \rightarrow \mathbb{R}^{n}$ which fulfills (2) on some interval $\mathfrak{I} \subseteq \mathbb{R}$. Due to the time-invariant nature of (2) we can always assume that $\mathfrak{I}=[0, T)$ for some $T \in(0, \infty]$.
Definition 2.1 (Consistency space). The consistency space of (2) is given by

$$
\mathfrak{C}_{E, f}:=\left\{\begin{array}{l|l}
x_{0} \in \mathbb{R}^{n} & \begin{array}{l}
\exists \operatorname{solution} x:[0, T) \rightarrow \mathbb{R}^{n} \\
\text { with } x(0)=x_{0}
\end{array}
\end{array}\right\}
$$

Each $x_{0} \in \mathfrak{C}_{E, f}$ is called a consistent initial value.
Time-invariance of (2) implies that all solutions $x$ of (2) evolve within $\mathfrak{C}_{E, f}$, i.e. $x(t) \in \mathfrak{C}_{E, f}$ for all $t \in[0, T)$. In general, it is not easy to characterize the solution behavior of (2) (for details see e.g. Reich, 1990; Rabier and Rheinboldt, 1994; Schlacher and Zehetleitner, 2004). Here we just assume that the solution behavior is not drastically different from the regular linear case:

Assumption. The nonlinear DAE (2) satisfies:
A1 $f(0)=0$, in particular $0 \in \mathfrak{C}_{E, f}$,
A2 $\mathfrak{C}_{E, f}$ is a closed manifold (possibly with boundary) in $\mathbb{R}^{n}$.

A3 For each $x_{0} \in \mathfrak{C}_{E, f}$ there exists a unique solution $x:[0, \infty) \rightarrow \mathbb{R}^{n}$ with $x(0)=x^{0}$ and $x \in\left(\mathcal{C}^{1} \cap \mathcal{C}_{\mathrm{pw}}^{\infty}\right)^{n}$.

Remark 2.2 (On A3). First note that we exclude systems which exhibit finite escape time. Secondly, the assumption that the differentiable solution is also piecewisesmooth is just a technical assumption which will be needed later for studying switched DAEs.

For the linear case

$$
\begin{equation*}
E \dot{x}=A x \tag{4}
\end{equation*}
$$

with $E, A \in \mathbb{R}^{n \times n}$ Assumptions A1 and A2 are fulfilled trivially (by linearity the consistency space is a linear subspace, see also the forthcoming Theorem 2.3), and A3 is fulfilled if and only if the matrix pair $(E, A)$ is regular, i.e. the polynomial $\operatorname{det}(E s-A) \in \mathbb{R}[s]$ is not the zero polynomial (for details see e.g. the textbook Kunkel and Mehrmann, 2006). Furthermore, regularity of the matrix pair $(E, A)$ is equivalent to the existence of invertible matrices $S, T \in \mathbb{R}^{n \times n}$ such that a coordinate transformation of the codomain and domain by $S$ and $T$ yields the quasiWeierstrass form

$$
(S E T, S A T)=\left(\left[\begin{array}{ll}
I & 0  \tag{5}\\
0 & N
\end{array}\right],\left[\begin{array}{ll}
J & 0 \\
0 & I
\end{array}\right]\right),
$$

where $J \in \mathbb{R}^{n_{1} \times n_{1}}, n_{1} \in \mathbb{N}$, is some matrix and $N \in$ $\mathbb{R}^{n_{2} \times n_{2}}, n_{2}=n-n_{1}$, is nilpotent, i.e. $N^{n_{2}}=0$. We call (5) quasi-Weierstrass form (following Berger, Ilchmann, and Trenn, 2010) because we do not assume that $J$ and $N$ are in Jordan canonical form as is the case for the Weierstrass canonical form (Weierstraß, 1868; Gantmacher, 1959). The smallest number $\nu \in \mathbb{N}$ such that $N^{\nu}=0$ is called the index of the corresponding linear DAE $E \dot{x}=A x$. It is not difficult to see that the consistency space $\mathfrak{C}_{E, A}$ is spanned by the first $n_{1}$ columns of $T$. A convenient way to calculate the matrices $S$ and $T$ is the usage of the Wong sequences of subspaces (named after Wong, 1974) ${ }^{1}$.

Theorem 2.3 (Armentano (1986) ${ }^{2}$ ). Consider a regular matrix pair $(E, A)$ with index $\nu$ and define the associated Wong sequences by, $i \in \mathbb{N}$,

$$
\begin{array}{cc}
\mathcal{V}_{0}:=\mathbb{R}^{n}, & \mathcal{V}_{i+1}:=A^{-1}\left(E \mathcal{V}_{i}\right),
\end{array} \quad \mathcal{V}^{*}:=\bigcap_{i} \mathcal{V}_{i}, ~\left(\mathcal{W}_{0}:=\{0\}, \quad \mathcal{W}_{i+1}:=E^{-1}\left(A \mathcal{W}_{i}\right), \quad \mathcal{W}^{*}:=\mathcal{W}_{i} .\right.
$$

The Wong sequences are nested and get stationary after exactly $\nu$ steps. For any full rank matrices $V, W$ with $\operatorname{im} V=\mathcal{V}^{*}=\mathcal{V}_{\nu}$ and $\operatorname{im} W=\mathcal{W}^{*}=\mathcal{W}_{\nu}$ the matrices $T:=(V, W)$ and $S:=(E V, A W)^{-1}$ are invertible and put $(E, A)$ into the quasi-Weierstrass form (5). In particular

$$
\mathfrak{C}_{E, A}=\mathcal{V}^{*}
$$

Remark 2.4 (Linear index-one case). From the quasiWeierstrass form (5) it can be deduced that the (classical) solutions of (4) do not depend on $N$, or, in other

[^1]words, the solutions remain the same when $N$ is set to be the zero matrix. Assuming that $N$ is the zero matrix is by definition equivalent to assuming that the matrix pair $(E, A)$ is index-one. The importance of $N$ only shows up when studying switched DAEs, where a non-zero $N$ might produce impulses in the solutions (we will study impulse-free solutions in more detail in Section 3.2). An easy way to exclude impulsive behaviors is an index-one assumption for all subsystems, i.e. assuming that in each quasi-Weierstrass form (5) the nilpotent matrix is the zero matrix. However this assumption excludes a large class of interesting switched DAEs. For example, if all subsystems have the same consistency space, then all solutions of the corresponding switched systems will have neither jumps nor impulses, independently of whether or not the subsystems are index-one. In Section 3 we propose Assumption A4, whose linear equivalent (13) ensures impulse-free solutions and is implied by the above two stricter conditions (index-one or same consistency spaces).

### 2.2. Stability and Lyapunov functions

We call the DAE (2) asymptotically stable when all solutions converge to zero as $t \rightarrow \infty$ and for every $\varepsilon>0$ there exists a $\delta>0$ such that for each consistent initial value $x^{0} \in \mathfrak{C}_{E, f}$ with $\left|x^{0}\right|<\delta$ the corresponding solution $x:[0, \infty) \rightarrow \mathfrak{C}_{E, f}$ fulfills $|x(t)|<\varepsilon$ for all $t \geq 0$. The only difference with the classical definition of asymptotic stability is the restriction to consistent initial values. Later, in the switched case, we have to reconsider this restriction, because due to the switching it is not guaranteed that the initial value at a switching instant is consistent.

Definition 2.5 (Lyapunov function). Consider the DAE (2) satisfying A1-A3. Any continuously differentiable nonnegative function $V: \mathfrak{C}_{E, f} \rightarrow \mathbb{R}_{\geq 0}$ fulfilling the following properties is called Lyapunov function for (2):

L1 $V$ is positive definite, i.e. $V(x)=0 \Leftrightarrow x=0$, and for all $x \in \mathfrak{C}_{E, f}$ each sublevel set $V^{-1}[0, V(x)] \subseteq \mathfrak{C}_{E, f}$ is bounded (hence compact by A2),
$\mathbf{L} 2$ there exists a continuous $F: \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow \mathbb{R}$ such that $\nabla V(x) z=F(x, E(x) z)$ for all $x \in \mathfrak{C}_{E, f}, z \in T_{x} \mathfrak{C}_{E, f}$, where $T_{x} \mathfrak{C}_{E, f}$ is the tangent space of $\mathfrak{C}_{E, f}$ at $x$,

L3 defining $\dot{V}(x):=F(x, f(x))$ we have $\dot{V}(x)<0$ for all $x \in \mathfrak{C}_{E, f} \backslash\{0\}$.
Note that in the linear case (4) the tangent space $T_{x} \mathfrak{C}_{E, A}$ is identical to the consistency space $\mathfrak{C}_{E, A}$ for all $x \in \mathfrak{C}_{E, A}$, hence L2 simplifies in this case. Furthermore, for any non-trivial solution $x$ of (2) with a Lyapunov function $V$ it holds that

$$
\begin{array}{r}
\frac{\mathrm{d}}{\mathrm{~d} t} V(x(t))=\nabla V(x(t)) \dot{x}(t) \stackrel{\mathrm{L} 2}{=} F(x(t), E(x(t)) \dot{x}(t)) \\
=F(x(t), f(x(t))=\dot{V}(x(t)) \stackrel{\mathbf{L} \mathbf{3}}{<} 0 \tag{6}
\end{array}
$$

Remark 2.6 (Weaker version of L2). In $\mathbf{L} 2$ one could also work with the weaker assumption $\nabla V(x) z \leq F(x, E(x) z)$ instead of $\nabla V(x) z=F(x, E(x) z)$. However, the definition in $\mathbf{L} \mathbf{3}$ for $\dot{V}$ would then be misleading, because $F(x, f(x))$ would only be an upper bound of $V$. In order to keep the spirit of the classical concept of a Lyapunov function we chose to use $\mathbf{L} 2$ but all results here hold true also for the weaker version. Furthermore, $\mathbf{L} 2$ could be formulated with $z \in E(x)^{-1}(f(x))$ instead of $z \in T_{x} \mathfrak{C}_{E, f}$ because $z$ is a placeholder for $\dot{x}$ when applied later and therefore all relevant $z$ must be solutions of $E(x) z=f(x)$. Depending on the specific problem it might be easier or more difficult to characterize $E(x)^{-1}(f(x))$ instead of $T_{x} \mathfrak{C}_{E, f}$.

Theorem 2.7 (Lyapunov's direct method). Consider the DAE (2) satisfying A1-A3. If there exists a Lyapunov function for (2) then (2) is (globally) asymptotically stable.

## Proof. Stability

For $\varepsilon>0$ consider the set $B_{\varepsilon}:=\left\{x \in \mathfrak{C}_{E, f}| | x \mid=\varepsilon\right\}$ which is empty or compact by Assumption A2. If $B_{\varepsilon}=\emptyset$ then each solution starting within the set enclosed by $B_{\varepsilon}$ cannot leave this set, hence stability follows in this case. Otherwise, let $b:=\min _{x \in B_{\varepsilon}} V(x)$ where positive definiteness of $V$ implies $b>0$. Continuity of $V$ and $V(0)=0$ guarantees the existence of $\delta>0$ such that $V(x)<b$ for all $|x|<\delta$, in particular $\delta<\varepsilon$. From (6) it follows that $t \mapsto V(x(t))$ is decreasing for any solution $x$ of (2), hence any solution $x$ with $|x(0)|<\delta$ fulfills $V(x(t))<b$ for all $t \geq 0$. Seeking a contradiction, assume there exists $t>0$ such that $|x(t)| \geq \varepsilon$, then, by continuity of $x$ together with $|x(0)|<\delta<\varepsilon$, there exists $t_{1} \in(0, t)$ such that $\left|x\left(t_{1}\right)\right|=\varepsilon$ which leads to $b \leq V\left(x\left(t_{1}\right)\right)<b$.
Convergence to zero
Step 1: $V(x(t)) \rightarrow 0$ as $t \rightarrow \infty$.
Let $x:[0, \infty) \rightarrow \mathfrak{C}_{E, f}$ be any non-trivial solution, then the non-negative function $t \mapsto v(t):=V(x(t)) \geq 0$ is strictly decreasing by (6). Therefore, $\underline{v}=\lim _{t \rightarrow \infty} v(t)$ is well defined. Seeking a contradiction, assume $\underline{v}>0$. Then $v(t) \in[\underline{v}, v(0)]$ for all $t \geq 0$. By $\mathbf{L} 1$ and continuity of $V, \mathcal{K}:=V^{-1}[\underline{v}, v(0)]$ is a compact set, hence $\mathcal{M}:=\dot{V}(\mathcal{K}) \subseteq \mathbb{R}$ is also compact (since $\dot{V}$ is continuous) and $0 \notin \mathcal{M}$. This implies that $m:=-\max \mathcal{M}>0$ and, in particular, $v^{\prime}(t)=\frac{\mathrm{d}}{\mathrm{d} t} V(x(t))=\dot{V}(x(t)) \leq-m<0$ for all $t \geq 0$. Hence $v(t) \leq v(0)-m t$ for all $t \geq 0$, which contradicts $v(t) \geq 0$ for all $t \geq 0$, hence $\underline{v}=0$ must hold. Step 2: $x(t) \rightarrow 0$ as $t \rightarrow \infty$.
Seeking a contradiction, assume $x(t) \nrightarrow 0$, then there exists a sequence $\left(t_{n}\right)_{n \in \mathbb{N}}$ with $t_{n} \rightarrow \infty$ as $n \rightarrow \infty$ and $\varepsilon>0$ such that $\left|x\left(t_{n}\right)\right|>\varepsilon$. By L1 and (6), each solution $x$ evolves within the compact set $V^{-1}[0, V(x(0))]$, hence there exists a convergent subsequence of $x\left(t_{n}\right)$ with limit $x^{*} \neq 0$. By continuity and positive definiteness of $V$ we arrive at the contradiction $0=\lim _{t \rightarrow \infty} V(x(t))=V\left(x^{*}\right)>0 . \quad$ qed

Remark 2.8 (The linear case). In the linear, regular
hence $V$ is decreasing along solutions.
case it is well-known (Owens and Debeljkovic, 1985) ${ }^{3}$ that $E \dot{x}=A x$ is asymptotically stable if, and only if, there exists a solution $(P, Q) \in \mathbb{R}^{n \times n} \times \mathbb{R}^{n \times n}$ of the generalized Lyapunov equation

$$
\begin{equation*}
A^{\top} P E+E^{\top} P A=-Q \tag{7}
\end{equation*}
$$

where $P=P^{\top}$ is positive definite and $Q=Q^{\top}$ is positive definite on $\mathfrak{C}_{E, A}$. In fact, it is easy to see that then $V(x)=(E x)^{\top} P E x$ is a Lyapunov function in the sense of Definition 2.5 with

$$
\nabla V(x) z=(E x)^{\top} P E z+(E z)^{\top} P E x=: F(x, E z)
$$

and

$$
\dot{V}(x)=x^{\top}\left(E^{\top} P A+A^{\top} P E\right) x=-x^{\top} Q x<0 \text { on } \mathfrak{C}_{E, A} .
$$

If the linear system $E \dot{x}=A x$ is index-one, i.e. $N=0$ in the quasi-Weierstrass form (5), it is shown in (Takaba et al., 1995; Ishihara and Terra, 2002) that asymptotic stability is also equivalent to the existence of a solution $P \in \mathbb{R}^{n \times n}$ of

$$
P^{\top} A+A^{\top} P=-Q, \quad P^{\top} E=E^{\top} P \geq 0
$$

for any positive definite $Q \in \mathbb{R}^{n \times n}$. The corresponding "asymmetric" Lyapunov function ${ }^{4}$ is given by $V(x)=$ $(E x)^{\top} P x$, with $\nabla V(x) z=(E x)^{\top} P z+(E z)^{\top} P x=x^{\top} P^{\top} E z+$ $(E z)^{\top} P x=: F(x, E z)$ and $\dot{V}(x)=x^{\top}\left(P^{\top} A+A^{\top} P\right) x=$ $-x^{\top} Q x<0$.

We conclude this section with an example which illustrates the application of Theorem 2.7.
Example 2.9. Consider the nonlinear DAE

$$
\left[\begin{array}{ccc}
\sin x_{3} & \cos x_{3} & 0  \tag{8}\\
0 & 0 & 0 \\
0 & 0 & 0
\end{array}\right]\left(\begin{array}{l}
\dot{x}_{1} \\
\dot{x}_{2} \\
\dot{x}_{3}
\end{array}\right)=\left(\begin{array}{c}
-x_{1} \sin x_{3}-x_{2} \cos x_{3} \\
x_{1} \cos x_{3}-x_{2} \sin x_{3} \\
x_{3}-x_{1}^{2}-x_{2}^{2}
\end{array}\right)
$$

which fulfills our Assumptions A1, A2 and A3. The consistency space is given by the equation $x_{3}=x_{1}^{2}+x_{2}^{2}$ and $x_{1} \cos x_{3}=x_{2} \sin x_{3}$; the projection to the $x_{1}-x_{2}$-plane is illustrated in Figure 1. Note that the consistency space can be parameterized by

$$
\mathfrak{C}_{E, f}=\left\{\left(\theta \sin \theta^{2}, \theta \cos \theta^{2}, \theta^{2}\right)^{\top} \mid \theta \in \mathbb{R}\right\}
$$

The corresponding tangent space is given by, for $x \neq 0$,

$$
\begin{equation*}
T_{x} \mathfrak{C}_{E, f}=\operatorname{span}\left\{\left(x_{1}+2 x_{2} x_{3}, x_{2}-2 x_{1} x_{3}, 2 x_{3}\right)^{\top}\right\} \tag{9}
\end{equation*}
$$

and $T_{0} \mathfrak{C}_{E, f}=\operatorname{span}\left\{(0,1,0)^{\top}\right\}$. We propose the following Lyapunov function candidate:

$$
V(x)=x_{3}
$$

[^2]

Figure 1: Consistency space of Example 2.9 in the $x_{1}-x_{2}$-plane (left) and in the $x_{2}-x_{3}$-plane (right); the dynamics within the consistency space are shown by the arrows.

For all $x \in \mathfrak{C}_{E, f}$ it follows that $x_{3}=x_{1}^{2}+x_{2}^{2}$, hence $V$ fulfills L1. Aiming for a function $F: \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow \mathbb{R}$ satisfying L2, i.e. for $x \in \mathfrak{C}_{E, f}$ and $z \in T_{x} \mathfrak{C}_{E, f}$,

$$
\begin{equation*}
F(x, E(x) z)=z_{3}=\nabla V(x) z \tag{10}
\end{equation*}
$$

we choose, for $x \in \mathfrak{C}_{E, f}$ and $v \in E(x)^{-1}\left(T_{x} \mathfrak{C}_{E, f}\right)$,

$$
F(x, v):=\frac{2 x_{3} v_{1}}{x_{1} \sin x_{3}+x_{2} \cos x_{3}}
$$

Then by using (9) as well as $x_{1} \cos x_{3}=x_{2} \sin x_{3}$ we indeed obtain (10). Finally,

$$
\begin{aligned}
\dot{V}(x)=F(x, f(x))= & F\left(\left(\begin{array}{l}
x_{1} \\
x_{2} \\
x_{3}
\end{array}\right),\left(\begin{array}{c}
-x_{1} \sin x_{3}-x_{2} \cos x_{3} \\
x_{1} \cos x_{3}-x_{2} \sin x_{3} \\
x_{3}-x_{1}^{2}-x_{2}^{2}
\end{array}\right)\right) \\
& =\frac{2 x_{3}\left(-x_{2} \sin x_{3}-x_{2} \cos x_{3}\right)}{x_{1} \sin x_{3}+x_{2} \cos x_{3}}=-2 x_{3}
\end{aligned}
$$

hence $\mathbf{L} \mathbf{3}$ is fulfilled and $V$ is a Lyapunov function for (8) and Theorem 2.7 shows that (8) is globally asymptotically stable.

## 3. Solutions of switched DAEs

### 3.1. The general nonlinear case

Recall the switched nonlinear DAE (1)

$$
E_{\sigma}(x) \dot{x}=f_{\sigma}(x)
$$

where each subsystem $E_{p}(x) \dot{x}=f_{p}(x), p=\{1, \ldots, \mathrm{p}\}$, fulfills Assumptions A1-A3 and $\sigma \in \Sigma$ is the switching signal. As an underlying solution framework for (1), we will use the space $\mathbb{D}_{\mathrm{pw} \mathcal{C}}$ of piecewise smooth distributions which was introduced in (Trenn, 2009a,b) for studying linear switched DAE. For a short summary of the basic definition and the main properties of piecewise-smooth distributions see the Appendix.

Definition 3.1 (Solution of (1)). A solution of (1) on some interval $\mathfrak{I} \subseteq \mathbb{R}$ is any piecewise-smooth function $x \in$ $\left(\mathcal{C}_{\mathrm{pw}}^{\infty}\right)^{n}$ such that (1) restricted to $\mathfrak{I}$ holds as an equation of piecewise-smooth distributions, i.e.

$$
\left(E_{\sigma}(x)\left(x_{\mathbb{D}}\right)^{\prime}\right)_{\mathfrak{J}}=\left(f_{\sigma}(x)_{\mathbb{D}}\right)_{\mathfrak{J}}
$$

The product $E(x)\left(x_{\mathbb{D}}\right)^{\prime}$ in Definition 3.1 is well defined, since by assumption $t \mapsto E(x(t))$ is piecewise smooth and $\left(x_{\mathbb{D}}\right)^{\prime}$ is a piecewise-smooth distribution. Note that this definition of a solution does not allow for Dirac impulses in the solution. There are two reasons for this: 1) It is not clear how a nonlinear function of a Dirac impulse should be defined in general and 2) for a stability analysis the existence of Dirac impulses in the solution can be interpreted as an undesired unstable solution. However, in Section 3.2 we will also study solutions with impulses for linear switched DAEs. The following assumption is essential for existence and uniqueness of solutions of the switched DAE (1).

Assumption. The switched DAE (1) and the corresponding consistency spaces $\mathfrak{C}_{p}:=\mathfrak{C}_{E_{p}, f_{p}}, p \in\{1, \ldots, \mathrm{p}\}$, satisfy:

A4 $\forall p, q \in\{1, \ldots, \mathrm{p}\} \forall x_{0}^{-} \in \mathfrak{C}_{p} \exists$ unique $x_{0}^{+} \in \mathfrak{C}_{q}$ : $x_{0}^{+}-x_{0}^{-} \in \operatorname{ker} E_{q}\left(x_{0}^{+}\right)$.

Assumption A4 makes it possible to define nonlinear consistency projectors $\Pi_{q}, q \in\{1, \ldots, \mathrm{p}\}$ :

$$
\Pi_{q}: \bigcup_{p} \mathfrak{C}_{p} \rightarrow \mathfrak{C}_{q}, \quad x_{0}^{-} \mapsto x_{0}^{+}
$$

where $x_{0}^{+}$is the unique value given by Assumption A4. In particular, $\Pi_{q}(x)=x$ for all $x \in \mathfrak{C}_{q}$.

Remark 3.2 (Motivation of Assumption A4). For a motivation of Assumption A4 consider the situation where the system switches from subsystem $p \in\{1, \cdots, \mathrm{p}\}$ to subsystem $q \in\{1, \ldots, \mathrm{p}\}$ at some switching time $t \in \mathbb{R}$. Any solution $x$ (in the sense of Def. 3.1) fulfills $x(t-) \in \mathfrak{C}_{p}$ and $x(t+) \in \mathfrak{C}_{q}$ and the impulsive part of $\dot{x}$ at $t$ is given by $\dot{x}[t]=(x(t+)-x(t-)) \delta_{t}$, where $\delta_{t}$ denotes the Dirac impulse at $t$. Since the right-hand side $f_{\sigma}(x)_{\mathbb{D}}$ does not contain impulses it follows that $E_{\sigma}(x) \dot{x}[t]=0$ must hold. This directly implies the existence part of Assumption A4. Furthermore, uniqueness of $x$ (for a given past) follows when $x(t+)$ is uniquely given by $x(t-)$, hence Assumption A4 is a necessary condition for existence and uniqueness of a solution $x$ of (1) in the sense of Def. 3.1. In the ODE case, i.e. $E_{\sigma}(x) \equiv I$, Assumption A4 is trivially fulfilled with $x_{0}^{+}:=x_{0}^{-}$. In the linear case an easy check for Assumption A4 is possible, see Section 3.2.

The following Theorem shows that Assumption A4 is also sufficient for existence and uniqueness of solutions of (1).

Theorem 3.3 (Existence and uniqueness). Consider the switched nonlinear DAE (1) satisfying $\boldsymbol{A} 4$ and $\boldsymbol{A} 1$ - $\boldsymbol{A} 3$ for each subsystem. Then for every switching signal $\sigma \in \Sigma$ and every $x_{0} \in \mathfrak{C}_{\sigma(0-)}$ there exists a unique solution $x \in\left(\mathcal{C}_{\mathrm{pw}}^{\infty}\right)^{n}$ of $(1)$ on $[0, \infty)$ with $x(0-)=x_{0}$. Furthermore, for all $t \in[0, \infty)$ and all solutions $x$ of (1),

$$
x(t)=\Pi_{\sigma(t)}(x(t-)),
$$

where $\Pi_{p}, p \in\{1, \ldots, \mathrm{p}\}$, are the consistency projectors induced by A4. In particular, on each interval which does not contain a switching time, $x$ is a classical solution of the corresponding subsystem.

Proof. Step 1: Existence of a solution.
Let $t_{0}=0$ and $t_{i}>0, i=1,2, \ldots$ be the ordered switching times of $\sigma$ after $t_{0}$ and let $p_{i}:=\sigma\left(t_{i}\right)$. Inductively and invoking Assumption A3 choose $x^{i} \in\left(\mathcal{C}^{1} \cap \mathcal{C}_{\mathrm{pw}}^{\infty}\right)^{n}$, $i \in \mathbb{N}$, such that $x^{i}$ is the unique (classical) solution of $E_{p_{i}}\left(x^{i}\right) \dot{x}^{i}=f_{p_{i}}\left(x^{i}\right)$ on the interval $\left[t_{i}, t_{i+1}\right)$ with $x^{i}\left(t_{i}\right)=$ $\Pi_{p_{i}}\left(x^{i-1}\left(t_{i}-\right)\right)$, where $x^{-1}\left(t_{0}-\right):=x_{0}$. We show that any $x \in\left(\mathcal{C}_{\mathrm{pw}}^{\infty}\right)^{n}$ with $x(0-)=x_{0}$ and $x_{\left[t_{i}, t_{i+1}\right)}=x^{i}{ }_{\left[t_{i}, t_{i+1}\right)}$ for $i \in \mathbb{N}$ solves the switched DAE (1) on $[0, \infty)$. By definition $x$ solves (1) on each open interval $\left(t_{i}, t_{i+1}\right)$ and it remains to check that

$$
\left(E_{\sigma}(x)\left(x_{\mathbb{D}}\right)^{\prime}\right)\left[t_{i}\right]=\left(f_{\sigma}(x)_{\mathbb{D}}\right)\left[t_{i}\right]=0 \quad \text { for all } i \in \mathbb{N}
$$

where $D[t]$ denotes the impulsive part of $D \in\left(\mathbb{D}_{\mathrm{pw} \mathcal{C}^{\infty}}\right)^{n}$ at $t \in \mathbb{R}$ (see Appendix for details). Invoking the properties of piecewise-smooth distributions, it follows that

$$
\begin{aligned}
& \left(E_{\sigma}(x)\left(x_{\mathbb{D}}\right)^{\prime}\right)\left[t_{i}\right]=E_{p_{i}}\left(x\left(t_{i}\right)\right)\left(x\left(t_{i}\right)-x\left(t_{i}-\right)\right) \delta_{t_{i}} \\
& \quad=E_{p_{i}}\left(\Pi_{p_{i}}\left(x\left(t_{i}-\right)\right)\right)\left(\Pi_{p_{i}}\left(x\left(t_{i}-\right)\right)-x\left(t_{i}-\right)\right) \delta_{t_{i}}=0
\end{aligned}
$$

where the last equation follows from Assumption A4. Hence $x$ is a solution of $(1)$ on $[0, \infty)$.
Step 2: Uniqueness of the solution.
With the notation as in Step 1 it suffices to show that the solution $x$ as constructed above is unique on $\left[0, t_{1}\right)$, uniqueness on $\left[t_{1}, \infty\right)$ follows then inductively. Let $z \in\left(\mathcal{C}_{\mathrm{pw}}^{\infty}\right)^{n}$ be a solution of $(1)$ on $\left[0, t_{1}\right)$ with $z(0-)=x_{0}$. With a similar argument as in Step 1 it follows that

$$
E_{p_{0}}(z(0))\left(z(0)-x_{0}\right)=0
$$

hence Assumption A4 ensures $z(0)=\Pi_{p_{0}}\left(x_{0}\right)=x(0)$. Furthermore, Assumption A4 also implies that $z(t)=$ $z(t-)$ for all $t \in\left(0, t_{1}\right)$, hence $z$ is continuous on $\left(0, t_{1}\right)$ which together with Assumption A3 implies that $z_{\left(0, t_{1}\right)}=$ $x_{\left(0, t_{1}\right)}$. Hence uniqueness of the solution is shown. qed
Remark 3.4 (Assumption A4 for a single system). Note that Assumption A4 applied to each single system, i.e. $p=q$, additionally restricts the possible nonlinear DAEs even without switching: In A4 one can always choose $x_{0}^{+}=x_{0}^{-}$if $p=q$ and the asserted uniqueness of $x_{0}^{+}$implies therefore

$$
\begin{equation*}
\forall x_{0}^{+} \in \mathfrak{C}_{p}: \operatorname{ker} E_{p}\left(x_{0}^{+}\right) \cap\left\{x_{0}^{+}-x_{0}^{-} \mid x_{0}^{-} \in \mathfrak{C}_{p}\right\}=\{0\} \tag{11}
\end{equation*}
$$

So in addition to A1-A3 each subsystem must also fulfill (11). In the linear case it can be shown that A3 already implies (11), but in the general case this is not true as the following example shows:

$$
\begin{aligned}
x_{2} \dot{x}_{1} & =0 \\
\dot{x}_{2} & =1
\end{aligned}
$$

With initial value $x_{2}(0)=t_{0} \in \mathbb{R}$, we get the unique solution $x_{2}(t)=t+t_{0}$ and $\dot{x}_{1}(t)=0$ for all $t \neq-t^{0}$. The only classical solution of the latter is $x_{1}(t)=x_{1}^{0}$, where $x_{1}^{0} \in \mathbb{R}$ is the initial value for $x_{1}$. Hence A3 holds. However, condition A4 is not fulfilled because (11) does not hold. In fact, when allowing jumps in solutions (as in the case for switched DAEs) uniqueness of solutions is lost, because $x$ can have an arbitrary jump at $t=t_{0}$ without violating the DAE (in a distributional sense).

Remark 3.5 (Index-one systems). If the nonlinear DAE (2) can be written as (e.g. via a (nonlinear) coordinate transformation)

$$
\begin{aligned}
\dot{x}_{1} & =g\left(x_{1}, x_{2}\right) \\
0 & =h\left(x_{1}, x_{2}\right)
\end{aligned}
$$

where $h$ is such that $x_{2}$ can be solved in terms of $x_{1}$ then (in analogy to the linear case) (2) is said to have index one. In this case, Assumption A4 clearly holds with consistency projector $\binom{x_{1}}{x_{2}} \mapsto\binom{x_{1}}{\bar{h}\left(x_{1}\right)}$, where the function $\bar{h}$ is such that $x_{2}=\bar{h}\left(x_{1}\right)$ is the unique solution of $h\left(x_{1}, x_{2}\right)=0$. However, Assumption A4 is weaker than the index-one assumption because it could hold even when not all subsystems are index-one; see also Remark 2.4.

### 3.2. The linear case

Consider the linear switched DAE (3) with switching signal $\sigma \in \Sigma$. As already mentioned above, the Assumptions A1-A3 for each subsystem reduce to the regularity condition $\operatorname{det}\left(E_{p} s-A_{p}\right) \not \equiv 0$ for each subsystem. Under this assumption (in particular without assuming A4) it already follows from (Trenn, 2009a,b) that existence and uniqueness of solutions of (3) is guaranteed. However, these solutions are then elements of the space of piecewisesmooth distributions and will therefore, in general, contain Dirac impulses and their derivatives. The following example illustrates this phenomenon.

Example 3.6. Consider (3) with subsystems given by

$$
\begin{array}{ll}
\left(E_{1}, A_{1}\right)= & \left(E_{2}, A_{2}\right)= \\
\left(\left[\begin{array}{lll}
1 & 0 & 0 \\
0 & 1 & 0 \\
0 & 0 & 1
\end{array}\right],\left[\begin{array}{lll}
0 & 0 & 0 \\
0 & 0 & 0 \\
0 & 0 & 0
\end{array}\right]\right), & \left(\left[\begin{array}{lll}
0 & 0 & 0 \\
1 & 0 & 0 \\
0 & 1 & 0
\end{array}\right],\left[\begin{array}{lll}
1 & 0 & 0 \\
0 & 1 & 0 \\
0 & 0 & 1
\end{array}\right]\right) .
\end{array}
$$

Then the switching signal $\sigma(t)=1$ on $[0,1)$ and $\sigma(t)=2$ on $[1, \infty)$ together with the initial condition $x(0)=(1,0,0)$ enforces a jump to zero in $x_{1}$ at the switching time. At the switching time the second system is already active, in particular $x_{2}=\dot{x}_{1}$ holds, hence $x_{2}$ is the derivative of a jump, i.e. $x_{2}=-\delta_{1}$ where $\delta_{t} \in \mathbb{D}_{\mathrm{pw} \mathcal{C}^{\infty}}$ denotes the Dirac impulse at time $t$. Furthermore, also the equation $x_{3}=\dot{x}_{2}$ must hold which yields that $x_{3}=-\delta_{1}^{\prime}$, i.e. $x_{3}$ contains the derivative of the Dirac impulse.

Since the presence of impulses in solutions can be seen as an undesired unstable behavior (see the next section),
we would like to give an easily checkable condition which ensures that for arbitrary switching all solutions of (3) are impulse-free (but may still exhibit jumps). It will turn out that this condition is equivalent to Assumption A4 but is easier to check in the linear case. For the formulation of this condition, we define the linear consistency projector of a regular matrix pair $(E, A)$.

Definition 3.7 (Linear consistency projector). Consider a regular matrix pair $(E, A) \in \mathbb{R}^{n \times n} \times \mathbb{R}^{n \times n}$ and, invoking Theorem 2.3, choose invertible matrices $S, T \in \mathbb{R}^{n \times n}$ such that ( $S E T, S A T$ ) is in quasi-Weierstrass form (5) with $n_{1} \times n_{1}$ and $n_{2} \times n_{2}$ the corresponding diagonal block sizes. The linear consistency projector is then given by

$$
\Pi_{E, A}:=T\left[\begin{array}{ll}
I & 0 \\
0 & 0
\end{array}\right] T^{-1}
$$

where $I$ is an $n_{1} \times n_{1}$ identity matrix.
Let $\mathcal{V}^{*}$ and $\mathcal{W}^{*}$ be the limits of the Wong sequences as in Theorem 2.3. Then it is easy to see that the definition of $\Pi_{E, A}$ is independent of the choice of $T$ and that it is a linear projection onto $\mathcal{V}^{*}=\mathfrak{C}_{E, A}$ along $\mathcal{W}^{*}$, i.e.

$$
\begin{equation*}
\Pi_{E, A}^{2}=\Pi_{E, A}, \operatorname{im} \Pi_{E, A}=\mathcal{V}^{*}, \operatorname{ker} \Pi_{E, A}=\mathcal{W}^{*} \tag{12}
\end{equation*}
$$

With the help of the linear consistency projectors it is now possible to give an easily checkable characterization of Assumption A4.

Theorem 3.8 (Linear version of Assumption A4). Consider the switched linear DAE (3) with regular matrix pairs $\left(E_{p}, A_{p}\right)$ and corresponding consistency projectors $\Pi_{p}, p \in\{1, \ldots, \mathrm{p}\}$ as in Definition 3.7. Then Assumption $\boldsymbol{A} 4$ is equivalent to

$$
\begin{equation*}
\forall p, q \in\{1, \ldots, \mathrm{p}\}: \quad E_{q}\left(\Pi_{q}-I\right) \Pi_{p}=0 \tag{13}
\end{equation*}
$$

and the linear mapping $x_{0}^{-} \mapsto x_{0}^{+}:=\Pi_{q} x_{0}^{-}$coincides with the consistency projector associated with Assumption A4.

Proof. Let $p, q \in\{1, \ldots, \mathrm{p}\}$ and $x_{0}^{-} \in \mathfrak{C}_{p}:=\mathfrak{C}_{E_{p}, A_{p}}$ be arbitrary and fixed in the rest of the proof.
Step 1: We show $(13) \Rightarrow \mathbf{A 4}$.
Let $x_{0}^{+}:=\Pi_{q} x_{0}^{-} \in \mathfrak{C}_{q}:=\mathfrak{C}_{E_{q}, A_{q}}$, then, since $\Pi_{p} x_{0}^{-}=x_{0}^{-}$,
$E_{q}\left(x_{0}^{+}-x_{0}^{-}\right)=E_{q}\left(\Pi_{q} \Pi_{p} x_{0}^{-}-\Pi_{p} x_{0}^{-}\right)=E_{q}\left(\Pi_{q}-I\right) \Pi_{p} x_{0}^{-} \stackrel{(13)}{=} 0$,
hence the existence assertion of Assumption A4 is shown. To show uniqueness of $x_{0}^{+}$, let $z \in \mathfrak{C}_{q}$ be such that

$$
z-x_{0}^{-} \in \operatorname{ker} E_{q} \subseteq \mathcal{W}_{q}^{*}=\operatorname{ker} \Pi_{q}
$$

where $\mathcal{W}_{q}^{*}$ is the limit of the corresponding Wong sequence for $\left(E_{q}, A_{q}\right)$ as in Theorem 2.3. Together with $\Pi_{q} z=z$ this implies $z=\Pi_{q} x_{0}^{-}=x_{0}^{+}$.
Step 2: We show A4 $\Rightarrow$ (13).
Choose $x_{0}^{+} \in \mathfrak{C}_{q}$ such that $x_{0}^{+}-x_{0}^{-} \in \operatorname{ker} E_{q} \subseteq \mathcal{W}_{q}^{*}=$ $\operatorname{ker} \Pi_{q}$, hence $x_{0}^{+}=\Pi_{q} x_{0}^{+}=\Pi_{q} x_{0}^{-}$. Therefore, by $\Pi_{p} x_{0}^{-}=$ $x_{0}^{-}$,
$0=E_{q}\left(x_{0}^{+}-x_{0}^{-}\right)=E_{q}\left(\Pi_{q} \Pi_{p} x_{0}^{-}-\Pi_{p} x_{0}^{-}\right)=E_{q}\left(\Pi_{q}-I\right) \Pi_{p} x_{0}^{-}$.

Since $x_{0}^{-} \in \mathfrak{C}_{p}=\mathcal{V}_{p}^{*}$ is arbitrary it follows from $\mathcal{V}_{p}^{*} \oplus \mathcal{W}_{p}^{*}=$ $\mathbb{R}^{n}$ together with $\mathcal{W}_{p}^{*}=\operatorname{ker} \Pi_{p}$ that $E_{q}\left(\Pi_{q}-I\right) \Pi_{p}=0$, hence (13) holds.

Combining Theorems 3.3 and 3.8 yields that for every switched linear DAE (3) with regular matrix pairs $\left(E_{p}, A_{p}\right), p=1, \ldots, \mathrm{p}$, satisfying (13) there exists a solution $x \in\left(\mathcal{C}_{\mathrm{pw}}^{\infty}\right)^{n}$, unique in this class of functions. By definition, this solution also solves (3) in the distributional framework of (Trenn, 2009a,b). Since the switched DAE (3) with regular pairs $\left(E_{p}, A_{p}\right), p=1, \ldots, \mathrm{p}$, has a unique distributional solution (for a fixed initial value $x(0-)$ ) we obtain the following result.

Corollary 3.9 (Impulse-free solutions for (3)).
Consider the switched DAE (3) with arbitrary switching signal $\sigma \in \Sigma$ and regular matrix pairs $\left(E_{p}, A_{p}\right)$ with corresponding consistency projectors $\Pi_{p} \in \mathbb{R}^{n \times n}$ given by Definition 3.7. If (13) holds, then every distributional solution $x \in\left(\mathbb{D}_{\mathrm{pw}} \mathcal{C}^{\infty}\right)^{n}$ of (3) is impulse-free.

## 4. Asymptotic stability of switched DAEs

Asymptotic stability for (1), with a fixed switching signal $\sigma$, can be defined basically in the same way as for the non-switched case, see Section 2.2; the only difference is that the solutions might have jumps, so we have to decide where to evaluate the initial value. In view of Theorem 3.3, we consider the initial value $x(0-)$. Note that in the linear case Assumption A4 excludes impulses in the solution, which is reasonable for the definition of stability, because an impulse can be interpreted as an infinite peak which remains infinite even when the corresponding solution is scaled so that $|x(0-)|$ gets arbitrarily small.

Theorem 4.1 (Arbitrary switching). Consider the switched DAE (1) satisfying Assumption A4 and Assumptions A1-A3 for each subsystem with corresponding consistency spaces $\mathfrak{C}_{p}:=\mathfrak{C}_{E_{p}, f_{p}}$ and consistency projectors $\Pi_{p}, p \in\{1, \ldots, \mathrm{p}\}$ induced by A4. Assume for each subsystem that there exists a Lyapunov function $V_{p}: \mathfrak{C}_{p} \rightarrow$ $\mathbb{R}_{\geq 0}$ in the sense of Definition 2.5. If

$$
\begin{equation*}
\forall p, q \in\{1, \ldots, \mathrm{p}\} \forall x \in \mathfrak{C}_{p}: \quad V_{q}\left(\Pi_{q}(x)\right) \leq V_{p}(x) \tag{14}
\end{equation*}
$$

then the switched $D A E(1)$ is asymptotically stable for any switching signal $\sigma \in \Sigma$.

Proof. Step 1: Definition of a Lyapunov function candidate.
If $x \in \mathfrak{C}_{p} \cap \mathfrak{C}_{q}$ for some $p, q \in\{1, \ldots, \mathrm{p}\}$ then $x=\Pi_{p}(x)=$ $\Pi_{q}(x)$ hence (14) implies $V_{p}(x)=V_{q}(x)$. Therefore

$$
V: \bigcup_{p} \mathfrak{C}_{p} \rightarrow \mathbb{R}, \quad x \mapsto V_{p}(x) \quad \text { for } x \in \mathfrak{C}_{p}
$$

is well defined.
Step 2: $V(x(t)) \rightarrow 0$ as $t \rightarrow \infty$.
Fix $\sigma \in \Sigma$ and let $x:[0, \infty) \rightarrow \mathbb{R}^{n}$ be a solution of (1)
in the sense of Theorem 3.3. Consider an interval $\mathfrak{I} \subseteq \mathbb{R}$ without switching times then $x$ is also a classical (local) solution of $E_{p}(x) \dot{x}=f_{p}(x)$ on $\mathfrak{I}$ where $p:=\sigma(\tau)$ for $\tau \in \mathfrak{I}$. From $x(\tau) \in \mathfrak{C}_{p}$ for all $\tau \in \mathfrak{I}$ it follows that $V(x(\tau))=$ $V_{p}(x(\tau))$ and, by Definition 2.5 together with (6),

$$
\frac{\mathrm{d}}{\mathrm{~d} t} V_{p}(x(\tau))=\dot{V}_{p}(x(\tau))<0 \quad \forall \tau \in \mathfrak{I}
$$

Let $t \in \mathbb{R}$ be a switching time of $\sigma$, then $x(t)=\Pi_{\sigma(t)}(x(t-))$ and $x(t-) \in \mathfrak{C}_{\sigma(t-)}$, yields, invoking (14),

$$
\begin{aligned}
V(x(t)) & =V_{\sigma(t)}(x(t))=V_{\sigma(t)}\left(\Pi_{\sigma(t)}(x(t-))\right) \\
& \leq V_{\sigma(t-)}(x(t-))=V(x(t-))
\end{aligned}
$$

Hence $t \mapsto v(t)=V(x(t))$ is monotonically decreasing and therefore $\underline{v}:=\lim _{t \rightarrow \infty} v(t) \geq 0$ is well defined. Seeking a contradiction, assume $\underline{v}>0$. Analogously to the proof of Theorem 2.7 let $\mathcal{K}_{p}:=V_{p}^{-1}[\underline{v}, v(0)], \mathcal{M}_{p}:=\dot{V}\left(\mathcal{K}_{p}\right)$ and $m_{p}:=-\max \mathcal{M}_{p}>0$. Let $m=\min _{p} m_{p}>0$ then $\frac{\mathrm{d}}{\mathrm{d} t} v(t)<-m<0$ for all non-switching (hence almost all) times $t \geq 0$, which contradicts $v(t) \geq 0$ and the assertion of Step 2 is shown.
Step 3: $x(t) \rightarrow 0$ as $t \rightarrow \infty$.
Seeking a contradiction, assume $x(t) \nrightarrow 0$. Then there exists $\varepsilon>0$ and a sequence $\left(s_{i}\right)_{i \in \mathbb{N}} \in \mathbb{R}^{\mathbb{N}}$ with $s_{i} \rightarrow \infty$ as $i \rightarrow \infty$ such that $\left|x\left(s_{i}\right)\right|>\varepsilon$ for all $i \in \mathbb{N}$. There is at least one $p \in\{1, \ldots, \mathrm{p}\}$ such that the set $\left\{i \in \mathbb{N} \mid \sigma\left(s_{i}\right)=p\right\}$ has infinitely many elements, therefore, without loss of generality, assume that $\sigma\left(s_{i}\right)=p$ for some $p$ and all $i \in$ $\mathbb{N}$. Since each $x\left(s_{i}\right)$ is contained within the compact set $V_{p}^{-1}[0, V(x(0))]$, the same argument as in the proof of Theorem 2.7 shows existence of $x^{*} \neq 0$ such that we arrive at the contradiction $0=\lim _{t \rightarrow \infty} V(x(t))=\lim _{i \rightarrow \infty} V_{p}\left(x\left(s_{i}\right)\right)=$ $V_{p}\left(x^{*}\right) \neq 0$.
Step 4: Stability of the switched DAE.
We first show that for all $\varepsilon>0$ there exists $b_{\varepsilon}>0$ such that for all $p \in\{1, \ldots, \mathrm{p}\}$ and all $x \in \mathfrak{C}_{p}$ :

$$
\begin{equation*}
V_{p}(x)<b_{\varepsilon} \Rightarrow|x|<\varepsilon \tag{15}
\end{equation*}
$$

Assume the contrary, then there exists $\varepsilon>0$ and sequences $\left(p_{n}\right)_{n \in \mathbb{N}}$ and $\left(x_{n}\right)_{n \in \mathbb{N}}$ such that $V_{p_{n}}\left(x_{n}\right)<1 / n$ and $\left|x_{n}\right| \geq \varepsilon$. There exist at least one $p \in\{1, \ldots, \mathrm{p}\}$ which occurs infinitely often in the sequence $\left(p_{n}\right)$, so we can, without loss of generality, assume that $p_{n}=p$ for all $n \in \mathbb{N}$ and some $p \in\{1, \ldots, \mathrm{p}\}$. Then, by L1, all $x_{n}$ are contained in the compact set $V_{p}^{-1}\left[0, V_{p}\left(x_{n_{\max }}\right)\right]$ where $n_{\max }:=\operatorname{argmax}_{n} V_{p}\left(x_{n}\right)<\infty$. This implies that there exists $x^{*} \in \mathfrak{C}_{p}$ which is a limit of a subsequence of $\left(x_{n}\right)$ and with $\left|x^{*}\right| \geq \varepsilon$. Hence we arrive at the contradiction $0=\lim _{n \rightarrow \infty} V_{p}\left(x_{n}\right)=V_{p}\left(x^{*}\right) \neq 0$ and the claim (15) is shown.

For a given $\varepsilon>0$ choose $b_{\varepsilon}>0$ according to (15). Let $p_{0}:=\sigma(0-)$, then by continuity of $V_{p_{0}}$ there exists $\delta>0$ such that $|x|<\delta$ implies $V_{p_{0}}(x)<b_{\varepsilon}$ for all $x \in$ $\mathfrak{C}_{p_{0}}$. In Step 2 it was shown that $t \mapsto V_{\sigma(t-)}(x(t-))$ is monotonically decreasing, hence $V_{\sigma(t-)}(x(t-))<b_{\varepsilon}$ for all $t \geq 0$. Hence (15) yields $|x(t-)|<\varepsilon$ for all $t \geq 0$. qed

Condition (14) implies that any two Lyapunov functions $V_{p}$ and $V_{q}$ coincide on the intersection $\mathfrak{C}_{p} \cap \mathfrak{C}_{q}$, hence Theorem 4.1 is a generalization of the switched ODE case where the existence of a common Lyapunov function is sufficient to ensure stability under arbitrary switching (Liberzon, 2003, Thm. 2.1). However, without condition (14) the existence of a common Lyapunov function is not enough (Liberzon and Trenn, 2009) for asymptotic stability of the switched DAE (1). Under arbitrary switching, solutions will in general exhibit jumps; these jumps are described by the consistency projectors, and these projectors must "fit together" with the Lyapunov functions in the sense of (14) to ensure stability of the switched DAE under arbitrary switching. Finally, with some additional effort it can be shown that the hypotheses of Theorem 4.1 guarantee uniformity of the asymptotic stability with respect to the switching signal.

It is well-known for switched ODEs that by restricting the class of switching signals one can obtain asymptotic stability also in cases where no common Lyapunov function exists. Denote by $N_{\sigma}(t, T)$ the number of switchings of $\sigma$ in the interval $[t, T)$ and define the class of average dwell time switching signals with average dwell time $\tau_{a}>0$ (Hespanha and Morse, 1999)

$$
\Sigma_{\tau_{a}}:=\left\{\begin{array}{l|l}
\sigma \in \Sigma & \begin{array}{l}
\exists N_{0}>0 \forall t \in \mathbb{R} \forall \Delta t>0: \\
N_{\sigma}(t, t+\Delta t)<N_{0}+\frac{\Delta t}{\tau_{a}}
\end{array}
\end{array}\right\} .
$$

The number $N_{0}>0$ is called chatter bound of the switching signal $\sigma \in \Sigma_{\tau_{a}}$. Note that the subset of average dwell time switching signals with chatter bound $N_{0}=1$ is precisely the class of switching signals with dwell time $\tau_{a}$.
Theorem 4.2 (Average dwell time switching). Consider ${ }_{B}$ the switched DAE (1) with corresponding consistency space $\mathfrak{C}_{p}$ and consistency projectors $\Pi_{p}, p \in\{1, \ldots, \mathrm{p}\}$. Assume that all subsystems permit a Lyapunov function $V_{p}$, $p \in\{1, \ldots, \mathrm{p}\}$, which additionally fulfill

1. $\exists \lambda>0: \dot{V}_{p}(x) \leq-\lambda V_{p}(x)$ for all $p \in\{1, \ldots, \mathrm{p}\}$, $x \in \mathfrak{C}_{p}$ and
2. $\exists \mu>0: V_{q}\left(\Pi_{q}(x)\right) \leq \mu V_{p}(x)$ for all $p, q \in\{1, \ldots, \mathrm{p}\}$, $x \in \mathfrak{C}_{p}$.
Then the switched DAE (1) with switching signal $\sigma \in \Sigma_{\tau_{a}}$ is asymptotically stable if

$$
\begin{equation*}
\tau_{a}>\frac{\ln \mu}{\lambda} \tag{16}
\end{equation*}
$$

Proof. With standard arguments (see e.g. (Liberzon, 2003)) it follows that the non-negative function $t \mapsto V_{\sigma(t-)}(x(t-))$ is bounded by an exponentially decreasing function and hence converges to zero. Arguments analogous to those in Step 3 and Step 4 of the proof of Theorem 4.1 now conclude the proof.

In the linear case the Lyapunov function can be chosen according to Remark 2.8; in this case it is possible to express the inequality (16) for the average dwell time directly in terms of eigenvalues of corresponding matrices.

Lemma 4.3 (The linear case). Consider the linear switched $D A E(3)$ with the regular matrix pairs $\left(E_{p}, A_{p}\right), p \in\{1, \ldots, \mathrm{p}\}$ with corresponding consistency spaces $\mathfrak{C}_{p}$, and let $\left(P_{p}, Q_{p}\right)$ be the solutions of the corresponding generalized Lyapunov equation (7). Choose a matrix $O_{p}$ with orthonormal columns such that $\operatorname{im} O_{p}=\operatorname{im} \Pi_{p}=\mathfrak{C}_{p}$, where $\Pi_{p}$ is the linear consistency projector corresponding to $\left(E_{p}, A_{p}\right)$ Then, for $p, q \in\{1, \ldots, \mathrm{p}\}$,

$$
\forall x \in \mathfrak{C}_{p}: \quad V_{q}\left(\Pi_{q} x\right) \leq \mu_{p, q} V_{p}(x)
$$

where $\mu_{p, q}:=\frac{\lambda_{\max }\left(O_{p}^{\top} \Pi_{q}^{\top} E_{q}^{\top} P_{q} E_{q} \Pi_{q} O_{p}\right)}{\lambda_{\min }\left(O_{p}^{\top} E_{p}^{\top} P_{p} E_{p} O_{p}\right)} \geq 0$ and

$$
\forall x \in \mathfrak{C}_{p}: \quad \dot{V}_{p}(x) \leq-\lambda_{p} V_{p}(x)
$$

where $\lambda_{p}:=\frac{\lambda_{\min }\left(O_{p}^{\top} Q_{p} O_{p}\right)}{\lambda_{\max }\left(O_{p}^{\top} E_{p}^{\top} P_{p} E_{p} O_{p}\right)}>0$ and where $\lambda_{\min }(\cdot)$ and $\lambda_{\max }(\cdot)$ denote the minimal and maximal eigenvalue of $a$ symmetric matrix, respectively.

Proof. Let $d_{p}:=\operatorname{dim} \mathfrak{C}_{p}$, i.e. $O_{p} \in \mathbb{R}^{n \times d_{p}}$, then $x \in$ $\mathfrak{C}_{p}$ if, and only if, there exists a unique $z \in \mathbb{R}^{d_{p}}$ with $x=O_{p} z, O_{p}^{\top} x=z$ and $|x|=|z|$. Hence, by choosing $z$ corresponding to $x \in \mathfrak{C}_{p}$ as above,

$$
\begin{aligned}
V_{p}(x) & =z^{\top} O_{p}^{\top} E_{p}^{\top} P_{p} E_{p} O_{p} z=: z^{\top} P_{p}^{z} z \\
& \geq \lambda_{\min }\left(P_{p}^{z}\right)|z|^{2}=\lambda_{\min }\left(P_{p}^{z}\right)|x|^{2} \\
V_{p}(x) & \leq \lambda_{\max }\left(P_{p}^{z}\right)|x|^{2} \\
V_{q}\left(\Pi_{q} x\right) & =z^{\top} O_{p}^{\top} \Pi_{q}^{\top} E_{q}^{\top} P_{q} E_{q} \Pi_{q} O_{p} z=: z^{\top} M_{p, q}^{z} z \\
& \leq \lambda_{\max }\left(M_{p, q}^{z}\right)|x|^{2} \\
\dot{V}_{p}(x) & =-z^{\top} O_{p}^{\top} Q_{p} O_{p} z=:-z^{\top} Q_{p}^{z} z \leq-\lambda_{\min }\left(Q_{p}^{z}\right)|x|^{2}
\end{aligned}
$$

By assumption, the matrices $Q_{p}^{z}=Q_{p}^{z \top} \in \mathbb{C}^{d_{p} \times d_{p}}$ and $P_{p}^{z}=P_{p}^{z^{\top}} \in \mathbb{C}^{d_{p} \times d_{p}}$ are positive definite, hence $\lambda_{\min }\left(Q_{p}^{z}\right)>$ 0 and $\lambda_{\max }\left(P_{p}^{z}\right) \geq \lambda_{\min }\left(P_{p}^{z}\right)>0$. Therefore,

$$
\mu_{p, q}:=\frac{\lambda_{\max }\left(M_{p, q}^{z}\right)}{\lambda_{\min }\left(P_{p}^{z}\right)} \geq 0, \lambda_{p}:=\frac{\lambda_{\min }\left(Q_{p}^{z}\right)}{\lambda_{\max }\left(P_{p}^{z}\right)}>0
$$

are well defined. Note that $\lambda_{\max }\left(M_{p, q}^{z}\right)=0$ is possible, however $\lambda_{\max }\left(M_{p, p}^{z}\right)=\lambda_{\max }\left(P_{p}^{z}\right) \geq \lambda_{\min }\left(P_{p}^{z}\right)$, hence $\mu_{p, p} \geq 1$ and $\max _{p, q} \ln \mu_{p, q} \geq 0$.
Corollary 4.4 (Average dwell time, linear case). For the switched linear $D A E$ (3) with asymptotically stable subsystems, let $\mu_{p, q}$ and $\lambda_{p}, p, q \in\{1, \ldots, \mathrm{p}\}$ be given as in Lemma 4.3. Then the linear switched DAE (3) is asymptotically stable if $\sigma \in \Sigma_{\tau_{a}}$ with

$$
\tau_{a}>\frac{\max _{p, q} \ln \mu_{p, q}}{\min _{p} \lambda_{p}}
$$

Note that the obtained results cannot in general be expressed in terms of the eigenvalues of the matrices $Q_{p}$ and $P_{p}$ (or $E_{p}^{\top} P_{p} E_{p}$ ); the consistency projectors and basis transformation must be incorporated as formulated in Lemma 4.3. We show the application of Corollary 4.4 with a simple linear example, which is based on Example 1 from (Liberzon and Trenn, 2009).

## Example 4.5. Let

$$
\left(E_{1}, A_{1}\right)=\left(\left[\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right],\left[\begin{array}{ll}
0 & -1 \\
1 & -1
\end{array}\right]\right),\left(E_{2}, A_{2}\right)=\left(\left[\begin{array}{ll}
1 & 1 \\
0 & 0
\end{array}\right],\left[\begin{array}{cc}
-1 & -1 \\
1 & 0
\end{array}\right]\right) .
$$

The corresponding consistency spaces and consistency projectors are given by

$$
\mathfrak{C}_{1}:=\mathfrak{C}_{E_{1}, A_{1}}=\operatorname{im}\left[\begin{array}{l}
1 \\
1
\end{array}\right], \quad \mathfrak{C}_{2}:=\mathfrak{C}_{E_{2}, A_{2}}=\operatorname{im}\left[\begin{array}{l}
0 \\
1
\end{array}\right]
$$

and $\Pi_{1}=\left[\begin{array}{ll}0 & 1 \\ 0 & 1\end{array}\right], \Pi_{2}=\left[\begin{array}{ll}0 & 0 \\ 1 & 1\end{array}\right]$. In Liberzon and Trenn (2009) it is shown that the corresponding switched DAE is not asymptotically stable under arbitrary switching. However, we can apply the result of Corollary 4.4. As basis matrices for the consistency space choose $O_{1}=\frac{1}{2}\left[\begin{array}{c}\sqrt{2} \\ \sqrt{2}\end{array}\right], O_{2}=$ $\left[\begin{array}{l}0 \\ 1\end{array}\right]$. Consider the Lyapunov functions $V_{1}(x)=\frac{1}{2} x_{2}^{2}$ and $V_{2}(x)=\frac{1}{2}\left(x_{1}+x_{2}\right)^{2}$, corresponding to $P_{1}=P_{2}=\frac{1}{2}\left[\begin{array}{cc}1 & 0 \\ 0 & 0\end{array}\right]$ and $Q_{1}=\left[\begin{array}{ll}0 & 0 \\ 0 & 1\end{array}\right], Q_{2}=\left[\begin{array}{ll}1 & 1 \\ 1 & 1\end{array}\right]$. Then

$$
\begin{aligned}
O_{1}^{\top} E_{1}^{\top} P_{1} E_{1} O_{1} & =\frac{1}{4}, & O_{2}^{\top} E_{2}^{\top} P_{2} E_{2} O_{2} & =\frac{1}{2} \\
O_{1}^{\top} \Pi_{2}^{\top} E_{2}^{\top} P_{2}^{\top} E_{2} \Pi_{2} O_{1} & =1, & O_{2}^{\top} \Pi_{1}^{\top} E_{1}^{\top} P_{1}^{\top} E_{1} \Pi_{1} O_{2} & =\frac{1}{2} \\
O_{1}^{\top} Q_{1} O_{1} & =\frac{1}{2}, & O_{2}^{\top} Q_{2} O_{2} & =1
\end{aligned}
$$

hence $\mu:=\max _{p, q} \mu_{p, q}=2$ and $\lambda:=\min _{p} \lambda_{p}=2$. Therefore the corresponding switched DAE is asymptotically stable for all switching signals $\sigma \in \Sigma_{\tau_{a}}$ with $\tau_{a}>\frac{\ln 2}{2}$. This bound is actually sharp in this example.

## 5. Conclusion

We have studied switched nonlinear DAEs with respect to solution and stability theory. For the non-switched nonlinear DAE subsystems we generalized the classical Lyapunov's Direct Method, in particular, we defined a Lyapunov function for quasi-linear DAEs in general terms. Furthermore, we studied existence and uniqueness of solutions of a switched nonlinear DAE, provided the subsystems are regular in a certain sense. Finally, we were able to generalize existing stability results on switched ODEs to switched DAEs.

## Appendix: Piecewise smooth distributions

We assume familiarity with the definitions and properties of classical distributions as formalized by Schwartz (Schwartz, 1950, 1951). We denote the space of test functions (i.e., smooth functions $\varphi: \mathbb{R} \rightarrow \mathbb{R}$ with compact support) by $\mathcal{C}_{0}^{\infty}$, then the space of distributions is the dual space of the space of test functions, i.e.

$$
\mathbb{D}:=\left\{D: \mathcal{C}_{0}^{\infty} \rightarrow \mathbb{R} \mid D \text { is linear and continuous }\right\}
$$

The main two properties of distributions are 1) that they can be interpreted as generalized functions and 2) that they are arbitrarily often differentiable. To be more precise, let $\mathcal{L}_{1, \text { loc }}$ be the space of locally integrable functions, then the mapping

$$
\mathcal{L}_{1, \text { loc }} \rightarrow \mathbb{D}, f \mapsto f_{\mathbb{D}}:=\left(\varphi \mapsto \int_{\mathbb{R}} f \varphi\right)
$$

is well defined (i.e. $f_{\mathbb{D}}$ is indeed a distribution) and an injective homomorphism. The simplest distribution which is not induced by a function is the Dirac impulse given by $\delta(\varphi):=\varphi(0)$, or, in general for $t \in \mathbb{R}, \delta_{t}(\varphi):=\varphi(t)$ for $\varphi \in \mathcal{C}_{0}^{\infty}$. The derivative of an arbitrary distribution $D \in \mathbb{D}$ is given by $D^{\prime}(\varphi):=-D\left(\varphi^{\prime}\right)$ for $\varphi \in \mathcal{C}_{0}^{\infty}$. Distributions can be multiplied with smooth functions:

$$
(\alpha D)(\varphi):=D(\alpha \varphi), \quad \alpha \in \mathcal{C}^{\infty}, D \in \mathbb{D}, \varphi \in \mathcal{C}_{0}^{\infty}
$$

Let $\mathcal{C}_{\mathrm{pw}}^{\infty}$ be the space of piecewise-smooth function, where $\alpha: \mathbb{R} \rightarrow \mathbb{R}$ is called piecewise-smooth when there exists a locally finite ordered set $S=\left\{s_{i} \in \mathbb{R} \mid i \in \mathbb{Z}\right\}$ and smooth functions $\alpha_{i} \in \mathcal{C}^{\infty}, i \in \mathbb{Z}$, such that $\alpha=$ $\sum_{i \in \mathbb{Z}}\left(\alpha_{i}\right)_{\left[s_{i}, s_{i+1}\right)}$. Here, $f_{I}$ denotes the restriction (or truncation) of a function $f: \mathbb{R} \rightarrow \mathbb{R}$ to the interval $I \subseteq \mathbb{R}$, i.e. $f_{I}(\tau)=f(\tau)$ for $\tau \in I$ and $f_{I}(\tau)=0$ otherwise. The space of piecewise-smooth distributions is then given by
$\mathbb{D}_{\mathrm{pw} \mathcal{C}^{\infty}}:=\left\{f_{\mathbb{D}}+\sum_{\tau \in T} D_{\tau} \left\lvert\, \begin{array}{l}f \in \mathcal{C}_{\mathrm{pw}}^{\infty}, T \subseteq \mathbb{R} \text { loc. finite } \\ \forall \tau \in T: D_{\tau} \in \operatorname{span}\left\{\delta_{\tau}, \delta_{\tau}^{\prime}, \delta_{\tau}^{\prime \prime}, \ldots\right\}\end{array}\right.\right\}$.
The properties of $\mathbb{D}_{\mathrm{pw}} \mathcal{C}^{\infty}$ and corresponding definitions are summarized in the following, $D=f_{\mathbb{D}}+\sum_{\tau \in T} D_{\tau} \in \mathbb{D}_{\mathrm{pwC}}{ }^{\infty}$ and $t \in \mathbb{R}$ :

1. Closed under differentiation: $D^{\prime} \in \mathbb{D}_{\mathrm{pw}} \mathcal{C}^{\infty}$.
2. Left- and right-evaluation: $D(t+):=f(t), D(t-):=$ $f(t-)$.
3. Impulsive part: $D[t]:=D_{t}$ if $t \in T, D[t]=0$ otherwise.
4. Restriction to interval: $D_{I}:=\left(f_{I}\right)_{\mathbb{D}}+\sum_{\tau \in T \cap I} D_{t}$, where $I \subseteq \mathbb{R}$ is some interval.
5. Multiplication with piecewise-smooth function: $\alpha D:=$ $\sum_{i \in \mathbb{Z}} \alpha_{i} D_{\left[s_{i}, s_{i+1}\right)}$, where $\alpha=\sum_{i \in \mathbb{Z}}\left(\alpha_{i}\right)_{\left[s_{i}, s_{i+1}\right)}$ as above; in particular, $\alpha \delta_{t}=\alpha(t) \delta_{t}$.
For more details see (Trenn, 2009a,b). In the proof of Theorem 3.3 we actually need the fact that for any $\alpha \in$ $\mathcal{C}_{\mathrm{pw}}^{\infty},(\alpha D)[t]=\alpha D[t]$ and

$$
\left(\alpha_{D}\right)^{\prime}=\sum_{i \in \mathbb{Z}}\left(\alpha_{i}^{\prime}\right)_{\left[s_{i}, s_{i+1}\right)}+\sum_{i \in \mathbb{Z}}\left(\alpha_{i}\left(s_{i}\right)-\alpha_{s_{i-1}}\left(t_{i}\right)\right) \delta_{s_{i}},
$$

where $\alpha=\sum_{i \in \mathbb{N}}\left(\alpha_{i}\right)_{\left[s_{i}, s_{i+1}\right)}$ as above.

## References

Acary, V., Brogliato, B., Goeleven, D., 2008. Higher order Moreau's sweeping process: mathematical formulation and numerical simulation. Mathematical Programming 113 (1), 133-217.
Aplevich, J. D., 1991. Implicit Linear Systems. No. 152 in Lecture Notes in Control and Information Sciences. Springer-Verlag, Berlin.
Armentano, V. A., 1986. The pencil $(s E-A)$ and controllabilityobservability for generalized linear systems: a geometric approach. SIAM J. Control Optim. 24, 616-638.
Berger, T., Ilchmann, A., Trenn, S., 2010. The quasi-Weierstraß form for regular matrix pencils. Lin. Alg. Appl.In press, preprint available online, Institute for Mathematics, Ilmenau University of Technology, Preprint Number 09-21.

Çamlıbel, M. K., Heemels, W. P. M. H., van der Schaft, A. J., Schumacher, J. M. H., 2003. Switched networks and complementarity. IEEE Trans. Circuits Syst., I, Fundam. Theory Appl. 50 (8), 1036-1046.
Chua, L. O., Rohrer, R. A., 1965. On the dynamic equations of a class of nonlinear RLC networks. IEEE Trans. Circ. Theory CT12 (4), 475-489.
Cobb, J. D., 1982. On the solution of linear differential equations with singular coefficients. J. Diff. Eqns. 46, 310-323.
Dieudonné, J., 1946. Sur la réduction canonique des couples des matrices. Bull. de la Societé Mathématique de France 74, 130-146.
Frasca, R., Çamlıbel, M. K., Goknar, I. C., Iannelli, L., Vasca, F., December 2010. Linear passive networks with ideal switches: Consistent initial conditions and state discontinuities. IEEE Trans. Circuits Syst., I, Fundam. Theory Appl. 57 (12), 3138-3151.
Gantmacher, F. R., 1959. The Theory of Matrices (Vol. I). Chelsea, New York.
Geerts, A. H. W. T., Schumacher, J. M. H., 1996a. Impulsive-smooth behavior in multimode systems. Part I: State-space and polynomial representations. Automatica 32 (5), 747-758.
Geerts, A. H. W. T., Schumacher, J. M. H., 1996b. Impulsive-smooth behavior in multimode systems. Part II: Minimality and equivalence. Automatica 32 (6), 819-832.
Goebel, R., Sanfelice, R. G., Teel, A. R., 2009. Hybrid dynamical systems. IEEE Control Systems Magazine 29 (2), 28-93.
Heemels, W. P. M. H., Schumacher, J. M. H., Weiland, S., MarchApril 2000. Linear complementarity systems. SIAM J. Appl. Math. 60 (4), 1234-1269.
Hespanha, J. P., Morse, A. S., 1999. Stability of switched systems with average dwell-time. In: Proc. 38th IEEE Conf. Decis. Control. pp. 2655-2660.
Ishihara, J. Y., Terra, M. H., November 2002. On the Lyapunov theorem for singular systems. IEEE Trans. Autom. Control 47 (11), 1926-1930.
Kuijper, M., 1994. First-Order Representations of Linear Systems. Birkhäuser, Boston.
Kunkel, P., Mehrmann, V., 2006. Differential-Algebraic Equations. Analysis and Numerical Solution. EMS Publishing House, Zürich, Switzerland.
Liberzon, D., 2003. Switching in Systems and Control. Systems and Control: Foundations and Applications. Birkhäuser, Boston.
Liberzon, D., Trenn, S., December 2009. On stability of linear switched differential algebraic equations. In: Proc. IEEE 48th Conf. on Decision and Control. pp. 2156-2161.
Liu, W. Q., Yan, Y., Teo, K. L., September 1995. On initial instantaneous jumps of singular systems. IEEE Trans. Autom. Control 40 (9), 1650-1655.
Meng, B., 2006. Observability conditions of switched linear singular systems. In: Proceedings of the 25th Chinese Control Conference. Harbin, Heilongjiang, China, pp. 1032-1037.
Meng, B., Zhang, J.-F., 2006. Reachability conditions for switched linear singular systems. IEEE Trans. Autom. Control 51 (3), 482488.

Nešić, D., Skafidas, E., Mareels, I. M. Y., Evans, R. J., 1999. Minimum phase properties for input nonaffine nonlinear systems. IEEE Trans. Autom. Control 44 (4), 868-872.
Owens, D. H., Debeljkovic, D. L., 1985. Consistency and Liapunov stability of linear descriptor systems: A geometric analysis. IMA J. Math. Control \& Information, 139-151.

Rabier, P. J., Rheinboldt, W. C., 1994. On the computation of impasse points of quasi-linear differential-algebraic equations. Math. Comp. 62 (205), 259-293.
Raouf, J., Michalska, H. H., December 2010. Exponential stabilization of singular systems by controlled switching. In: Proc. 49th IEEE Conf. Decis. Control, Atlanta, USA. IEEE Control Systems Society, IEEE, pp. 414-419.
Reich, S., 1990. On a geometrical interpretation of differentialalgebraic equations. Circuits Systems Signal Process. 9 (4), 367382.

Schiehlen, W. (Ed.), 1990. Multibody Systems Handbook. SpringerVerlag, Heidelberg, Germany.

Schlacher, K., Zehetleitner, K., 2004. Formale Methoden für implizite dynamische Systeme. Automatisierungstechnik 52 (9), 446-455.
Schwartz, L., 1950, 1951. Théorie des Distributions I,II. No. IX,X in Publications de l'institut de mathématique de l'Universite de Strasbourg. Hermann, Paris.
Takaba, K., Morihira, N., Katayama, T., 1995. A generalized Lyapunov theorem for descriptor systems. Syst. Control Lett. 24 (1), 49-51.
Trenn, S., 2009a. Distributional differential algebraic equations. Ph.D. thesis, Institut für Mathematik, Technische Universität Ilmenau, Universitätsverlag Ilmenau, Ilmenau, Germany. URL http://www.db-thueringen.de/servlets/DocumentServlet?id=13581
Trenn, S., 2009b. Regularity of distributional differential algebraic equations. Math. Control Signals Syst. 21 (3), 229-264.
van der Schaft, A. J., Schumacher, J. M. H., 1996. The complementary-slackness class of hybrid systems. Math. Control Signals Syst. 9, 266-301.
Verghese, G. C., Levy, B. C., Kailath, T., August 1981. A generalized state-space for singular systems. IEEE Trans. Autom. Control AC26 (4), 811-831.
Weierstraß, K., 1868. Zur Theorie der bilinearen und quadratischen Formen. Monatsh. Akademie. Wiss., 310-338.
Wong, K.-T., 1974. The eigenvalue problem $\lambda T x+S x$. J. Diff. Eqns. 16, 270-280.
Wunderlich, L., 2008. Analysis and numerical solution of structured and switched differential-algebraic systems. Ph.D. thesis, Fakultät II Mathematik und Naturwissenschaften, Technische Universität Berlin, Berlin, Germany.
Zhai, G., Kou, R., Imae, J., Kobayashi, T., 2006. Stability analysis and design for switched descriptor systems. In: Proceedings of the 2006 IEEE International Symposium on Intelligent Control. pp. 482-487.


[^0]:    ${ }^{i}$ This work was supported by NSF grants CNS-0614993, ECCS0821153 and DFG grant Wi1458/10-1.

    * Corresponding author

    Email addresses: liberzon@illinois.edu (Daniel Liberzon), stephan.trenn@mathematik.uni-wuerzburg.de (Stephan Trenn)

[^1]:    ${ }^{1}$ These sequences can be traced back to Dieudonné (1946); however, he only implicitly considers the second Wong sequence via a duality argument. Although some authors use these sequences (Aplevich, 1991; Kuijper, 1994; van der Schaft and Schumacher, 1996), the connection between them and the quasi-Weierstrass form, as established by Theorem 2.3, seems not to be very well known.
    ${ }^{2}$ See also Berger et al. (2010).

[^2]:    ${ }^{3}$ Actually, in Owens and Debeljkovic (1985) only the complexvalued case is studied; however, by considering the real part of the generalized Lyapunov equation (7) we also obtain real-valued matrix pairs $(P, Q)$ with the desired properties.
    ${ }^{4}$ We thank Emilia Fridman for making us aware of this Lyapunov function construction.

