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ANALYSIS OF THE IMPLICIT EULER TIME-DISCRETIZATION OF PASSIVE LINEAR DESCRIPTOR COMPLEMENTARITY SYSTEMS

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Abstract This article is largely concerned with the time-discretization of descriptor-variable systems coupled with complementarity constraints. They are named linear descriptor complementarity systems (LDCS). More specifically passive LDCS with minimal state space representation are studied. The Euler implicit discretization of LDCS is analysed: the one-step non-smooth problem (OSNSP), that is a generalized equation, is shown to be well-posed under some conditions. Then the convergence of the discretized solutions is studied. Several examples illustrate the applicability and the limitations of the developments.

12 1 INTRODUCTION

The analysis of non-smooth dynamical systems with set-valued right-hand sides satisfying maximal 13 monotone properties has been the object of many contributions, see [26] for a survey. Linear and 14 nonlinear complementarity dynamical systems, various kinds of projected dynamical systems, of 15 differential variational inequalities, differential inclusions with maximal monotone right-hand sides, 16 Moreau's sweeping processes, and some switching dynamics can be recast in this class of set-valued 17 dynamical systems. In this article we focus on singular, or linear descriptor complementarity systems 18 (LDCS). They may be viewed either as an extension of "classical" descriptor-variable systems which 19 are widely used in Automatic Control and Circuit Theory, or of "classical" linear complementarity 20 systems. They can also be seen as differential-algebraic equations (DAEs) coupled with complementarity 21 constraints. Motivations for studying LDCS are in circuits with set-valued components [2, Chapters 3 22 and 5] [32], and in a more abstract framework in state-dependent switching DAEs [71, 70], after a kind 23 of convexification (or filling-in the graph) procedure has been made. They are also closely related to 24 DAEs whose equality constraints are nonsmooth, as they occur in chemistry or thermo-dynamical 25 systems [78, 79, 77]. 26 It has to be noted that the study of LDCS has not received a lot of attention yet, and is still at an 27 embryonic stage of analysis. In particular the analysis of the implicit Euler method, which has received 28 a lot of attention in several fields [10, 14, 57, 55, 2, 3, 67, 69], is an open issue for LDCS. Our goal 29 is to analyze the well-posedness of the one-step nonsmooth problem (OSNSP) obtained after an 30

implicit (backward) Euler discretization, and the convergence of the discrete-time solutions towards a

³² continuous-time limit which is a solution to the continuous-time system. The second step is useful

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- not only for proving the existence of solutions to the continuous-time system, but also, perhaps most
- ³⁴ importantly, to show that the proposed time-discretization makes sense for numerical simulations.
- ³⁵ Passivity properties (passive systems, positive real transfer functions, Lur'e equations) are pivotal
- throughout the article. In section 2, a class of passive LDCS (which extends descriptor variable systems $\frac{1}{2}$
- with a positive real transfer function) is introduced. A specific Weierstrass form associated with these systems, when the state-space realization is minimal, is used. The well-posedness of the OSNSP
- associated with the implicit Euler discretization, is analysed in section 3. Convergence of the discrete-
- time solutions is studied in sections 4 and 5. Examples that illustrate the theoretical developments, are
- given in section 6. Conclusions are drawn in section 7. The appendix is dedicated to recall various
- 42 mathematical tools.

Notation and definitions: for any vector $x \in \mathbb{R}^n$ and any matrix $M \in \mathbb{R}^{m \times n}$, ||x|| is the Euclidean 43 norm and ||M|| is the Frobenius norm, which are compatible norms [16, Proposition 9.3.5], *i.e.*, $||Mx|| \leq 1$ 44 ||M|| ||x||. Let $M \in \mathbb{R}^{n \times m}$, then Im(M) is its range, Ker(M) is its null space. We use $\langle x, y \rangle = x^{\top} y$, so 45 $\langle x, x \rangle = ||x||^2$. Positive definite matrix: M > 0 if $x^{\top}Mx > 0$ for all $x \neq 0$, positive semidefinite matrix: 46 $M \ge 0$ if $x^{\top}Mx \ge 0$ for all x (such M is not necessarily symmetric). The maximum singular value 47 is denoted as $\sigma_{\max}(M)$, and the minimum and maximum eigenvalues as $\lambda_{\min}(M)$ and $\lambda_{\max}(M)$. $M_{\bullet i}$ 48 denotes the *i*-th column of *M*. More mathematical results and definitions are given in Appendices A, B, 49 C, D, E and G. 50

⁵¹ 2 THE CLASS OF PASSIVE LDCS

Let us consider the following singular dynamical system, that we may name a linear descriptor complementarity system (LDCS):

(2.1)
$$\begin{cases} P\dot{x}(t) = Ax(t) + B\lambda(t) + E(t), \text{ a.e. } t \ge 0\\ 0 \le \lambda(t) \perp w(t) = Cx(t) + D\lambda(t) + F(t) \ge 0, \text{ for all } t \ge 0 \end{cases}$$

with $x(t) \in \mathbb{R}^n$, $\lambda(t) \in \mathbb{R}^m$, A, B, C, D constant matrices of appropriate dimensions, $P \in \mathbb{R}^{n \times n}$ has rank p < n. Dissipativity is a fundamental property in Systems and Control [25, 87, 88]. The material in this article relies on a result in [28, 32, 51] on a special Weierstrass form for passive descriptor variable systems [25, Section 3.1.7] which possess a *minimal state-space realization* [28, Theorem 3.1]. The passivity of the quintuple (P, A, B, C, D) means the passivity of the operator $\lambda \mapsto w$, see section D. Assuming that E(t) = 0 and F(t) = 0, passivity and minimality, this special form writes as (see [51] [32, Equation (16)] [28, Proposition A.3]):

(2.2)
$$\begin{cases} (a) \quad \dot{x}_1(t) = A_1 x_1(t) + B_1 \lambda(t) \\ (b) \quad \dot{x}_2(t) = z(t) + B_2 \lambda(t) \\ (c) \quad 0 = x_2(t) + B_3 \lambda(t) \\ (d) \quad 0 \le \lambda(t) \perp w(t) = C_1 x_1(t) + C_2 x_2(t) + C_3 z(t) + D\lambda(t) \ge 0, \end{cases}$$

where the quadruple (A_1, B_1, C_1, \hat{D}) is itself passive, $\hat{D} \stackrel{\Delta}{=} D - C_3 B_2 - C_2 B_3, x_1(t) \in \mathbb{R}^{n_1}, x_2(t) \in \mathbb{R}^{n_2}, z(t) \in \mathbb{R}^{n_2}, n = n_1 + 2n_2, \lambda(t) \in \mathbb{R}^m$, and there exists $X_3 = X_3^\top \ge 0$ such that $B_3^\top X_3 = -C_3$. The matrix X_3 is unique when the state-space realization is minimal [28, Theorem 5.4]. The LMI in (D.2) holds. Due to the passivity and the complementarity conditions which imply $\lambda(t)^\top w(t) = 0$ for all times, it follows that for all $t \ge 0$ one obtains $V(x_1(t), x_2(t)) \stackrel{\Delta}{=} \frac{1}{2}(x_1^\top X_1 x_1 + x_2^\top X_3 x_2)(t) \le V(x_1(0), x_2(0)), [32]$ and [28, Proof of Theorem 5.1], with $X_1 > 0$ if the system is strongly SPR and minimal [28, Theorem

- ⁶⁷ 5.4], see section D. In the latter case $x_1(\cdot)$ is bounded, A_1 is Hurwitz¹, $\hat{D} > 0$ and B_3 is full row rank
- 68 $(\Rightarrow n_2 \le m)$ [28, Proposition A.4].
- ⁶⁹ Remark 2.1. Notice that the differentiation index (shortly the index in the following) of the DAE part of
- ⁷⁰ the system (2.2) is equal to 2, since one needs to differentiate the third line twice (assuming $\lambda(t) = 0$)
- to recover the variable $\dot{z}(t)$. It is clear from (2.2) that the LDCS (2.1) may be seen as a DAE coupled
- ⁷² with nonsmooth constraints that take the form of complementarity conditions. It is noteworthy that
- ⁷³ singular zero order sweeping processes (named in ZOSwP [26]) have been studied [85, 6, 9], but they
- ⁷⁴ form another class of singular systems than (2.1) or (2.2).
- ⁷⁵ Remark 2.2. The system in (2.1) can be rewritten equivalently as:

(2.3)
$$P\dot{x}(t) - Ax(t) - E(t) \in B(N_{K^{\star}} + D)^{-1}(-Cx(t) - F(t))$$

with $K = K^* = \mathbb{R}^m_+$, which is, as pointed out in [26, section 3.16], a class of singular systems that has not yet received much attention. The same applies to (2.2) which is rewritten as:

(2.4)
$$\begin{pmatrix} I_{n_1} & 0 & 0 \\ 0 & I_{n_2} & 0 \\ 0 & 0 & 0 \end{pmatrix} \begin{pmatrix} \dot{x}_1(t) \\ \dot{x}_2(t) \\ \dot{z}(t) \end{pmatrix} - \begin{pmatrix} A_1 & 0 & 0 \\ 0 & 0 & I_{n_2} \\ 0 & I_{n_2} & 0 \end{pmatrix} \begin{pmatrix} x_1(t) \\ x_2(t) \\ z(t) \end{pmatrix} \in \begin{pmatrix} B_1 \\ B_2 \\ B_3 \end{pmatrix} (\mathcal{N}_{K^{\star}} + D)^{-1} (-w(t)).$$

⁷⁸ Exogenous terms can also be considered in (2.2) and (2.4), though the analyses in sections 3 and 4

⁷⁹ are made without such terms. Also, an alternative way of analysing (2.2) is presented in section 5,

⁸⁰ incorporating exogenous signals. An example is treated in section 6.1 where exogenous terms are

taken into account in the special form (2.2), see (6.3) and (6.4). In the case D = 0, both (2.4) and (2.3) can be rewritten, under some assumptions, as a singular first-order sweeping process (FOSwP [26]).

can be rewritten, under some assumptions, as a singular first-order sweeping process (FOSwP [26]).
 The optimal control of FOSwP as treated in [42, 41, 36], relies on the analysis of the time-discretization

and the convergence of solutions. Such an important issue could be tackled for singular FOSwP, where

⁸⁵ the results in this article may serve as a starting point.

3 THE EULER IMPLICIT SCHEME AND THE OSNSP

In the sequel of the article, we let $t \in [0, T]$, $0 < T < +\infty$, $h = \frac{T}{N}$, $N \in \mathbb{N} \setminus \{0\}$, $t_k = kh$, $k \in \{0, 1, ..., N\}$, and $f_k = f(t_k)$. The studied implicit Euler time-discretization of (2.2) is as follows:

(3.1)
$$\begin{cases} x_{1,k+1} = x_{1,k} + hA_1x_{1,k} + hB_1\lambda_{k+1} \\ x_{2,k+1} = x_{2,k} + hz_{k+1} + hB_2\lambda_{k+1} \\ 0 = x_{2,k+1} + B_3\lambda_{k+1} \\ 0 \le \lambda_{k+1} \perp w_{k+1} = C_1x_{1,k+1} + C_2x_{2,k+1} + C_3z_{k+1} + D\lambda_{k+1} \ge 0, \end{cases}$$

for all $k \in \{0, 1, ..., N-1\}$, where $A_1x_{1,k}$ could be replaced by $A_1x_{1,k+1}$ without modifying much the next developments. Nevertheless we choose to call the method in (3.1) an implicit method despite of the fact that the linear term is explicitly calculated, because our objective is to focus solely on the set-valued part of the system (*i.e.*, the complementarity conditions). The set of equations in (3.1) makes a mixed linear complementarity problem (MLCP). After few calculations from (3.1) the following OSNSP is obtained (where the variable hw_{k+1} is used instead of w_{k+1}) as an LCP:

$$0 \le \lambda_{k+1} \perp \bar{D}_h \lambda_{k+1} + h\beta_k \ge 0 \quad \Leftrightarrow \quad \bar{D}_h \lambda_{k+1} + h\beta_k \in -\mathcal{N}_{\mathbb{R}^m_+}(\lambda_{k+1})$$
$$\Leftrightarrow \quad \lambda_{k+1} \in (\bar{D}_h + \mathcal{N}_{\mathbb{R}^m_+})^{-1}(-h\beta_k),$$

(3.2)

¹That is, all its eigenvalues have a negative real part.

- with $\bar{D}_h \stackrel{\Delta}{=} h\hat{D} + h^2 C_1 B_1 + B_3^\top X_3 B_3$, $\beta_k \stackrel{\Delta}{=} C_1 (I_{n_1} + hA_1) x_{1,k} \frac{1}{h} C_3 x_{2,k}$. It is noteworthy that z_{k+1} is eliminated and does not appear in β_k . The inclusion (3.2) is a VI as (B.2) with $\varphi(\cdot) = \Psi_{\mathbb{R}^m_+}(\cdot)$. Obviously the well-posedness of the OSNSP depends on the properties of \bar{D}_h . From the passivity we know that $\hat{D} \ge 0$ and $B_3^\top X_3 B_3 \ge 0$. Using Theorem A.1 and Corollary A.2 it is possible to characterize the set of
- admissible hC_1B_1 that do not destroy the positive definiteness of \overline{D}_h in case $\hat{D} > 0$.
- Remark 3.1. We do not investigate here whether or not the discrete time system (3.1) is dissipative. Such a study is made in [55] for the case of LCS discretized with a θ -method.
- As alluded to in Remark 2.2, the cone \mathbb{R}^m_+ can be replaced in (2.1), (2.2) and (3.1) by any closed non empty convex cone $K \subset \mathbb{R}^m$ and its dual cone K^* . Then the complementarity conditions become
- empty convex cone $K \subset \mathbb{R}^m$ and its dual cone K^* . Then the complementarity control $K^* \ni \lambda(t) \perp w(t) = C_1 x_1(t) + C_2 x_2(t) + C_3 z(t) + D\lambda(t) \in K$, yielding the inclusion
- $104 \quad \mathbf{K} = \mathcal{I}_{\mathbf{K}}(t) + \mathcal{I}_{\mathbf{$
 - (3.3) $\bar{D}_h \lambda_{k+1} + h \beta_k \in -\mathcal{N}_{K^*}(\lambda_{k+1})$ $\Leftrightarrow \lambda_{k+1} \in (\bar{D}_h + \mathcal{N}_{K^*})^{-1}(-h \beta_k),$

which is a linear cone complementarity problem (LCCP), instead of (3.2). One can also consider timevarying terms in (2.2), in both the dynamics and the complementarity conditions, and adapt β_k accordingly. Thus in the next proposition a more general complementarity framework is considered where

the positive orthant is replaced by a closed convex cone.

- Proposition 3.2. Let h > 0 and $k \ge 0$ be given. Assume that $\overline{D}_h \ge 0$, not necessarily symmetric, and let
- 110 $K \subset \mathbb{R}^m$ be a non empty closed convex cone (where the complementarity conditions in (2.2) are extended 111 to $K^* \ni \lambda(t) \perp w(t) = C_1 x_1(t) + C_2 x_2(t) + C_3 z(t) + D\lambda(t) \in K$). Let $S = K^* \cap \text{Ker}(\bar{D}_h + \bar{D}_h^\top) \cap \{w \in \mathbb{R}^m \mid \bar{D}_h w \in K\}$. Then:
- 1. If $S = \{0\}$, for each $h\beta_k$ the LCCP in (3.3) has at least a solution.
- 114 2. If $S \neq \{0\}$, and there exists $\xi_0 \in \mathbb{R}^m$ such that $\langle h\beta_k \bar{D}_h^\top \xi_0, v \rangle + \Psi_{K^*}(v) > 0$ for all $v \neq 0, v \in S$, 115 then the LCCP in (3.3) has at least a solution.
- 3. If λ_{k+1}^1 and λ_{k+1}^2 are two solutions of the LCCP, then $\lambda_{k+1}^1 \lambda_{k+1}^2 \in \text{Ker}(\bar{D}_h + \bar{D}_h^\top)$.
- **Proof:** Let us consider Proposition B.2. Comparing (3.3) and (B.2) one has $\varphi(\cdot) = \Psi_{K^*}(\cdot)$, $\mathbf{M} = \overline{D}_h$,
- q = $h\beta_k$. Using Proposition B.1 c) d), it follows that $(dom(\varphi))_{\infty} = K^*$. The second set in S is obvious.
- The third set necessitates the calculation of $(\text{dom}(\varphi_{\infty}))^{\star}$, see (B.3). Again the result of the calculation
- follows from Proposition B.1 c) and d). Then the proof follows from Proposition B.2 a), b) and c).
- Once the existence of λ_{k+1} is established for any $k \ge 0$, then the existence of $x_{1,k+1}$, $x_{2,k+1}$ and z_{k+1} can be inferred from (3.1) for any h > 0.
- Remark 3.3. If $K = \mathbb{R}^m_+$ then other results like [43, Theorem 3.8.6] may be used to study the LCP (3.2), using for instance the copositivity of \overline{D}_h .
- Let us state a consequence of Proposition 3.2. The definition of a strongly SPR (SSPR) system is recalled in Appendix D.
- ¹²⁷ Corollary 3.4. Let the quintuple (P, A, B, C, D) be SSPR. Then, the LCCP in (3.3) has a unique solution for ¹²⁸ any $h\beta_k$, for h > 0 small enough.
- **Proof:** This follows from [28, Theorem 5.4] ($\Rightarrow \hat{D} > 0$), the fact that $B_3^{\top} X_3 B_3 \ge 0$, Corollary A.2

 $(\Rightarrow \bar{D}_h > 0 \text{ for } 0 < h < h^* \stackrel{\Delta}{=} \frac{1}{||C_1B_1||_2 \|(\frac{\dot{D}+\dot{D}^{\top}}{2})^{-1}\|_2} \Rightarrow S = \{0\}), \text{ and items 1 and 3 in Proposition 3.2. }$

It is noteworthy that the SSPRness is only sufficient to guarantee that $\hat{D} > 0$, which is itself a sufficient condition for the existence of solutions to the LCCP (3.3).

133 4 BOUNDEDNESS AND CONVERGENCE ANALYSIS

From now on, we assume that the complementarity conditions are given as in (2.2) and (3.1), *i.e.*, $K = K^* = \mathbb{R}^m_+$, see Remark 4.2. The next result characterizes least-norm solutions of the LCP.

Corollary 4.1. Assume that the conditions for existence of a solution in Proposition 3.2 are satisfied. Then the set of solutions to the LCP in (3.2) is nonempty convex polyhedral and contains a unique least-norm element $\lambda_{k+1}^{\min}(x_{1,k}, x_{2,k})$ such that $||\lambda_{k+1}^{\min}|| \leq ||\lambda_{k+1}||$ for any other solution λ_{k+1} . Moreover there exists a constant $\eta > 0$ depending only on $\overline{D}_h(=h\hat{D}+h^2C_1B_1+B_3^{T}X_3B_3)$ such that $||\lambda_{k+1}^{\min}|| \leq \eta h||\beta_k||$.

- Proof: Follows from [43, Theorem 3.1.7] and [57, Lemma 1].
- The constant η generally depends on h since it depends on \overline{D}_h , and we may denote it as $\eta(h)$ (as the
- examples in section 6 show, it may also be constant in some cases). This result certainly is important
- because it will subsequently allow us to get upper bounds on the iterates of the implicit scheme (3.1).
- Notice that if $\bar{D}^h = 0$ then $\beta_{k,i} > 0 \Rightarrow \lambda_{k+1,i} = 0$, and $\beta_{k,i} = 0 \Rightarrow \lambda_{k+1,i} \in \mathbb{R}_+$, hence $\lambda_{k+1}^{\min} = 0$ so that n = 0. If h = 0 then $\bar{D}^0 = P^\top X P$, $\beta_{k+1,i} = 0$, and $\eta_{k+1,i} = 0$, where $\lambda_{k+1,i} \in \mathbb{R}_+$, hence $\lambda_{k+1}^{\min} = 0$ so
- that $\eta = 0$. If h = 0 then $\overline{D}^0 = B_3^\top X_3 B_3$, $\beta_0 = -C_3 x_{2,k}$, and $\eta = \eta_0$ such that $||\lambda_{k+1}^{\min}|| \le \eta_0 ||C_3 x_{2,k}||$ is independent of h. If the LCP $(\overline{D}_h, h\beta_k)$ enjoys the stability of solutions property [43, Theorem 7.3.12],
- independent of *h*. If the LCP(D_h , $h\beta_k$) enjoys the stability of solutions property [43, Theorem 7.3.12], then a small perturbation of \bar{D}_h implies a small variation of the solution, hence of λ_{k+1}^{\min} , and consequently
- of the constant η . Assume that $C_3 = 0$ (this is Assumption 2 below). Then solutions of the LCP $(\bar{D}_h, h\beta_k)$
- are also solutions of LCP $(\hat{D} + hC_1B_1, C_1(I_{n_1} + hA_1)x_{1,k})$. If $\hat{D} + hC_1B_1$ is a P-matrix (which occurs if \hat{D} is
- a P-matrix and h > 0 is small enough, see Corollary A.2), then it follows from [3, Theorem B.3] [34]
- that the unique solution of LCP($\hat{D} + hC_1B_1, C_1(I_{n_1} + hA_1)x_{1,k}$) (which is also the least-norm solution)
- satisfies $||\lambda_{k+1}^{\min}|| \leq \eta ||C_1(I_{n_1} + hA_1)x_{1,k}||$ for some constant $\eta > 0$.
- Remark 4.2. As long as a minimal norm element exists and satisfies the property stated in the corollary,
- one may consider an LCCP as in (3.3) instead of an LCP (3.2).

155 4.1 CONVERGENCE ANALYSIS

From (3.2) and under the conditions of Corollary 4.1, the mapping $\beta_k \mapsto \lambda_{k+1}(-h\beta_k) = \lambda_{k+1}^{\min}$ with λ_{k+1}^{\min}

the least-norm element of the convex polyhedral set $(\bar{D}_h + \mathcal{N}_{\mathbb{R}^m_+})^{-1}(-h\beta_k)$, is a single-valued mapping.

The implicit scheme (3.1) may be modified as follows (which may be named the "minimal norm" implicit scheme):

$$(4.1) \qquad \begin{cases} x_{1,k+1} = (I_{n_1} + hA_1)x_{1,k} + hB_1\lambda_{k+1}^{\min}(x_{1,k}, x_{2,k}) \\ x_{2,k+1} = x_{2,k} + hz_{k+1} + hB_2\lambda_{k+1}^{\min}(x_{1,k}, x_{2,k}) \\ 0 = x_{2,k+1} + B_3\lambda_{k+1}^{\min}(x_{1,k}, x_{2,k}) \\ 0 \le \lambda_{k+1}^{\min} \perp w_{k+1} = C_1x_{1,k+1} + C_2x_{2,k+1} + C_3z_{k+1} + D\lambda_{k+1}^{\min} \ge 0 \\ ||\lambda_{k+1}^{\min}(x_{1,k}, x_{2,k})|| \le \eta ||h\beta_k||, \quad h\beta_k = hC_1(x_{1,k} + hA_1x_{1,k}) - C_3x_{2,k}, \quad \eta > 0 \end{cases}$$

for all $k \in \{0, ..., N-1\}$. The above results explore the well-posedness of the OSNSP, for fixed h > 0. Our goal is now to analyse the boundedness of the iterates on the integration interval [0, T], T > 0.

For all $k \ge 1$ we have $x_{2,k} = -B_3 \lambda_k^{\min}$, and $hz_{k+1} = -B_3 \lambda_{k+1}^{\min} - x_{2,k} - hB_2 \lambda_{k+1}^{\min}$, thus it is inferred that

 $x_{2,k+1} = -B_3(\lambda_{k+1}^{\min} + \lambda_k^{\min}) - x_{2,k}$. Consequently (4.1) is rewritten equivalently as:

$$(4.2) \qquad \begin{cases} x_{1,k+1} = (I_{n_1} + hA_1)x_{1,k} + hB_1\lambda_{k+1}^{\min}(x_{1,k}, x_{2,k}) \\ x_{2,k+1} = -x_{2,k} - B_3(\lambda_{k+1}^{\min}(x_{1,k}, x_{2,k}) + \lambda_k^{\min}(x_{1,k-1}, x_{2,k-1})) \\ 0 \le \lambda_{k+1}^{\min} \perp w_{k+1} = C_1(I_{n_1} + hA_1)x_{1,k} - (C_2 + \frac{1}{h}C_3)x_{2,k} + \frac{1}{h}\bar{D}_h\lambda_{k+1}^{\min} - C_2B_3\lambda_k^{\min} \ge 0 \\ 0 = x_{2,k+1} + B_3\lambda_{k+1}^{\min}(x_{1,k}, x_{2,k}), \\ hz_{k+1} = -B_3\lambda_{k+1}^{\min}(x_{1,k}, x_{2,k}) - x_{2,k} - hB_2\lambda_{k+1}^{\min}(x_{1,k}, x_{2,k}) \\ ||\lambda_{k+1}^{\min}(x_{1,k}, x_{2,k})|| \le \eta ||h\beta_k||, \quad h\beta_k = hC_1(I_{n_1} + hA_1)x_{1,k} - C_3x_{2,k}, \quad \eta > 0, \end{cases}$$

for all $k \in \{1, N - 1\}$, given initial data $x_{1,0}$ and $x_{2,0}$. Starting from (4.2) the next result is obtained:

Lemma 4.3. Let $\alpha_{1,k} \stackrel{\Delta}{=} ||x_{1,k}||, \alpha_{2,k} \stackrel{\Delta}{=} ||x_{2,k}||, \alpha_k = (\alpha_{1,k} \ \alpha_{2,k})^\top$, then:

(4.3)
$$\begin{cases} \alpha_{k+1} \le (M_i M + M_{i-1} U) \alpha_{k-i} + M_i U \alpha_{k-i-1} \\ M_{i+1} = M_i M + M_{i-1} U, \ M_0 = I_2, M_{-1} = 0, \ 0 \le i \le k-1, \ k \ge 1, \end{cases}$$

where $M = \begin{pmatrix} a & b \\ c & d \end{pmatrix}$, $U = \begin{pmatrix} 0 & 0 \\ e & f \end{pmatrix}$, $a = ||I_{n_1} + hA_1|| (1 + h^2\eta ||B_1|| ||C_1||)$, $b = h\eta ||B_1|| ||C_3||$, $c = h\eta ||B_3|| ||C_1|| ||I_{n_1} + hA_1||$, $d = (1 + \eta ||B_3|| ||C_3||)$, e = c, $f = \eta ||B_3|| ||C_3||$. In particular

(4.4)
$$\alpha_{k+1} \le (M_{k-1}M + M_{k-2}U)\alpha_1 + M_{k-1}U\alpha_0$$

168 for all $k \ge 1$.

(4.5)

¹⁶⁹ **Proof:** The proof starts from (4.2). The following upper bounds can be derived:

 $\begin{aligned} (a) \ ||\lambda_{k+1}^{\min}|| &\leq \eta h ||C_1|| \ ||I_{n_1} + hA_1|| \ ||x_{1,k}|| + \eta ||C_3|| \ ||x_{2,k}|| \\ (b) \ ||x_{1,k+1}|| &\leq ||I_{n_1} + hA_1|| \ (1 + h^2\eta \ ||B_1|| \ ||C_1||) \ ||x_{1,k}|| + h\eta ||B_1|| \ ||C_3|| \ ||x_{2,k}|| \\ &= a ||x_{1,k}|| + b ||x_{2,k}|| \\ \end{aligned}$ $(c) \ ||x_{2,k+1}|| &\leq (1 + \eta \ ||B_3|| \ ||C_3||) \ ||x_{2,k}|| + \eta \ ||B_3|| \ ||C_3|| \ ||x_{2,k-1}|| \\ &+ h\eta ||B_3|| \ ||C_1|| \ ||I_{n_1} + hA_1|| \ (||x_{1,k}|| + ||x_{1,k-1}||) \\ &= d ||x_{2,k}|| + f ||x_{2,k-1}|| + c (||x_{1,k}|| + ||x_{1,k-1}||). \end{aligned}$

¹⁷⁰ One infers that:

$$(4.6) \qquad \qquad \alpha_{k+1} \le M\alpha_k + U\alpha_{k-1}$$

The upper bound in (4.3) therefore holds for i = 0. Let us prove by induction that it holds for all $0 \le i \le k - 1$. Let it hold for i > 0, we have $\alpha_{k-i} \le M\alpha_{k-i-1} + U\alpha_{k-i-2}$. Thus $\alpha_{k+1} \le (M_iM + M_{i-1}U)M\alpha_{k-i-1} + (M_iM + M_{i-1}U)U\alpha_{k-i-2} + M_iU\alpha_{k-i-1} = (M_{i+1}M + M_iU)\alpha_{k-i-1} + M_{i+1}U\alpha_{k-i-2}$. The result is proved. \boxtimes

¹⁷⁵ Our goal is now to examine (4.4). Let us come back to the matrices M and U in Lemma 4.3. It is clear

from (4.4) that the boundedness of the sequence $\{\alpha_k\}_{k\geq 2}$ depends on the boundedness of M_k , $k \geq 1$.

Apart from *d* and *f*, all the other terms in *M* and *U* are proportional to *h*, ηh or to $h_1 \stackrel{\Delta}{=} ||I_{n_1} + hA_1||$.

Moreover d > 1 whenever $B_3 \neq 0$ and $C_3 \neq 0$. The term in position (2, 2) of M^n is equal to $d^n + O(h_1h^2)$.

¹⁷⁹ Before stating the next result, let us introduce the following assumptions:

Assumption 1. The matrix A_1 is such that $h_1 \le 1$ for h > 0 small enough.

- Let us notice that the induced norm $|| \cdot ||_{2,2} = \sigma_{\max}(\cdot)$ [16, Proposition 9.4.9] is compatible with the
- vector norm $|| \cdot ||_2 = || \cdot ||$ [16, pp.607-609], and can be used instead of the Frobenius norm to compute h_1 . The next assumption is crucial in the body of this section.
- Assumption 2. One has $B_{3,\bullet i} \in \text{Ker}(X_3)$ for all $1 \le i \le m \implies C_3 = 0$.
- Let us remind that X_3 is uniquely defined under minimality and positive-realness [28, Theorem 5.4]. If
- Assumption 2 holds then d = 1 and b = f = 0. This assumption means that the variable z does not
- ¹⁸⁷ enter the complementarity conditions, see (2.2). Hence in a sense the complementarity part of the
- LDCS and the algebraic part are decoupled, and β_k no longer depends on $\frac{1}{h}C_3x_{2,k}$. The terms in (4.5)
- involve products with $\eta(h)$, and it is important to guarantee that they are bounded, when Assumption 2 holds.
- Assumption 3. The constant $\eta = \eta(h)$ is such that the products $\eta(h)h$ can be made arbitrarily small for small enough $h \ge 0$.
- Notice that this assumption is verified in the examples of section 6, where η is a constant (but these
- examples do not verify Assumption 2, however). The matrices M_k then involve only products of powers
- of the terms *a*, *b*, *c*, *d* = 1, *e*. These products are bounded for bounded *k*, and provided that both *h* and *h*₁ are small enough so that $a \le 1$, $b \le 1$, $c \le 1$, d = 1, $e \le 1$, they remain bounded as $k \to +\infty$. If the
- h₁ are small enough so that $a \le 1$, $b \le 1$, $c \le 1$, d = 1, $e \le 1$, they remain bounded as $k \to +\infty$. If the inequalities are strict (apart form d = 1), the first rows in both $M_{k-1}M + M_{k-2}U$ and $M_{k-1}U$ converge
- to zero, while the second rows converge to (0,1) (the 1 being the consequence of a term d^n in matrices
- M^n). The next proposition is therefore proved:
- Proposition 4.4. Assume that: a) the conditions of existence of solutions in Proposition 3.2 are verified, b)
- Assumptions 1, 2 and 3 hold true, and c) both h and h_1 are small enough, so that $a \le 1, b \le 1, c \le 1, d = 1$,
- $_{202}$ $e \leq 1$. Then the sequences $\{\alpha_{1,k}\}_k$ and $\{\alpha_{2,k}\}_k$ are uniformly bounded in k. If $A_1 < 0$, and h > 0 is small
- enough such that a < 1, c < 1, e < 1, then $\alpha_{1,k} \rightarrow 0$ as $k \rightarrow +\infty$.
- Recall that $\alpha_{1,k} = ||x_{1,k}||$ and $\alpha_{2,k} = ||x_{2,k}||$. The last assertion of the proposition follows from the material in Appendix A. Then we have the next corollary.

Corollary 4.5. Let the conditions a), b), c) of Proposition 4.4 hold. Then $||\lambda_{k+1}^{\min}||$ is uniformly bounded in k and h > 0, and it converges to zero as $k \to +\infty$ if $\alpha_{1,k} \to 0$ as $k \to +\infty$. Therefore the sequence $\left\{\frac{x_{1,k+1}-x_{1,k}}{h}\right\}_{k}$ is uniformly bounded also in k and h > 0.

Proof: The boundedness of $||\lambda_{k+1}^{\min}||$ and its convergence follow from (4.5) (a), the boundedness of $||x_{1,k}||$ and Assumption 2. From the first line in (4.2), one has $\frac{x_{1,k+1}-x_{1,k}}{h} = A_1x_{1,k} + B_1\lambda_{k+1}^{\min}$, hence the result is proved. \boxtimes

Let us now study the convergence of the piecewise-linear and step-function approximations defined as:

(4.7)
$$\begin{cases} x_1^h(t) = x_{1,k+1} + \frac{t_{k+1}-t}{h}(x_{1,k} - x_{1,k+1}) \\ \dot{x}_1^h(t) = \frac{x_{1,k+1}-x_{1,k}}{h} & \text{for all } t \in [t_k, t_{k+1}) \\ x_{1,h}^{\star}(t) = x_{1,k+1} & \end{cases}$$

and similarly for the other variables.

Corollary 4.6. There exists subsequences denoted as $\{x_1^h(\cdot)\}_{h\geq 0}$ which converge uniformly to a limit $x_1(\cdot)$ in $C^0([0,T]; \mathbb{R}^{n_1})$, such that almost everywhere one has $\dot{x}_1(t) = A_1x_1(t) + \lambda_1(t)$, for some function $\lambda_1(\cdot)$. The sequence $\{x_1^{h\star}(\cdot)\}_{h\geq 0}$ converges to $x_1(\cdot)$ strongly in $\mathcal{L}_2([0,T]; \mathbb{R}^{n_1})$. Moreover the sequences $\{x_{2,h}^{\star}(\cdot)\}_{h\geq 0}$ and $\{\lambda_h^{\min,\star}(\cdot)\}_{h\geq 0}$ converge weakly \star in $\mathcal{L}_{\infty}([0,T]; \mathbb{R}^{l})$, $l = n_2$ and l = m, respectively, to limits $x_2(\cdot)$ and $\lambda(\cdot)$ which satisfy the continuous-time constraints.

Proof: Both $\{x_1^h(\cdot)\}_h$ and $\{\dot{x}_1^h(\cdot)\}_h$ are uniformly bounded on [0, T]. The sequence of $C^0([0, T]; \mathbb{R}^{n_1})$ functions $\{x_1^h(\cdot)\}_h$ is thus equicontinuous. By the Arzela-Ascoli Theorem this sequence stays in a compact subset of $C^0([0, T]; \mathbb{R}^{n_1})$ and converges uniformly towards a continuous limit $x_1(\cdot) : [0, T] \rightarrow \mathbb{R}^{n_1}$ as $h \to 0, h > 0$. Let us notice that $x_1(t) \to 0$ as $t \to +\infty$. Indeed from the above results which

are stated for $k \in \{0, ..., N\}$ and $h > 0, k \to +\infty$ only if $T = Nh \to +\infty$. Let us now define the step functions $x_1^{h\star}(t) = x_{1,k+1}$ if $t \in [t_k, t_{k+1}), \lambda_{h\star}^{\min}(t) = \lambda_{k+1}^{\min}$ if $t \in [t_k, t_{k+1})$, and

(4.8) $\dot{x}_1^h(t) = A_1 x_1^{h\star}(t) + B_1 \lambda_{h\star}^{\min}(t), \text{ almost everywhere.}$

²²⁶ In addition one has:

$$(4.9) = \sum_{k=0}^{N-1} \int_{t_k}^{t_{k+1}} (t_{k+1} - t)^2 ||\dot{x}_{1,k}||^2 dt \le C^2 \sum_{k=0}^{N-1} \int_{t_k}^{t_{k+1}} (t_{k+1} - t)^2 dt \le C^2 \sum_{k=0}^{N-1} \int_{t_k}^{t_{k+1}} (t_{k+1} - t)^2 dt \le C^2 \sum_{k=0}^{N-1} \int_{t_k}^{t_{k+1}} (t_{k+1} - t)^2 dt$$

where $\dot{x}_{1,k} = \frac{x_{1,k+1} - x_{1,k}}{h}$, and $||\dot{x}_{1,k}|| \le C$ for all $k \ge 0$ for some C > 0 from Corollary 4.5. It follows that 227 $x_1^{h\star}(\cdot) \to x_1(\cdot)$ strongly in $\mathcal{L}_2([0,T];\mathbb{R}^{n_1})$ as $h \to 0$. On the other hand since $\dot{x}_1^h(\cdot) \in \mathcal{L}_{\infty}([0,T];\mathbb{R}^{n_1})$ 228 it follows from the Banach-Alaoglu-Bourbaki Theorem [18, Theorem III.15] that $\dot{x}_1^h(\cdot) \rightarrow \dot{x}_{1,\text{lim}}(\cdot)$ in the 229 weak \star topology as $h \to 0$. Since the considered domain [0, T] is bounded, one has $\mathcal{L}_1([0, T]; \mathbb{R}^{n_1}) \subseteq$ 230 $\mathcal{L}_2([0,T];\mathbb{R}^{n_1})$, thus $\dot{x}_1^h(\cdot)$ converges weakly in $\mathcal{L}_2([0,T];\mathbb{R}^{n_1})$. Moreover since $x_1^h(t) = x_1^h(0) + x_1^h(t) = x_1^h(t) = x_1^h(t) + x_1^h(t) + x_1^h(t) = x_1^h(t) + x_1^h(t) + x_1^h(t) + x_1^h(t) = x_1^h(t) + x_1^h(t)$ 231 $\int_{0}^{t} \dot{x}_{1}^{h}(t) dt$, one has $\dot{x}_{1,\text{lim}}(\cdot) = \dot{x}_{1}(\cdot)$ almost everywhere. Therefore we deduce using (4.8) that $(B_{1}\lambda_{h\star}^{\min})(\cdot)$ 232 converges weakly in $\mathcal{L}_2([0,T];\mathbb{R}^{n_1})$ as $h \to 0$ towards some limit $\lambda_1(\cdot)$, and we have $\dot{x}_1(t) = A_1 x_1(t) + A_1 x_1(t$ 233 $\lambda_1(t)$ almost everywhere on [0, T]234

Let us consider now the step-functions $x_{2,h}^{\star}(\cdot)$ and $\lambda_h^{\min,\star}(\cdot)$ defined as $x_{2,h}^{\star}(t) = x_{2,k+1}$ for all $t \in$ 235 $[t_k, t_{k+1})$, and similarly for the multiplier. It follows from Proposition 4.4 and from Corollary 4.5 236 that both sequences $\{x_{2,h}^{\star}(\cdot)\}_{h\geq 0}$ and $\{\lambda_h^{\min,\star}(\cdot)\}_{h\geq 0}$ are uniformly bounded. Hence from the Banach-Alaoglu-Bourbaki Theorem, there exist subsequences of both which converge weakly \star to limits 237 238 $x_2 \in \mathcal{L}_{\infty}([0,T];\mathbb{R}^{n_2})$ and $\lambda \in \mathcal{L}_{\infty}([0,T];\mathbb{R}^m)$, respectively. We infer that $0 = x_{2h}^{\star}(t) + B_3 \lambda_h^{\star}(t) \rightarrow 0$ 239 $x_2(t) + B_3\lambda(t)$ for almost all $t \ge 0$ in the weak \star sense. Next, using the second line in (4.1) it follows 240 that the sequence $\{hz_k\}_{h\geq 0}$ is also uniformly bounded. Denoting $\sigma_k = hz_k$, it follows that the sequence 241 $\{\sigma_h^{\star}\}_{h\geq 0}$ converges weakly \star to a limit denoted as $\sigma(t) \in \mathcal{L}_{\infty}([0,T]; \mathbb{R}^{n_2})$. Moreover one has $\sigma(t) = 0$ 242 almost everywhere since $x_{2,h}^{\star}(t) - x_{2,h}^{\star}(t-t_k) - hB_2\lambda_h^{\min,\star}(t) = \sigma_h^{\star}(t)$ almost everywhere in [0,T]. 243 Finally, the complementarity conditions in (4.1) are equivalently rewritten as $0 \le \lambda_h^{\star}(t) \perp C_1 x_{1h}^{\star}(t) +$ 244 $C_2 x_{2,h}^{\star}(t) + D\lambda_h^{\star}(t) \ge 0$ for all $t \in [t_k, t_{k+1})$, since $C_3 = 0$ by Assumption 2. It is deduced that in the 245 limit as $h \to 0$ one obtains $0 \le \lambda(t) \perp C_1 x_1(t) + C_2 x_2(t) + D\lambda(t) \ge 0$. 246

247 4.2 STABILITY ANALYSIS VIA A LYAPUNOV FUNCTION

To complete the foregoing section, it is natural to analyze how the Lyapunov function $V(x_1, x_2)$ for the continuous-time system in (2.2) (see also appendix D), could be used in the discrete-time context for (4.2). Indeed implicit Euler methods are known to be able to mimic the continuous-time passivity and stability properties, in some cases [55]. It is noteworthy that Assumption 2 implies that $X_3x_{2,k} = -X_3B_3\lambda_k = 0$ and $X_3x_2(t) = -X_3B_3\lambda(t) = 0$ for all k and t. Therefore $V(x_1, x_2) = \frac{1}{2}x_1^{\mathsf{T}}X_1x_1$, so that this Lyapunov function is not a positive definite function of the whole state. For the sake of completeness let us nevertheless make the calculations in the general setting. We set $V_k \stackrel{\Delta}{=} \frac{1}{2}x_{1,k}^{\mathsf{T}}X_1x_{1,k} + x_{2,k}^{\mathsf{T}}X_3x_{2,k}$. It follows 255 that:

$$V_{k+1} - V_{k} = \frac{1}{2} (x_{1,k+1} - x_{1,k})^{T} X_{1} (x_{1,k+1} + x_{1,k}) + (x_{2,k+1} - x_{2,k})^{T} X_{3} (x_{2,k+1} + x_{2,k}) \\ = \frac{1}{2} (hA_{1}x_{1,k} + hB_{1}\lambda_{k+1})^{T} X_{1} (2x_{1,k} + hA_{1}x_{1,k} + hB_{1}\lambda_{k+1}) \\ + (hz_{k+1} + hB_{2}\lambda_{k+1})^{T} X_{3} (2x_{2,k} + hz_{k+1} + hB_{2}\lambda_{k+1}) \\ = h(x_{1,k}^{T} A_{1}^{T} X_{1,k,k} + \lambda_{k+1}^{T} B_{1}^{T} X_{1,k,l}) + h^{2} (\frac{1}{2} x_{1,k}^{T} A_{1}^{T} X_{1,k,k} + \frac{1}{2} \lambda_{k+1}^{T} B_{1}^{T} X_{1}B_{1}\lambda_{k+1}) \\ + h^{2} x_{1,k}^{T} A_{1}^{T} X_{1}B_{1}\lambda_{k+1} + hz_{k+1}^{T} C_{3}\lambda_{k+1} + h\lambda_{k+1}^{T} C_{3}B_{2}\lambda_{k+1} + hz_{k+1}^{T} X_{3}x_{2,k} \\ + h\lambda_{k+1}^{T} B_{2}^{T} X_{3}x_{2,k} + h\lambda_{k+1}^{T} w_{k+1} - h\lambda_{k+1}^{T} (\Delta B_{1} - B_{2}^{T} X_{3}B_{2} - \frac{1}{2} B_{1}^{T} X_{1}B_{1}))\lambda_{k+1} \\ + hx_{1,k}^{T} (B_{1}^{T} X_{1} - C_{1} - hA_{1}^{T} C_{1}^{T} + hA_{1}^{T} X_{1}B_{1})\lambda_{k+1} + \frac{h}{2} x_{2,k}^{T} X_{3}B_{2}\lambda_{k+1} \\ + \frac{h}{2} z_{k+1}^{T} X_{3}x_{2,k} + h\lambda_{k+1}^{T} w_{k+1} \\ + \frac{h}{2} z_{k+1}^{T} X_{3}x_{2,k} + h\lambda_{k+1}^{T} w_{k+1} \\ + \frac{h}{2} z_{k+1}^{T} X_{1}x_{1,k} - h^{2} \lambda_{k+1}^{T} (CB_{1} - B_{2}^{T} X_{3}B_{2} - \frac{1}{2} B_{1}^{T} X_{1}B_{1})\lambda_{k+1} \\ + \frac{h}{2} x_{1,k}^{T} (A_{1}^{T} X_{1} A_{1}x_{1,k} - h^{2} \lambda_{k+1}^{T} (CB_{1} - B_{2}^{T} X_{3}B_{2} - \frac{1}{2} B_{1}^{T} X_{1}B_{1})\lambda_{k+1} \\ + h^{2} x_{1,k}^{T} (-A_{1}^{T} C_{1}^{T} + A_{1}^{T} X_{1}B_{1})\lambda_{k+1} + \frac{h}{2} x_{2,k}^{T} X_{3}B_{2}\lambda_{k+1} + \frac{h}{2} x_{k+1}^{T} X_{3}x_{2,k} + h\lambda_{k+1}^{T} w_{k+1} \\ = \frac{h}{2} (x_{1,k}^{T} \lambda_{k+1}^{T}) Q_{1} \begin{pmatrix} x_{1,k} \\ \lambda_{k+1} \end{pmatrix} + h\lambda_{k+1}^{T} w_{k+1} \\ + h^{2} \left\{ \frac{1}{2} x_{1,k}^{T} A_{1}^{T} X_{1}A_{1}x_{1,k} - h^{2} \lambda_{k+1}^{T} (C_{1}B_{1} - B_{2}^{T} X_{3}B_{2} - \frac{1}{2} B_{1}^{T} X_{1}B_{1})\lambda_{k+1} \\ + h^{2} \left\{ \frac{1}{2} x_{1,k}^{T} A_{1}^{T} X_{1}A_{1}x_{1,k} - \lambda_{k+1}^{T} (C_{1}B_{1} - B_{2}^{T} X_{3}B_{2} - \frac{1}{2} B_{1}^{T} X_{1}B_{1})\lambda_{k+1} \\ + h^{2} \left\{ \frac{1}{2} x_{1,k}^{T} A_{1}^{T} X_{1}A_{1}x_{1,k} - \lambda_{k+1}^{T} (C_{1}B_{1} - B_{2}^{T} X_{3}B_{2} \lambda_{k+1} + \frac{h}{2} z_{k+1}^{T} X_{3}x_{2,k} \\ \end{pmatrix} \right\}$$

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(4.10)
$$= \frac{h}{2} (x_{1,k}^{\top}, \lambda_{k+1}^{\top}) Q_1 \begin{pmatrix} x_{1,k} \\ \lambda_{k+1} \end{pmatrix} + h \lambda_{k+1}^{\top} w_{k+1} \\ + h^2 \left\{ \frac{1}{2} x_{1,k}^{\top} A_1^{\top} X_1 A_1 x_{1,k} - \lambda_{k+1}^{\top} (C_1 B_1 - B_2^{\top} X_3 B_2 - \frac{1}{2} B_1^{\top} X_1 B_1) \lambda_{k+1} \\ + x_{1,k}^{\top} (-A_1^{\top} C_1^{\top} + A_1^{\top} X_1 B_1) \lambda_{k+1} \right\} - \underbrace{\frac{1}{2} x_{2,k}^{\top} X_3 x_{2,k} + \lambda_{k+1}^{\top} B_3^{\top} X_3 B_3 \lambda_k}_{=0 \text{ if Assumption 2 holds.}}$$

where the last two terms are obtained using the second equality in (3.1). Remind that $\lambda_{k+1}^{\top} w_{k+1} = 0$. Assume that $Q_1 < 0$, then let Assumption 2 hold. Let us rewrite (4.10) as $V_{k+1} - V_k = \frac{h}{2} (x_{1,k}^{\top}, \lambda_{k+1}^{\top}) [Q_1 + hQ_2] \begin{pmatrix} x_{1,k} \\ \lambda_{k+1} \end{pmatrix}$, where the indefinite matrix

(4.11)
$$Q_2 = \begin{pmatrix} A_1^{\mathsf{T}} X_1 A_1 & -\frac{1}{2} A_1^{\mathsf{T}} C_1^{\mathsf{T}} + \frac{1}{2} A_1^{\mathsf{T}} X_1 B_1 \\ (-\frac{1}{2} A_1^{\mathsf{T}} C_1^{\mathsf{T}} + \frac{1}{2} A_1^{\mathsf{T}} X_1 B_1)^{\mathsf{T}} & -C_1 B_1 + B_2^{\mathsf{T}} X_3 B_2 - \frac{1}{2} B_1^{\mathsf{T}} X_1 B_1 \end{pmatrix},$$

and we remind that Assumption 2 implies that $V_k = \frac{1}{2} x_{1,k}^\top X_1 x_{1,k}$. The matrix Q_2 is treated as a perturbation to Q_1 . Using Corollary A.2, there exists $h_{\max} > 0$ such that for all $h \in [0, h_{\max})$ one has $Q_3 \stackrel{\Delta}{=} Q_1 + hQ_2 < 0$. Using classical bounds on quadratic forms and vector/matrices products $(x^\top y \le ||x|| ||y|| \le \frac{1}{2} ||x||^2 + \frac{1}{2} ||y||^2$ [16, Corollary 9.1.7], if $Q = Q^\top$ then $\lambda_{\min}(Q) ||x||^2 \le x^\top Qx \le$ $\lambda_{\max}(Q) ||x||^2$ [16, Corollary 8.4.2], $||Qx|| \le ||Q|| ||x||$ for compatible norms [16, section 9.3]), and noting that we can replace Q_3 by its negative definite symmetric part in the quadratic form, it follows that:

(4.12)
$$V_{k+1} \leq V_k + \frac{h}{2}\lambda_{\max}(Q_3) ||(x_{1,k},\lambda_{k+1})||^2 \leq V_k + \frac{h}{2}\lambda_{\max}(Q_3) ||x_{1,k}||^2 \leq V_k + \frac{h}{2}\frac{\lambda_{\max}(Q_3)}{\lambda_{\max}(X_1)}V_k = \left(1 + \frac{h}{2}\frac{\lambda_{\max}(Q_3)}{\lambda_{\max}(X_1)}\right)V_k, \text{ for all } k \geq 0.$$

Since $\lambda_{\max}(Q_3) < 0$ (which is the reason why the second inequality in (4.12) holds), one has $1 + \frac{h}{2} \frac{\lambda_{\max}(Q_3)}{\lambda_{\max}(X_1)} < 1$ for all h > 0. Therefore all V_k are uniformly bounded for any bounded V_0 , and $V_k \to 0$ as $k \to +\infty$ for any $h \in (0, h_{\max})$. Let us set $\lambda_{k+1} = \lambda_{k+1}^{\min}$. Since the dependence of λ_{k+1}^{\min} in $x_{2,k}$ through β_k is $C_3 x_{2,k}$ (see (4.2)), it follows from Assumptions 2 and 3 that $\lambda_{k+1}^{\min} = \lambda_{k+1}^{\min}(x_{1,k})$, and $\lambda_{k+1}^{\min}(x_{1,k}) \to 0$ as $k \to +\infty$. Using the third line in (4.2), we infer that the iterates $x_{2,k}$ (which are are uniformly bounded for any bounded for any bounded $\lambda_{k+1}^{\min}(x_{1,k}) \to 0$ as a bove Lyapunov analysis, we have therefore proved the following:

Lemma 4.7. Let Assumptions 2 and 3 hold, and the quintuple (P, A, B, C, D) be strongly SPR and minimal.

Then there exists $h_{\text{max}} > 0$ depending on A_1 , C_1 , X_1 , X_3 , B_1 , B_2 , such that for all $h \in (0, h_{\text{max}})$, and for

all bounded initial condition $x_{1,0}$, the sequences $\{x_{1,k}\}_k$, $\{x_{2,k}\}_k$, $\{\lambda_k^{\min}\}_k$ are uniformly bounded and $x_{1,k}, x_{2,k}, \lambda_{k+1}^{\min} \to 0$ as $k \to +\infty$.

- Remark 4.8. Concerning Assumption 1, let us examine the constants *a*, *b*, *c*, *d*, *f* in Lemma 4.3, and the inequalities in (4.5), as well as the above Lyapunov analysis. The term h_1 is present in both *a* and *c*. The inequality $a \le 1$ implies $h_1 < 1$, but $c \le 1$ does not imply that $h_1 < 1$. The constant *a* is present in (4.5) (b), however the Lyapunov analysis guarantees the boundedness of $x_{1,k}$ without this assumption (the analysis relies on the passivity of (A_1, B_1, C_1, \hat{D}) and $h \ge 0$ small enough). But, Lemma 4.7 relies on $Q_1 < 0$, which is a strong passivity condition and may be restrictive in applications. There exists some
- 283 cases where SSPRness can be relaxed, as shown next.
- Let (A_1, B_1, C_1, \hat{D}) be strictly state passive $(\Rightarrow X_1A_1 + A_1^\top X_1 = -\mu X_1, \mu > 0, X_1 = X_1^\top > 0, \mu$ is called the dissipation constant) and $\hat{D} + \hat{D}^\top = 0 \Rightarrow X_1B_1 = C_1^\top$, then Q_2 in (4.11) and Q_1 can be rewritten as:

(4.13)
$$Q_2 = \begin{pmatrix} A_1^{\top} X_1 A_1 & 0 \\ 0 & B_2^{\top} X_3 B_2 - \frac{3}{2} B_1^{\top} X_1 B_1 \end{pmatrix}, \quad Q_1 = \begin{pmatrix} -\mu X_1 & 0 \\ 0 & 0 \end{pmatrix}$$

Thus $Q_1 + hQ_2 < 0$ if and only if $\mu X_1 - hA_1^{\top}X_1A_1 > 0$ and $B_2^{\top}X_3B_2 - \frac{3}{2}B_1^{\top}X_1B_1 < 0$. Let us recall that if (A_1, B_1, C_1, \hat{D}) is minimal, then strict state passivity and SPRness are equivalent, see Appendix **D**.

Lemma 4.9. Assume that: 1) the quadruple (A_1, B_1, C_1, \hat{D}) is strictly state passive with $\hat{D} + \hat{D}^{\top} = 0$ and dissipation constant $\mu > 0$, 2) $B_2^{\top} X_3 B_2 - \frac{3}{2} B_1^{\top} X_1 B_1 < 0$. Let Assumptions 2 and 3 hold. Then there exists $h_{\text{max}} > 0$ such that for all $h \in [0, h_{\text{max}})$, the sequences $\{x_{1,k}\}_k, \{x_{2,k}\}_k$ and $\{\lambda_k^{\min}\}_k$ are uniformly bounded.

Proof: Denote $\tilde{X}_1 = \mu X_1 - h A_1^{\top} X_1 A_1$, then $\tilde{X}_1 > 0$ for $h \ge 0$ small enough, using Corollary A.2. Using 2) and Assumption 2, and after few calculations, there exists $h_{\text{max}} > 0$ such that:

(4.14)
$$V_{k+1} \le \left(1 - \frac{h}{2} \frac{\lambda_{\min}(\tilde{X}_1)}{\lambda_{\max}(X_1)}\right) V_k$$

for all $h \in [0, h_{\max})$. Thus $V_k \to 0$ as $k \to +\infty$, and $x_{1,k} \to 0$ as $k \to +\infty$. It can be inferred that $||x_{1,k}|| \le \sqrt{\frac{\lambda_{\max}(X_1)}{\lambda_{\min}(X_1)}} ||x_{1,0}|| \alpha^{\frac{k}{2}}$ with $\alpha < 1$. Fom (4.5) (a) and Assumption 3 it follows that $||\lambda_{k+1}^{\min}|| \to 0$ as $k \to +\infty$, and $\{\lambda_k^{\min}\}_k$ is uniformly bounded. From (4.1) it is inferred that $x_{2,k} \to 0$ as $k \to +\infty$ and $\{x_{2,k}\}_k$ is uniformly bounded. \boxtimes

Similar convergence results as in Corollary 4.6 can then be deduced, and are not developed here for the sake of briefness. One sees from the above and Corollary A.2 that the degree of strict SPRness

(represented by the constant μ) is related to the value of h_{max} : the greater μ , the greater h_{max} . Finally 300 it is noteworthy that the strict state passivity in item 1) of lemma 4.9 can again be relaxed, provided 301 that the matrix $A_1^{\top}X_1 + X_1A_1 + hA_1^{\top}X_1A_1 \prec 0$. Consider as an example $A_1 = \begin{pmatrix} -1 & 0 \\ 0 & 0 \end{pmatrix}$, $X_1 = I_2$, then 302 $A_1^{\top}X_1 + X_1A_1 + hA_1^{\top}X_1A_1 = \begin{pmatrix} -2+h & 0\\ 0 & 0 \end{pmatrix} < 0 \text{ for all } 0 \le h < 2.$ 303

5 ANOTHER APPROACH FOR PASSIVE SYSTEMS 304

The LDCS in (2.2) are suitably transformed in [32] to enable their analysis. The transformation may be 305 seen as a kind of relative degree reduction, involving the multiplier derivative. As will become clear 306 later, this new dynamics allows us to somewhat relax the stringent Assumption 2 (see Assumption 4 c) 307 below). More accurately, a close relationship is shown in [32, Lemma 3, Theorem 3] between (2.2) and 308

the following system: 309

(5.1)
$$\begin{cases} (a) \begin{pmatrix} I_{n_1} & 0\\ 0 & B_3^{\mathsf{T}} X_3 B_3 \end{pmatrix} \begin{pmatrix} \dot{x}_1(t)\\ \dot{\lambda}(t) \end{pmatrix} = \begin{pmatrix} A_1 & B_1\\ -C_1 & -\hat{D} \end{pmatrix} \begin{pmatrix} x_1(t)\\ \lambda(t) \end{pmatrix} + \begin{pmatrix} 0\\ I_m \end{pmatrix} w(t) \\ (b) & 0 \le \lambda(t) \perp w(t) = C_1 x_1(t) + C_2 x_2(t) + C_3 z(t) + D\lambda(t) \ge 0 \\ (c) & x_2(t) = -B_3 \lambda(t), \ z(t) = -B_2 \lambda(t) - B_3 \dot{\lambda}(t). \end{cases}$$

It is noteworthy that this dynamics is obtained from the original one by treating the multiplier λ as a 310 state variable. The "equivalence" between both systems is to be understood first as follows [32, Lemma 311

3]: 312

 $\{(x_1, \lambda, w) \mid (x_1, x_2, \lambda, w) \text{ is a solution of } (2.2) (a, b, c)\} = \{(x_1, \lambda, w) \mid (x_1, \lambda, w) \text{ is a solution of } (5.1) (a)\}$ 313 and $B_3\lambda$ is absolutely continuous}. 314

By solution of (2.2) it is meant that (x_1, x_2) is absolutely continuous and the dynamics holds Lebesgue 315 almost everywhere [32, Definition 3]. The initial conditions are chosen compatible with the constraints. 316 It is noteworthy that the variables z and λ are not shown to be absolutely continuous, and examples in 317 sections 6.1 and 6.3 will illustrate this fact. One has $w(t) \in -\mathcal{N}_{\mathbb{R}^m_+}(\lambda(t))$, therefore one obtains the 318 differential inclusion (DI): 319

(5.2)
$$\begin{pmatrix} I_{n_{1}} & 0\\ 0 & B_{3}^{\top}X_{3}B_{3} \end{pmatrix} \begin{pmatrix} \dot{x}_{1}(t)\\ \dot{\lambda}(t) \end{pmatrix} - \begin{pmatrix} A_{1} & B_{1}\\ -C_{1} & -\hat{D} \end{pmatrix} \begin{pmatrix} x_{1}(t)\\ \lambda(t) \end{pmatrix} \quad \in -\begin{pmatrix} 0\\ I_{m} \end{pmatrix} \mathcal{N}_{\mathbb{R}^{m}_{+}} \left(\lambda(t) \right) \\ = -\begin{pmatrix} 0\\ I_{m} \end{pmatrix} \mathcal{N}_{\mathbb{R}^{m}_{+}} \left((0 \ I_{m}) \begin{pmatrix} x_{1}(t)\\ \lambda(t) \end{pmatrix} \right) \\ = -\mathcal{N}_{\Phi} \left(\begin{pmatrix} x_{1}(t)\\ \lambda(t) \end{pmatrix} \right)$$

with $\Phi \stackrel{\Delta}{=} \{(x_1, \lambda) \in \mathbb{R}^{n_1} \times \mathbb{R}^m \mid \lambda \ge 0\}$. If this DI has a solution $(x_1(\cdot), \lambda(\cdot))$, then $x_2(\cdot), z(\cdot)$, and $w(\cdot)$ 320 can be obtained (provided indeed that $B_3\lambda(\cdot)$ is absolutely continuous). Multiplying the first line in 321

(5.2) by X_1 and adding an exogenous excitation G(t) to the left-hand side, one obtains 322

(5.3)
$$\underbrace{\begin{pmatrix} X_1 & 0\\ 0 & B_3^\top X_3 B_3 \end{pmatrix}}_{\triangleq \mathcal{P}} \begin{pmatrix} \dot{x}_1(t)\\ \dot{\lambda}(t) \end{pmatrix} - \underbrace{\begin{pmatrix} X_1 A_1 & X_1 B_1\\ -C_1 & -\hat{D} \end{pmatrix}}_{\triangleq \mathcal{A}} \begin{pmatrix} x_1(t)\\ \lambda(t) \end{pmatrix} + G(t) \in -\mathcal{N}_{\Phi} \left(\begin{pmatrix} x_1(t)\\ \lambda(t) \end{pmatrix} \right)$$

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which is similar to the form obtained in [32, Theorem 3]. Let us assume that an exogenous input $W(x) = (x + y)^T + (x + y)^T$

³²⁴ $V(t) = (v_{(a)}(t), v_{(b)}(t), v_{(c)}(t), v_{(d)}(t))^{\top}$ is added to (2.2), then:

(5.4)
$$G(t) = \begin{pmatrix} v_{(a)}(t) \\ C_3 \left(v_{(b)}(t) + \dot{v}_{(c)}(t) \right) + C_2 v_{(c)}(t) - v_{(d)}(t) \end{pmatrix}.$$

Notice that $B_3^{\top}X_3B_3 \ge 0$, therefore $\mathcal{P} = \mathcal{P}^{\top} \ge 0$. However \mathcal{P} has a specific block-diagonal structure. Under strong SPRness and minimality, $\mathcal{A} < 0$, since $\mathcal{A} + \mathcal{A}^{\top} < 0$ from the KYP Lemma [25] (then A_1 is a Hurwitz matrix, $X_1 = X_1^{\top} > 0$ and $\hat{D} > 0$). The exogeneous term G(t) stems from the exogeneous terms E(t) and F(t) in (2.1) and their time-derivatives (see (5.4) and (6.4) for an example).

In a second step the "equivalence" between both systems is to be understood as follows [32, Theorem 330
 3]:

 $\{(x_1, \lambda) \mid (x_1, x_2, \lambda, w) \text{ is a solution of } (2.2) \ (a, b, c, d)\} = \{(x_1, \lambda) \mid (x_1, \lambda) \text{ is a solution of } (5.3) \text{ and} \\ B_3\lambda \ (= -x_2) \text{ is absolutely continuous}\}.$

As will become clear later, the results that are obtained using this approach, and the approach in the

foregoing section, are not the same (especially concerning the convergence of the piecewise-linear

335 discrete solutions).

336 5.1 THE TIME-DISCRETIZATION AND THE OSNSP

As alluded to in (5.3), we can enlarge the analysis by considering time-varying terms, so as to embed

(5.3) into the DI: $\mathcal{P}\dot{\xi}(t) - \mathcal{R}\xi(t) + G(t) \in -\mathcal{N}_{\Phi}(\xi(t))$. The Euler implicit scheme associated with (5.3) is:

(5.5)
$$\mathcal{P}\xi_{k+1} - \underbrace{(\mathcal{P} + h\mathcal{A})}_{\triangleq \mathcal{P}_h} \xi_k + hG_k \in -\mathcal{N}_{\Phi}(\xi_{k+1}),$$

It is noteworthy that this scheme is not equal to the implicit scheme in (3.1), as can be shown after few manipulations, see Appendix H. Though the difference between both approaches is not analysed further in this paper, it will be tackled numerically through an example in section 6.3. In addition, a slightly more general convex set than in (5.5) is considered, for the sake of generality. It corresponds to changing the complementarity conditions in (5.1) (b) to cone complementarity conditions $K^* \ni \lambda(t) \perp w(t) \in K$, $K \subseteq \mathbb{R}^m$ a nonempty closed convex cone, similarly to what is done in Proposition 3.2. Thus $\Phi = \mathbb{R}^{n_1} \times K^*$.

Proposition 5.1. Let h > 0 and $k \ge 0$ be given. Let us consider the DI in (5.3) and its discretization in (5.5), with $\Phi = \mathbb{R}^{n_1} \times K^*$, with K a nonempty closed convex cone. Let us denote $q_k \stackrel{\Delta}{=} -\mathcal{P}_h \xi_k + hG_k$. Let us consider the set:

(5.6)
$$\mathcal{S} \stackrel{\Delta}{=} \{ (x_1, \lambda) \in \mathbb{R}^{n_1} \times \mathbb{R}^m \mid x_1 \in \operatorname{Ker}(X_1), \ \lambda \in \operatorname{Ker}(B_3^\top X_3 B_3) \cap (B_3^\top X_3 B_3)^{-1}(K) \cap K^\star \}$$

349 Then:

1. If
$$S = \{0\}$$
, for each $q_k \in \mathbb{R}^{n_1 \times m}$ the problem $VI(\mathcal{P}, q_k, \Psi_{\Phi})$ has at least one solution.

2. Assume that $S \neq \{0\}$. If there exists $\xi^0 \in \text{dom}(\Phi)$ such that $\langle q_k - \mathcal{P}^\top \xi^0, v \rangle > 0$ for all $v \in S, v \neq 0$, then the problem $VI(\mathcal{P}, q_k, \Psi_{\Phi})$ has at least one solution. In case $X_1 = X_1^\top$, one can take $\xi_0 = 0$.

353
3. If
$$\xi_{k+1}^1$$
 and ξ_{k+1}^2 are two solutions of $VI(\mathcal{P}, q_k, \psi_{\Phi})$, then $\xi_{k+1}^1 - \xi_{k+1}^2 \in Ker(\mathcal{P})$, hence $x_{1,k+1}$ is unique.
354
Moreover $\langle q_k, \begin{pmatrix} 0\\ \lambda_{k+1}^1 - \lambda_{k+1}^2 \end{pmatrix} \rangle = 0.$

- 4. If $X_1 > 0$ and $B_3^{\top} X_3 B_3 > 0$, then the problem $VI(\mathcal{P}, q_k, \Psi_{\Phi})$ always has a unique solution.
- 5. Assume that the quintuple (P, A, B, C, D) is minimal and strongly SPR, so that the quadruple (A_1, B_1, C_1, \hat{D}) is minimal and strongly SPR. If $B_3^T X_3 B_3 > 0$, then the problem VI $(\mathcal{P}, q_k, \Psi_{\Phi})$ always has a unique solution.

6. Let
$$K = K^{\star} = \mathbb{R}^{m}_{+}$$
. Then $S = \{(x_{1}, \lambda) \mid x_{1} \in \operatorname{Ker}(X_{1}), \lambda \in \mathbb{R}^{m}_{+} \cap \operatorname{Ker}(B_{3}^{\top}X_{3}B_{3})\}$.

Proof: From Proposition **B.1** d) one has $(K^{\star})_{\infty} = K^{\star}$ and $\Phi_{\infty} = \Phi$. Also $\varphi = \Psi_{\Phi}$ hence from 360 Proposition **B.1** c) and d), $\varphi_{\infty} = \Psi_{\mathbb{R}^{n_1} \times K^{\star}}$ and $(\operatorname{dom}(\varphi_{\infty}))^{\star} = \Phi^{\star} = \mathbb{R}^{n_1} \times K$. Then one has $S \stackrel{\Delta}{=} \Phi \cap$ 361 $\operatorname{Ker}(\mathcal{P}) \cap \{\xi \in \mathbb{R}^{n_1+m} \mid \mathcal{P}\begin{pmatrix} x_1 \\ \lambda \end{pmatrix} = \begin{pmatrix} 0 \\ \alpha \end{pmatrix}, \ \alpha \in K\} = \Phi \cap \operatorname{Ker}(\mathcal{P}) \cap \{\xi \in \mathbb{R}^{n_1+m} \mid x_1 \in \operatorname{Ker}(X_1), \ B_3^{\mathsf{T}}X_3B_3\lambda \in \mathbb{R}^{n_1+m} \mid x_1 \in \mathbb{R}^{$ 362 *K*}, where we used the block diagonal structure of \mathcal{P} . Therefore $\mathcal{S} = \{(x_1, \lambda) \in \mathbb{R}^{n_1} \times \mathbb{R}^m \mid x_1 \in \mathbb{R}^{n_1} \times \mathbb{R}^n$ 363 $\operatorname{Ker}(X_1), B_3^{\mathsf{T}}X_3B_3\lambda \in K, \ \lambda \in \operatorname{Ker}(B_3^{\mathsf{T}}X_3B_3), \ \lambda \in K^{\star}\}.$ In item 2) one uses that $\Psi_{\Phi}(v) = 0$ when $v \in S$, 364 and item b') in Proposition B.2. Let us now deal with item 4. In this case $\mathcal{P} = \mathcal{P}^{\top} > 0$ and the result 365 follows from items 1 and 3. Let us define $\mathcal{R} = \mathcal{R}^{\top} > 0$, $\mathcal{R}^2 = \mathcal{P}$, and $\zeta_k = \mathcal{R}\xi_k$. Then the VI in (5.5) is 366 equivalently rewritten as: 367

(5.7)
$$\zeta_{k+1} - \mathcal{R}^{-1}\mathcal{P}_h \mathcal{R}^{-1}\zeta_k + \mathcal{R}^{-1}G_k \in -\mathcal{R}^{-1}\mathcal{N}_{\Phi}(\mathcal{R}^{-1}\zeta_{k+1}) = -\mathcal{N}_{\tilde{\Phi}}(\zeta_{k+1})$$

where the last equality stems from the chain rule, with $\tilde{\Phi} = \{z \in \mathbb{R}^{n_1+m} \mid \mathcal{R}^{-1}z \in \Phi\}$. One infers from (5.7) that $\zeta_{k+1} = \operatorname{proj}[\tilde{\Phi}; \mathcal{R}^{-1}\mathcal{P}_h\mathcal{R}^{-1}\zeta_k - \mathcal{R}^{-1}G_k] \Rightarrow \zeta_{k+1} = \mathcal{R}^{-1}\operatorname{proj}[\tilde{\Phi}; \mathcal{R}^{-1}\mathcal{P}_h\zeta_k - \mathcal{R}^{-1}G_k]$. The first assertion in item 5 about strong SPRness follows from [28, Theorem 5.4 (A')], which guarantees that $X_1 > 0$ (see Appendix **D**). The last assertion of item 5 is a consequence of item 4. Consider finally item 6. One has $(B_3^{\top}X_3B_3)^{-1}(K) = \{\lambda \in \mathbb{R}^m \mid B_3^{\top}X_3B_3\lambda = \alpha, \alpha_i \ge 0, 1 \le i \le m\}$, see *e.g.*, [53, Example 15]. Hence $\operatorname{Ker}(B_3^{\top}X_3B_3) \cap (B_3^{\top}X_3B_3)^{-1}(K) = \operatorname{Ker}(B_3^{\top}X_3B_3)$.

In item 2, the sufficient condition $\langle q_k, v \rangle > 0$ for all $v \in S$ is equivalent to $q_k \in S^* \setminus S^{\perp}$. One notices

that under the conditions of item 4 in Proposition 5.1, the DI in (5.3) is equivalent to a DI with maximal monotone set-valued term:

(5.8)
$$\dot{\zeta}(t) - \mathcal{R}^{-1}\mathcal{A}\mathcal{R}^{-1}\zeta(t) + \mathcal{R}^{-1}G(t) \in -\mathcal{N}_{\tilde{\Phi}}(\zeta(t))$$

- with $\zeta = \mathcal{R}\xi$. This fits within the class of DIs analysed for instance in [14]. Under basic assumption on the regularity of $G(\cdot)$, (5.8) has Lipschitz continuous solutions with uniqueness.
- Proposition 5.2. Assume that $X_1 > 0$, and $G(\cdot)$ is bounded. The problem $VI(\mathcal{P}, q_k, \Psi_{\Phi})$ always has at least one solution for bounded ξ_k .

Proof: Consider item 2 in Proposition 5.1, with $\xi^0 = \begin{pmatrix} x^0 \\ \lambda^0 \end{pmatrix}$, $v = \begin{pmatrix} v_x \\ v_\lambda \end{pmatrix}$. Then:

(5.9)

$$\langle q_k - \mathcal{P}^{\top} \xi^0, v \rangle = \langle -\mathcal{P}_h \begin{pmatrix} x_{1,k} \\ \lambda_k \end{pmatrix} + hG_k, v \rangle - \langle B_3^{\top} X_3 B_3 \lambda^0, v_{\lambda} \rangle
- \langle X_1 x^0, v_x \rangle
= \langle - \begin{pmatrix} X_1 + hX_1 A_1 & hX_1 B_1 \\ -hC_1 & B_3^{\top} X_3 B_3 - h\hat{D} \end{pmatrix} \begin{pmatrix} x_{1,k} \\ \lambda_k \end{pmatrix} + hG_k, v \rangle
- \langle B_3^{\top} X_3 B_3 \lambda^0, v_{\lambda} \rangle - \langle X_1 x^0, v_x \rangle.$$

382 Since $\xi^0 \in \Phi$ it follows that $x^0 \in \mathbb{R}^{n_1}$ while $\lambda^0 \in \mathbb{R}^m_+$, and similarly $v_x \in \mathbb{R}^{n_1}$, $v_\lambda \in \mathbb{R}^m_+$. Since $v \neq 0$ let

us take $x^0 = -\alpha X_1^{-1} \frac{v_x}{\|v_x\|} - \beta X_1^{-1} \frac{v_x}{\|v_x\|^2} \|v_\lambda\|$ for some $\alpha \ge \beta \ge 0$, and $\lambda^0 = 0$. It follows that:

$$\langle q_{k} - \mathcal{P}^{\mathsf{T}} \xi^{0}, v \rangle = - \langle \begin{pmatrix} X_{1} + hX_{1}A_{1} & hX_{1}B_{1} \\ -hC_{1} & B_{3}^{\mathsf{T}}X_{3}B_{3} - h\hat{D} \end{pmatrix} \begin{pmatrix} x_{1,k} \\ \lambda_{k} \end{pmatrix}, v \rangle + \langle hG_{k}, v \rangle$$

$$\geq - \| \begin{pmatrix} X_{1} + hX_{1}A_{1} & hX_{1}B_{1} \\ -hC_{1} & B_{3}^{\mathsf{T}}X_{3}B_{3} - h\hat{D} \end{pmatrix} \begin{pmatrix} x_{1,k} \\ \lambda_{k} \end{pmatrix} \| \|v\| - \|hG_{k}\| \|v\|$$

$$+ \alpha \|v_{x}\| + \beta \|v_{\lambda}\|$$

$$\geq - \| \begin{pmatrix} X_{1} + hX_{1}A_{1} & hX_{1}B_{1} \\ -hC_{1} & B_{3}^{\mathsf{T}}X_{3}B_{3} - h\hat{D} \end{pmatrix} \begin{pmatrix} x_{1,k} \\ \lambda_{k} \end{pmatrix} \| (\|v_{x}\| + \|v_{\lambda}\|)$$

$$- \|hG_{k}\| (\|v_{x}\| + \|v_{\lambda}\|) + \alpha \|v_{x}\| + \beta \|v_{\lambda}\|.$$

Since $G(\cdot)$ is bounded, for large enough α and β it follows that $\langle q_k - \mathcal{P}^\top \xi^0, v \rangle > 0$ and this holds for any $v \neq 0$ inside S.

386 5.2 BOUNDEDNESS AND CONVERGENCE ANALYSIS

In the next analysis we assume that $B_3^{\top}X_3B_3$ is singular, which is the most general case (see the example

in section 6.2, see also section 5.3 for the invertible case), and that $K = \mathbb{R}^m_+$ and that $X_1 > 0$. From the

structure of \mathcal{P} and of the set-valued right-hand side, it follows from (5.5) that $x_{1,k+1}$ is calculable from $x_{1,k}$ and λ_k :

$$\begin{aligned} X_1 x_{1,k+1} &= (X_1 + hX_1A_1) x_{1,k} + hX_1B_1\lambda_k - hG_k^x \\ \Leftrightarrow & x_{1,k+1} &= (I_{n_1} + hA_1) x_{1,k} + hB_1\lambda_k - hX_1^{-1}G_k^x, \text{ for all } k \ge 0. \end{aligned}$$

³⁹¹ where $G_k = \begin{pmatrix} G_k^x \\ G_k^z \end{pmatrix}$.

(5.11)

Lemma 5.3. From (5.11), for all $k \ge 0$ one has:

(5.12)
$$x_{1,k+1} = (I_{n_1} + hA_1)^{k+1} x_{1,0} + h \sum_{i=0}^k (I_{n_1} + hA_1)^{k-i} B_1 \lambda_i - h \sum_{i=0}^k (I_{n_1} + hA_1)^{k-i} X_1^{-1} G_i^x.$$

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Proof: The relation is true for k = 0. Let it be true for $k \ge 1$. Then using (5.11) it follows that:

$$\begin{aligned} x_{1,k+2} &= (I_{n_1} + hA_1)x_{1,k+1} + hB_1\lambda_{k+1} - hX_1^{-1}G_{k+1}^x \\ &= (I_{n_1} + hA_1)\{(I_{n_1} + hA_1)^{k+1}x_{1,0} + h\sum_{i=0}^k (I_{n_1} + hA_1)^{k-i}B_1\lambda_i \\ &-h\sum_{i=0}^k (I_{n_1} + hA_1)^{k-i}X_1^{-1}G_i^x\} + hB_1\lambda_{k+1} - hX_1^{-1}G_{k+1}^x \\ &= (I_{n_1} + hA_1)^{k+2}x_{1,0} + h\sum_{i=0}^{k+1} (I_{n_1} + hA_1)^{k+1-i}B_1\lambda_i \\ &-h\sum_{i=0}^{k+1} (I_{n_1} + hA_1)^{k+1-i}X_1^{-1}G_i^x. \end{aligned}$$

By induction the relation in (5.12) is true for all $k \ge 0$.

³⁹⁶ Using (5.5) we can also express λ_{k+1} as the solution of a generalised equation (GE):

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for all $k \ge 0$. The second formulation shows the iterative feature of this GE for λ_k , which makes it quite different from the LCCP (3.3). It is equivalent to the LCP $(B_3^{\top}X_3B_3, q_k^{\lambda}(\lambda_k, x_{1,k}, G_k^{\lambda}))$. Under the conditions stated in Proposition 5.1 and Corollary 4.1, the set of solutions to this LCP is non empty, and it contains a least-norm element $\lambda_{k+1}^{\min}, k \ge 0$, satisfying $||\lambda_{k+1}^{\min}(\lambda_k, x_{1,k}, G_k^{\lambda})|| \le \eta ||q_k^{\lambda}(\lambda_k, x_{1,k}, G_k^{\lambda})||, k \ge 0$, where $\eta > 0$ depends only on $B_3^{\top}X_3B_3$. Therefore $||\lambda_{k+1}^{\min}(\lambda_k, x_{1,k}, G_k^{\lambda})|| \le \eta ||h\hat{D} - B_3^{\top}X_3B_3|| ||\lambda_k^{\min}|| + h\eta ||G_k^{\lambda}||$. Let us denote $h_3 \stackrel{\Delta}{=} ||h\hat{D} - B_3^{\top}X_3B_3||$.

Lemma 5.4. One has for all $k \ge 0$:

(5.15)
$$\begin{aligned} ||\lambda_{k+2}^{\min}(\lambda_{k+1}, x_{1,k+1}, G_{k+1}^{\lambda})|| &\leq (\eta h_3)^{k+1} ||\lambda_1^{\min}(\lambda_0, x_{1,0}, G_0^{\lambda})|| \\ &+ h\eta \sum_{i=0}^k (\eta h_3)^{k-i} (||C_1 x_{1,i+1}|| + ||G_{i+1}^{\lambda}||) \end{aligned}$$

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Proof: The inequality is satisfied for k = 0. Assume that it is true for some $k \ge 1$. Then (5.16)

$$\begin{split} ||\lambda_{k+3}^{\min}(\lambda_{k+2}, x_{1,k+2}, G_{k+2}^{\lambda})|| &\leq \eta ||h\hat{D} - B_{3}^{\top}X_{3}B_{3}|| \; ||\lambda_{k+2}^{\min}|| + h\eta ||C_{1}x_{1,k+2}|| + h\eta ||G_{k+2}^{\lambda}|| \\ &\leq \eta ||h\hat{D} - B_{3}^{\top}X_{3}B_{3}|| \left\{ \eta^{k+1}||h\hat{D} - B_{3}^{\top}X_{3}B_{3}||^{k+1}||\lambda_{1}^{\min}(\lambda_{0}, x_{1,0}, G_{0}^{\lambda})|| \\ &+ h\eta \sum_{i=0}^{k} \eta^{k-i}||h\hat{D} - B_{3}^{\top}X_{3}B_{3}||^{k-i}(||C_{1}x_{1,i+1}|| + ||G_{i+1}^{\lambda}||) \right\} \\ &+ h\eta ||C_{1}x_{1,k+2}|| + h\eta ||G_{k+2}^{\lambda}|| \\ &= \eta^{k+2}||h\hat{D} - B_{3}^{\top}X_{3}B_{3}||^{k+2}||\lambda_{1}^{\min}(\lambda_{0}, x_{1,0}, G_{0}^{\lambda})|| \\ &+ h\eta \sum_{i=0}^{k} \eta^{k+1-i}||h\hat{D} - B_{3}^{\top}X_{3}B_{3}||^{k+1-i}(||C_{1}x_{1,i+1}|| + ||G_{i+1}^{\lambda}||) \\ &+ h\eta ||C_{1}x_{1,k+2}|| + h\eta ||G_{k+2}^{\lambda}|| \\ &= \eta^{k+2}||h\hat{D} - _{3}^{\top}X_{3}B_{3}||^{k+2}||\lambda_{1}^{\min}(\lambda_{0}, x_{1,0}, G_{0}^{\lambda})|| \\ &+ h\eta \sum_{i=0}^{k} \eta^{k+1-i}||h\hat{D} - B_{3}^{\top}X_{3}B_{3}||^{k+1-i}(||C_{1}x_{1,i+1}|| + ||G_{i+1}^{\lambda}||). \end{split}$$

406 The result follows by induction. \boxtimes

⁴⁰⁷ The next step is to use Lemma 5.4 and Lemma 5.3 to get a new upper bound. Let us denote $h_1 \stackrel{\Delta}{=} ||I_{n_1} + hA_1||$.

408 Then using (5.12) and (5.15) one has:

(5.17)
$$\begin{cases} ||x_{1,k+1}|| \le h_1^{k+1}||x_{1,0}|| + h\sum_{i=0}^k h_1^{k-i}||B_1|| \ ||\lambda_{i+1}^{\min}|| + h\sum_{i=0}^k h_1^{k-i}||X_1^{-1}|| \ ||G_i^x| \\ ||\lambda_{k+2}^{\min}|| \le \eta^{k+1}h_3^{k+1}||\lambda_1^{\min}|| + h\eta\sum_{i=0}^k \eta^{k-i}h_3^{k-i}(||C_1x_{1,i+1}|| + ||G_{i+1}^{\lambda}||) \end{cases}$$

409 Therefore we obtain:

$$\begin{aligned} ||x_{,1,k+1}|| &\leq h_1^{k+1}||x_{1,0}|| + h\sum_{i=0}^k h_1^{k-i}||X_1^{-1}|| ||G_i^x|| \\ &+ h\sum_{i=0}^k h_1^{k-i}||B_1||\eta^i h_3^i||\lambda_1^{\min}|| + hh_1^k||B_1|| ||\lambda_1^{\min}|| \\ &+ h\sum_{i=1}^k h_1^{k-i}||B_1|| h\eta \sum_{j=0}^{i-1} \eta^{i-1-j} h_3^{i-1-j}||C_1||||G_{j+1}^{\lambda}|| \\ &+ h\sum_{i=1}^k h_1^{k-i}||B_1||h\eta \sum_{j=0}^{i-1} \eta^{i-1-j} h_3^{i-1-j}||C_1|| ||x_{1,j+1}||, \end{aligned}$$

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410 which we rewrite as:

(5.19)
$$||x_{1,k+1}|| \leq h_1^{k+1} ||x_{1,0}|| + h f_k + h^2 \eta ||C_1|| ||B_1|| \sum_{i=1}^k h_1^{k-i} \sum_{j=0}^{i-1} (\eta h_3)^{i-1-j} ||x_{1,j+1}||$$

411 with

(5.20)

$$\begin{aligned} f_k &= \sum_{i=0}^k h_1^{k-i} ||X_1^{-1}|| \ ||G_i^x|| + ||B_1|| \ ||\lambda_1^{\min}|| \ \sum_{i=0}^k h_1^{k-i} (\eta h_3)^i + h_1^k ||B_1|| \ ||\lambda_1^{\min}|| \\ &+ h\eta ||B_1|| \ ||C_1|| \ \sum_{i=1}^k h_1^{k-i} \ \sum_{j=0}^{i-1} (\eta h_3)^{i-1-j} \ ||G_{j+1}^{\lambda}||. \end{aligned}$$

The third term in the right-hand side of (5.19) involves terms from $||x_{1,1}||$ to $||x_{1,k}||$. The inequality (5.19) can be rewritten as:

(5.21)
$$||x_{1,k+1}|| \le h_1^{k+1} ||x_{1,0}|| + h f_k + h^2 \sum_{i=1}^k \beta_i ||x_{1,i}||,$$

for some positive sequence of real numbers $\{\beta_i\}$, each β_i being calculated from the expression in (5.19) $\eta ||C_1|| ||B_1|| \sum_{i=1}^k h_1^{k-i} \sum_{j=0}^{i-1} (\eta h_3)^{i-1-j} ||x_{1,j+1}||$, and thus being made of powers of h_1 and of $h_3\eta$. Calculations show that (see Appendix F)

(5.22)
$$\beta_i = \eta ||C_1|| ||B_1|| \sum_{i=i}^k h_1^{k-j} (\eta h_3)^{j-i}, \quad i \ge 1$$

Let us now apply a version of the discrete-time Gronwall inequality, due to Bachpatte [13] [47, Lemma

⁴¹⁸ 100], see section E. It is deduced that for all $k \ge 0$:

(5.23)
$$||x_{1,k+1}|| \le h_1^{k+1} ||x_{1,0}|| + h f_k + h^2 \sum_{i=0}^k (h_1^{i+1}\beta_i ||x_{1,0}|| + h\beta_i f_i) \prod_{j=i+1}^{k+1} (1 + h^2\beta_j).$$

Let us state the next assumption (recall that $h_3 = ||h\hat{D} - B_3^\top X_3 B_3||$ and $h_1 = ||I_{n_1} + hA_1||$).

- Assumption 4. One has a) $h_1 \le 1$ for h > 0 small enough, b) ||G(t)|| is uniformly bounded by $G_{\max} < +\infty$ on [0, T] for any T > 0, and $G(\cdot)$ is Lipschitz continuous with Lipschitz constant $k_G \ge 0$, c) $\eta h_3 \le 1$ for h > 0 small enough.
- Inequality a) is similar to Assumption 1. As illustrated through examples in section 6, it may hold in practice for reasonable values of the time step. See also the developments in Appendix A for conditions on its validity, and Remark 5.7 for its relaxation. Assumption 4 b) is a standard property of exogenous time-dependent terms. The inequality c) is the central assumption of Assumption 4, which quantifies the couplings between the DAE part and the complementarity part of the LDCS. Some examples of circuits in section 6 will prove that it can be satisfied in practice. Notice from (5.14) that if $B_3^T X_3 B_3 = 0$ (Assumption 2), then the minimal-norm element of the set of solutions to the LCP is always $\lambda_k^{\min} = 0$, hence n = 0 in that case. Using (5.20) the next upper bound can be found:

hence $\eta = 0$ in that case. Using (5.20) the next upper bound can be found:

431 Since $h = \frac{T}{N}$ and $k \in \{0, N - 1\}$:

$$(5.25) \quad hf_k \le T \frac{(k+1)}{N} (G_{\max}||X_1^{-1}|| + ||B_1|| ||\lambda_1^{\min}||) + \frac{T}{N} ||B_1|| ||\lambda_1^{\min}|| + \frac{T^2}{n^2} \eta ||B_1|| ||C_1|| G_{\max} \sum_{i=1}^k \sum_{j=0}^{i-1} j d_{\max} \sum_{j=0}^k d_{\max} \sum_{i=1}^k \sum_{j=0}^k d_{\max} \sum_{i=1}^k \sum_{j=0}^k d_{\max} \sum_{i=1}^k \sum_{j=0}^k d_{\max} \sum_{i=1}^k \sum_{j=0}^k d_{\max} \sum_{j=0}^$$

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- ⁴³² Now we have that $\sum_{i=1}^{k} \sum_{j=0}^{i-1} j = (k-1) + 2(k-1) + 3(k-3) + \ldots + (k-3)3 + (k-2)2 + (k-1)$, where
- the sum consists of k 1 terms. Hence $\sum_{i=1}^{k} \sum_{j=0}^{i-1} j = O(k^2)$, and $\frac{T^2}{N^2} \eta ||B_1|| ||C_1|| G_{\max} \sum_{i=1}^{k} \sum_{j=0}^{i-1} j = O(k^2)$
- ₄₃₄ $T^2 \eta ||B_1|| ||C_1|| G_{\max} \frac{O(k^2)}{N^2}$. Since $k \leq N$ we infer that this term is uniformly bounded in k, for each
- fixed bounded *T*. The same holds for $T\frac{(k+1)}{N}(G_{\max}||X_1^{-1}|| + ||B_1|| ||\lambda_1^{\min}||)$. Thus $hf_k \le \alpha_1 \max(T, T^2)$
- for some α_1 which depends on η , $||B_1||$, $||C_1||$, G_{max} , $||\lambda_1^{\min}||$.
- ⁴³⁷ The next step consists of computing an upper bound from (5.23) as:

(5.26)
$$||x_{1,k+1}|| \le \left(||x_{1,0}|| + \alpha_1 \max(T, T^2))(1 + \frac{T^2}{N^2} \sum_{i=0}^k \beta_i \prod_{j=i+1}^{k+1} (1 + \frac{T^2}{N^2} \beta_j)\right).$$

- 439 $\eta ||C_1|| ||B_1||N, i \ge 1, k \le N.$
- 440 Consequently:

(

$$\begin{aligned} ||x_{1,k+1}|| &\leq (||x_{1,0}|| + \alpha_1 \max(T, T^2))(1 + \frac{T^2}{N} \sum_{i=0}^k \eta ||C_1|| ||B_1|| \prod_{j=i+1}^{k+1} (1 + \frac{T^2}{N} \eta ||C_1|| ||B_1||)) \\ &\leq (||x_{1,0}|| + \alpha_1 \max(T, T^2))(1 + \eta ||C_1|| ||B_1|| \frac{T^2}{N} \sum_{i=0}^k \prod_{j=i+1}^{k+1} (1 + \frac{T^2}{N} \eta ||C_1|| ||B_1||)) \\ &\leq (||x_{1,0}|| + \alpha_1 \max(T, T^2))(1 + \eta ||C_1|| ||B_1|| \frac{T^2}{N} \sum_{i=0}^k (1 + \frac{T^2}{N} \eta ||C_1|| ||B_1||)^{k-i}) \\ &\leq (||x_{1,0}|| + \alpha_1 \max(T, T^2))(1 + \eta ||C_1|| ||B_1|| \frac{T^2}{N} \sum_{i=0}^k (1 + \frac{\alpha_2}{N})^{k-i}) \\ &\leq (||x_{1,0}|| + \alpha_1 \max(T, T^2))(1 + \eta ||C_1|| ||B_1|| \frac{T^2}{N} \sum_{i=0}^N (1 + \frac{\alpha_2}{N})^N) \\ &\leq (||x_{1,0}|| + \alpha_1 \max(T, T^2))(1 + \eta ||C_1|| ||B_1|| \frac{T^2}{N} (N + 1)(1 + \frac{\alpha_2}{N})^N) \\ &\leq (||x_{1,0}|| + \alpha_1 \max(T, T^2))(1 + T^2\eta ||C_1|| ||B_1|| \frac{N+1}{N} (1 + \frac{\alpha_2}{N})^N), \end{aligned}$$

where $\alpha_2 = T^2 \eta ||C_1|| ||B_1||$. One has $(1 + \frac{\alpha_2}{N})^N \to e^{\alpha_2}$ as $N \to +\infty$ [73, p.196]. Therefore the sequence { $x_{1,k}$ } is uniformly bounded over $k \in \{0, N-1\}$. It directly follows from Lemma 5.4 that the sequence { λ_{k}^{\min} } is uniformly bounded as well. It follows from (4.1) that the sequence $\{\frac{x_{1,k+1}-x_{1,k}}{h}\}_k$ is uniformly bounded also. Therefore we are ready to prove the following result. The piece-wise linear approximations are defined in (4.7), and the step approximations $x_1^{h\star}(t) = x_{1,k+1}$ if $t \in [t_k, t_{k+1})$, $\lambda_{h\star}^{\min}(t) = \lambda_k^{\min}$ if $t \in [t_k, t_{k+1})$.

Proposition 5.5. Let Assumption 4 hold true, as well as the conditions for the well-posedness of the OSNSP. Then the sequence $\{x_1^h(\cdot)\}_{h\geq 0}$ converges uniformly towards a continuous limit $x_1(\cdot)$. Moreover $x_1^{h\star}(\cdot)$ converges strongly in $\mathcal{L}_2([0,T]; \mathbb{R}^{n_1})$ towards $x_1(\cdot)$, and $\dot{x}_1^h(\cdot)$ converges weakly in $\mathcal{L}_2([0,T]; \mathbb{R}^{n_1})$ towards $\dot{x}_1(\cdot)$. Therefore $\{(B_1\lambda_{h\star}^{\min})(\cdot)\}_{h\geq 0}$ converges weakly in $\mathcal{L}_2([0,T]; \mathbb{R}^{n_1})$ towards a limit $\lambda_1(\cdot)$, and $\dot{x}_1(t) = A_1x_1(t) + G^x(t) + \lambda_1(t)$ almost everywhere on [0,T].

Proof: Both sequences $\{x_1^h(\cdot)\}$ and $\{\dot{x}_1^h(\cdot)\}$ are uniformly bounded on [0, T], the sequence $\{x_1^h(\cdot)\}$ of $C^0([0, T]; \mathbb{R}^{n_1})$ functions is thus equicontinuous. From the Ascoli-Arzela Theorem, $\{x_1^h(\cdot)\}$ stays in a compact subset of $C^0([0, T]; \mathbb{R}^{n_1})$ and indeed it converges uniformly towards a continuous limit $x_1(\cdot)$. Let us now examine the sequences of piecewise-linear functions $\{G^h(\cdot)\}_{h\geq 0}$ and $\{\dot{G}^h(\cdot)\}_{h\geq 0}$, designed similarly as (4.7). From Assumption 4 b) one infers that $G^h(t)$ is uniformly bounded on

- $[0,T], \frac{\|G_{k+1}-G_k\|}{h} \le k_E$ implies that $\dot{G}^h(t)$ is also uniformly bounded in h on [0,T], thus $\{G^h(\cdot)\}_{h\ge 0}$ is
- a sequence of Lipschitz continuous, uniformly bounded, equicontinuous functions. From the Ascoli-
- Arzela Theorem, $\{G^h(\cdot)\}$ stays in a compact subset of $C^0([0,T];\mathbb{R}^{n_1})$ and it converges uniformly

- towards the continuous limit $G(\cdot)$ as $h \to 0$. The remaining part of the proof is similar to the proof of Corollary 4.6 (adding $G^{x,h}(\cdot)$ with the required properties) and it is omitted.
- ⁴⁶² Remark 5.6. Assumption 4 c) is clearly less restrictive than Assumption 2. It still imposes some constraint
- on the matrix C_3B_3 . Examples are treated next which show that it is satisfied in some practical cases.
- Let us now prove the uniform boundedness of the sequences $\{\lambda_k^{\min}\}_k$ and $\{x_{2,k}\}_k$. To this aim we shall
- use first the second line in (5.17) and (5.27). Let us denote the right-hand side of (5.27) as $x_{1,\max}(T,\eta,x_{1,0})$.
- ⁴⁶⁶ Then from (5.17) and Assumption 4 we obtain:

(5.28)
$$\begin{aligned} \|\lambda_{k+1}^{\min}\| &\leq \|\lambda_{0}^{\min}\| + h\eta \sum_{i=0}^{k} (\|C_{1}\| x_{1,\max}(T,\eta,x_{1,0}) + \|G_{i+1}^{\lambda}\|) \\ &\leq \|\lambda_{0}^{\min}\| + h\eta \|C_{1}\| x_{1,\max}(T,\eta,x_{1,0})(k+1) + h\eta \sum_{i=0}^{k} \|G_{i+1}^{\lambda}\| \\ &\leq \|\lambda_{0}^{\min}\| + \frac{T(k+1)}{N}\eta \|C_{1}\| x_{1,\max}(T,\eta,x_{1,0}) + \eta \frac{T(k+1)}{N}G_{\max} \\ &\leq \|\lambda_{0}^{\min}\| + \frac{T(N+1)}{N}\eta \|C_{1}\| x_{1,\max}(T,\eta,x_{1,0}) + \eta \frac{T(N+1)}{N}G_{\max}. \end{aligned}$$

- ⁴⁶⁷ Therefore $\|\lambda_{k+1}^{\min}\|$ is uniformly bounded for any bounded T and initial data. The boundedness of $x_{2,k}$
- follows from the equality constraint in (4.1). A result similar to Corollary 4.6 holds, but this time verifying the complementarity conditions is less easy due to the presence of z in the complementarity
- 470 variable *w*.

Remark 5.7. [Relaxation of assumptions] Assumption 4 a) is convenient to simplify the upper bounds 471 calculations, however it is somewhat restrictive. One path to relax it, is to study the terms h_1^k another way, 472 since for any $a_1 > 0$ one has $(1 + \frac{T}{N}a_1)^N \rightarrow e^{a_1T}$ as $N \rightarrow +\infty$, while this sequence is increasing. Consider 473 first (5.20). Taking into account Assumption 4 b) and c), one sees that the first two terms in hf_k are upper 474 bounded by constants times $h \sum_{i=0}^{k} h_1^{k-i} = h \sum_{i=0}^{k} h_1^i \le hk(1+ha_1)^k$, where $h_1 \le ||I_{n_1}|| + h||A_1|| = 1 + ha_1$ for some $a_1 \ge 0$. We have $hk(1+ha_1)^k \le \frac{T}{N}k(1+\frac{Ta_1}{N})^N \to Te^{a_1T}$ as $N \to +\infty$. The third term of hf_k is upper bounded by a constant times $hh_1^k \le \frac{T}{N}(1+ha_1)^N \to 0$ as $N \to +\infty$. The last term of hf_k 475 476 477 is upper bounded by a constant times $h^2 \sum_{i=1}^k h_1^{k-i} k \leq \frac{T^2}{N^2} k^2 (1+ha_1)^k \leq \frac{T^2}{N^2} N^2 (1+ha_1)^N \to T^2 e^{Ta_1}$ as $N \to +\infty$. These first calculations allow us to upper bound the term hf_k in (5.23). The first term 478 479 $|x_{1,0}|| = (1 + ha_1)^{k+1} ||x_{1,0}|| \le (1 + ha_1)^N ||x_{1,0}|| \to ||x_{1,0}||e^{a_1T} \text{ as } N \to +\infty.$ 480 The third term in (5.23) is more complex: $h^2 \sum_{i=0}^k (h_1^{i+1}\beta_i ||x_{1,0}|| + h\beta_i f_i) \prod_{j=i+1}^{k+1} (1 + h^2\beta_j)$. Since from As-481 sumption 4 c) $\eta h_3 \leq 1$, it follows that $\beta_i \leq \eta ||C_1|| ||B_1|| \sum_{j=1}^k h_1^{k-j} \leq \eta ||C_1|| ||B_1|| \sum_{j=1}^k (1+ha_1)^{k-j}$ for some $a_1 > 0$. Thus $h^2 \beta_j \leq \eta ||C_1|| ||B_1||h^2 k(1+ha_1)^{k-1} \leq \eta ||C_1|| ||B_1|| \frac{T^2}{N^2} N(1+\frac{T}{N}a_1)^N = \eta ||C_1|| ||B_1|| \frac{T^2}{N} (1+\frac{T}{N}a_1)^N$. Hence $\prod_{j=i+1}^{k+1} (1+h^2\beta_j) \leq (1+\eta ||C_1|| ||B_1|| \frac{T^2}{N} (1+\frac{T}{N}a_1)^N \to e^{\eta ||C_1|| ||B_1||T^2}e^{a_1T}$. 482 483 484 The term $h^2 \sum_{i=0}^k \beta_i h f_i = \sum_{i=0}^k (h^2 \beta_i) (hf_i) \le N\eta ||C_1|| ||B_1|| \frac{T^2}{N} (\alpha_1 T + \alpha_2 \frac{T}{N} + \alpha_3 T^2) (1 + \frac{T}{N} a_1)^{2N} \rightarrow \eta ||C_1|| ||B_1|| T^2 (\alpha_1 T + \alpha_3 T^2) e^{2a_1 T}$, hence is bounded for any bounded T. The term $h^2 \sum_{i=0}^k h_1^{i+1} \beta_i ||x_{1,0}|| = ||x_{1,0}|| \sum_{i=0}^k h_1^{i+1} h^2 \beta_i \le \eta ||C_1|| ||B_1|| ||x_{1,0}|| \sum_{i=0}^k (1+ha_1)^k h^2 k (1+ha_1)^{k-1} \le \eta ||C_1|| ||B_1|| ||x_{1,0}|| \sum_{i=0}^N (1+ha_1)^k h^2 k (1+ha_1)^{k-1} \le \eta ||C_1|| ||B_1|| ||x_{1,0}|| \sum_{i=0}^N (1+ha_1)^k h^2 k (1+ha_1)^{k-1} \le \eta ||C_1|| ||B_1|| ||x_{1,0}|| \sum_{i=0}^N (1+ha_1)^k h^2 k (1+ha_1)^{k-1} \le \eta ||C_1|| ||B_1|| ||x_{1,0}|| \sum_{i=0}^N (1+ha_1)^k h^2 k (1+ha_1)^{k-1} \le \eta ||C_1|| ||B_1|| ||x_{1,0}|| \sum_{i=0}^N (1+ha_1)^k h^2 k (1+ha_1)^{k-1} \le \eta ||C_1|| ||B_1|| ||x_{1,0}|| \sum_{i=0}^N (1+ha_1)^k h^2 k (1+ha_1)^{k-1} \le \eta ||C_1|| ||B_1|| ||x_{1,0}|| \sum_{i=0}^N (1+ha_1)^k h^2 k (1+ha_1)^k h^2$ 485 486 487 $(ha_1)^{2N} \frac{T^2}{N^2} N \leq \eta ||C_1|| ||B_1|| ||x_{1,0}|| T^2 e^{2a_1T}$. Consequently from (5.23) it is inferred that if Assumption 488 4 b) and c) holds, then $||x_{1,k+1}||$ is bounded by a term that depends only on T, initial conditions and 489 constants. 490 491 $|B_3^{\top}X_3B_3||$. Let us assume that \hat{D} is invertible. Then $||h\hat{D} - B_3^{\top}X_3B_3|| = ||\hat{D}(h - \hat{D}^{-1}B_3^{\top}X_3B_3)|| \leq ||\hat{D}(h - \hat{D}^{-1}B_3^{\top}X_3B_3)||$ 492

- Therefore it follows that $h_3^k \leq (||\hat{D}|| ||\hat{D}^{-1}|| ||B_3^{\top}X_3B_3)||)^k \left(1 + \frac{T}{||\hat{D}^{-1}|| ||B_3^{\top}X_3B_3)||\frac{1}{N}}\right)^k$. Consequently (ηh_3)^k $\leq (\eta ||\hat{D}|| ||\hat{D}^{-1}|| ||B_3^{\top}X_3B_3)||)^k \left(1 + \frac{T}{||\hat{D}^{-1}|| ||B_3^{\top}X_3B_3)||\frac{1}{N}}\right)^k$ which converges as $N \to +\infty$ to
- ⁴⁹⁷ $(\eta ||\hat{D}|| ||\hat{D}^{-1}|| ||B_3^{\mathsf{T}}X_3B_3)||)^N e^{\frac{T}{||\hat{D}^{-1}|| ||B_3^{\mathsf{T}}X_3B_3)||}}$. A necessary condition for the limit to be bounded is

⁴⁹⁸ $\eta ||\hat{D}|| ||\hat{D}^{-1}|| ||B_3^\top X_3 B_3|| \le 1$: it is unclear whether or not this inequality is less conservative that the ⁴⁹⁹ one in Assumption 4 c).

We conclude that while the relaxation of Assumption 4 a) is reasonable, the relaxation of Assumption

 $_{501}$ 4 c) is not obvious.

502 5.3 CASE $X_1 > 0$ and $B_3^{\top} X_3 B_3 > 0$

This is the case studied in item 4 of Proposition 5.1 and may be verified in practice, see the example in 503 section 6.1. Then the DI in (5.3) is equivalently rewritten as in (5.8). Consequently one can apply the 504 implicit Euler scheme and all the existence, convergence and order results presented in [14] (see that 505 article for details on the discretization method and theoretical results), to (5.8). It is inferred that in 506 this case (5.8) possesses unique solutions in $C^0([0,T];\mathbb{R}^{n_1+m})$ with essentially bounded derivatives [14, 507 Proposition 2.6]. It is deduced that $x_2(\cdot)$ is also unique in $C^0([0, T]; \mathbb{R}^{n_1+m})$ with essentially bounded 508 derivative, while the algebraic variable $z(\cdot)$ (which depends on $\dot{\lambda}(\cdot)$, see (5.1)) is essentially bounded 509 (hence *a priori* it may be discontinuous). 510

511 6 EXAMPLES AND NUMERICAL SIMULATIONS

⁵¹² Two examples of circuits with ideal diodes, and one academic example, are developed in this section.

Let us comment that the analysis of circuits applies to hydraulics, due to the strong analogy between

⁵¹⁴ hydraulic and electrical circuits [68, 56] [52, Chapter 1], where check valves play the role of diodes.

515 6.1 FIRST EXAMPLE



Figure 1: An RLCD passive circuit.

Let us consider the circuit in Figure 1. Let us assume first that $R_1 = R_3 = 0$ Ω . Let $\xi_1 = i_2, \xi_2 = \int_0^t i_3(t) dt, \xi_3 = i_1$. The dynamics is given by:

(6.1)
$$\begin{cases} \dot{\xi}_1(t) = -\frac{R_2}{L}\xi_1(t) + \frac{1}{L}u_D + \frac{1}{L}u_1(t) \\ \dot{\xi}_2(t) = \xi_1(t) - \xi_3(t) \\ 0 = \frac{1}{C}\xi_2(t) + u_D(t) + u_2(t) \\ 0 \le u_D \perp \xi_3(t) \ge 0. \end{cases}$$

It is noteworthy that this circuit dynamics agrees with [2, Theorem 3.6] [48] about the index (indeed it

cannot be of index 1 due to the loop with a capacitive element). In our notations we have $x = (\xi_1, \xi_2)^{\top}$, $z = \xi_3, C = (0 \ 0 \ 1), B = (\frac{1}{L} \ 0 \ 1)^{\top}$. The transfer function $H(s) = \frac{z(s)}{u_D(s)} = \frac{Cs^2 + \frac{CR}{L}s + \frac{1}{L}}{s + \frac{R}{T}} = Cs + \frac{1}{Ls + R}$. It has a relative degree r = -1. It can be checked that H(s) is positive real (PR) as long as $R_2 > 0$ [25, Theorem

⁵²² 2.45], but it is not strictly positive real (SPR) (the first condition in item 3.b in [25, Theorem 2.61] does ⁵²³ not hold, despite of the fact that the second condition in the same item holds), and hence it is not

⁵²⁴ strongly SPR. Let us verify the minimality of (6.1):

525 1. rank
$$((A B) - sP B) = rank \begin{pmatrix} \frac{-R}{L} - s & 0 & 0 & \frac{1}{L} \\ 1 & -s & -1 & 0 \\ 0 & \frac{1}{C} & 0 & 1 \end{pmatrix} = 3 \text{ for all } s \in \mathbb{C}.$$

526 2. rank(*P B*) = rank $\begin{pmatrix} 1 & 0 & 0 & \frac{1}{L} \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}$ = 3.

527 3.
$$\operatorname{rank}(A^{\top} - sP^{\top} \ C^{\top}) = \operatorname{rank}\begin{pmatrix} \frac{-R}{L} - s & 1 & 0 & 0\\ 0 & -s & \frac{1}{C} & 0\\ 0 & -1 & 0 & 1 \end{pmatrix} = 3 \text{ for all } s \in \mathbb{C}.$$

528 4. rank
$$(P^{\top} \ C^{\top}) = rank \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix} = 3.$$

529 5.
$$(A \ B) \operatorname{Ker}(P) = \begin{pmatrix} \frac{-R}{L} & 0 & 0\\ 1 & 0 & -1\\ 0 & \frac{1}{C} & 0 \end{pmatrix} \operatorname{Ker}(P) = \{\xi \in \mathbb{R}^3 \mid \xi_1 = \xi_3 = 0, \xi_2 \in \mathbb{R}\} \subset \operatorname{Im}(P) = \{\xi \in \mathbb{R}^3 \mid \xi_3 = 0, \xi_2 \in \mathbb{R}, \xi_1 \in \mathbb{R}\},\$$

Therefore we infer that (6.1) admits a special Weiertrass form as (2.2). Let us choose $L = 4.7 \ 10^{-4}$ H, $R = 1\Omega$, $C = 10^{-1}$ F. Let us define $S = \begin{pmatrix} -1.414 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 10^{-4} \end{pmatrix}$, $T = \begin{pmatrix} -0.707 & 0 & 0 \\ 0 & 0 & 1 \\ -0.707 & -1 & 0 \end{pmatrix}$. Then $SAT = \begin{bmatrix} -2.12 \ 10^3 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 1 \end{pmatrix}$, $SB = \begin{pmatrix} -3 \ 10^{-3} \\ 0 \\ 10^{-4} \end{pmatrix}$, $CT = (-0.707 - 1 \ 0)$, $S\begin{pmatrix} \frac{1}{L} & 0 \\ 0 & 0 \\ 0 & 1 \end{pmatrix} = \begin{pmatrix} 0.301 \ 10^4 & 0 \\ 0 & 0 \\ 0 & 10^{-4} \end{pmatrix}$. Let us make the state variable change $\zeta = T^{-1}\zeta$, then we obtain the transformed dynamics:

(6.2)
$$\begin{cases} \dot{\zeta}_1(t) = -2.12 \ 10^3 \zeta_1(t) - 3 \ 10^{-3} \lambda(t) + 0.301 \ 10^4 u_1(t) \\ \dot{\zeta}_3(t) = \zeta_2(t) \\ 0 = \zeta_3(t) + 10^{-4} \lambda(t) + 10^{-4} u_2(t) \end{cases}$$

 $\begin{pmatrix} 0 - \zeta_3(t) + 10 & \lambda(t) + 10 & u_2(t) \\ 0 \le \lambda(t) \perp w(t) = -0.707\zeta_1(t) - \zeta_2(t) \ge 0. \end{cases}$

Let us now define $x_1 = \zeta_1$, $x_2 = \zeta_3$, $x_3 = \zeta_2$, then we obtain

(6.3)
$$\begin{cases} \dot{x}_1(t) = -2.12 \ 10^3 x_1(t) - 3 \ 10^{-3} \lambda(t) + 0.301 \ 10^4 u_1(t) \\ \dot{x}_2(t) = x_3(t) \\ 0 = x_2(t) + 10^{-4} \lambda(t) + 10^{-4} u_2(t) \\ 0 \le \lambda(t) \perp w(t) = -0.707 x_1(t) - x_3(t) \ge 0, \end{cases}$$

which is the canonical form in (2.2) with $z = x_3$, $B_1 = -3.10^{-3}$, $B_2 = 0$, $B_3 = 10^{-4}$, $C_1 = -0.707$, $C_2 = 0$, $C_3 = -1$, D = 0. We have $X_3 = 10^4$, $A_1 = -2.12 \ 10^3$, $X_1 = \frac{C_1}{B_1} > 0$, $\hat{D} = D - C_3 B_2 - C_2 B_3 = 0$, $\bar{D}_h = h^2 C_1 B_1 + B_3^{\mathsf{T}} X_3 B_3 = 2.121 \ 10^{-3} h^2 + 10^{-4}$, $h_3 = B_3^{\mathsf{T}} X_3 B_3 = 10^{-4}$, $h_1 = |1 + hA_1| = |1 - 2.12 \ 10^3 \ h|$ which verifies Assumption 4 for h > 0 small enough. Notice that the set S in Propositions 3.2 and 5.1 verifies $S = \{0\}$, hence the OSNSP is well-posed with unique solution. Provided that Assumption 4 is satisfied, the material in section 5 applies. In a similar way to (5.1), the system (6.3) is transformed into:

(6.4)
$$\begin{cases} \dot{x}_1(t) = -2.12 \ 10^3 x_1(t) - 3 \ 10^{-3} \lambda(t) + 0.301 \ 10^4 u_1(t) \\ 10^{-4} \dot{\lambda}(t) = 0.707 x_1(t) - 10^{-4} \dot{u}_2(t) + w(t) \\ x_2(t) = -10^{-4} \lambda(t) - 10^{-4} u_2(t) \\ z(t) = -10^{-4} \dot{\lambda}(t) - 10^{-4} \dot{u}_2(t) . \end{cases}$$

⁵⁴² From (5.14), the following LCP is obtained:

(6.5)
$$\begin{cases} 0 \le 10^{-4} \lambda_{k+1} + q_k \perp \lambda_{k+1} \ge 0 \\ q_k = -10^{-4} \lambda_k - 0.707 h x_{1,k} - 10^{-4} h u_{2,k} , \end{cases}$$

and $x_{1,k+1}$, $x_{2,k+1}$, z_{k+1} as in (H.1). Let us check Assumptions 1 and 4 c). Since $\hat{D} = 0$ it is easy to show that 543 $\eta = 10^4$ while $h_3 = 10^{-4}$, using the LCP in (5.14) that can be solved explicitly. Thus $\eta h_3 \leq 1$ for all $h \geq 0$. 544 However, $h_1 = ||1 - 2.12 \, 10^3 h|| \le 1$ if and only if $h < 4.7 \, 10^{-4}$. It is also noteworthy that this system 545 verifies the conditions of section 5.3, and item 4 of Proposition 5.1, since $X_1 > 0$ and $B_3^{\top}X_3B_3 = 10^{-4} > 0$. 546 Hence the state space form (5.8), that is equivalent to (5.3) since $\mathcal{P} > 0$, can be calculated. Also the 547 generalized equation (5.14) always has a unique solution λ_{k+1} . As we shall see through numerical 548 simulations, the variable $z(\cdot) = x_3(\cdot)$ may be discontinuous. Actually this is in agreement with the 549 theoretical conclusions from section 5.3, which do not state that this variable is continuous. It also shows that a complete convergence analysis, including z, should rely on other tools than those used in 551 this paper. 552

In Figure 2, the results of the numerical simulation of (6.3) with an implicit Euler scheme are depicted². Numerical simulations are performed with the INRIA software package SICONOS [1]. At each time step, the Lemke solver [65, 64] is applied to solve the LCP (3.1) resulting from the discretization scheme. The initial conditions are set to $\mathbf{x}(0) = (0.01, -10^{-4}, -7.07 \cdot 10^{-3})^{\mathrm{T}}$, and the time step is $h = 5 \cdot 10^{-6}$ s. The input functions $u_1(t)$ and $u_2(t)$ are taken as sinusoidal functions: $u_1(t) = \sin(t)$ and $u_2(t) = \cos(10t)$. Let us observe that the variable $z(t) = x_3(t)$ appears to be discontinuous at $t \approx 5.9$ s. The multiplier $\lambda(t)$ remains continuous, see Figure 3. Finally, an order of convergence equal to one is deduced from the data in Figure 4. The error formula used is:

(6.6)
$$\operatorname{error} = \|\mathbf{x}(t) - \mathbf{x}^{h}(t)\|_{\infty}.$$

561

Remark 6.1 (Modelling). One can obtain a relative degree one LCS by eliminating the variable $u_D(t)$ in (6.1). However, for the sake of generality, and since in practice equality constraints are not always

eliminated, we consider here a poorly formulated version of the circuit's dynamics, where the variables

are badly matched to the constraints. Such bad matching could be avoided using structural analysis.

566 6.2 SECOND EXAMPLE

This circuit example depicted in Figure 5 is inspired from [32]. The so-called modified nodal analysis is used to derive the dynamics [2, chapter 3]. It is assumed that the voltage source delivers a constant voltage $u(t) = u_c$, while the other voltage source $u_1(t)$ is time-varying Lipschitz continuous and bounded.

570 Notice that the current/voltage characteristic of this device can be written as $i_5(t) \in -\partial \psi_{\{u_c\}}(u(t))$, or

²Source code of the examples can be found at https://github.com/siconos/siconos-tutorials/tree/master/.sandbox/mlcspaper



Figure 2: Implicit Euler numerical simulation of (6.3) with $h = 5 \cdot 10^{-6}$ s. The numerical simulations of the variables $x_1(t)$, $x_2(t)$ and $x_3(t) = z(t)$ are depicted.



Figure 3: Implicit Euler numerical simulation of (6.3) with $h = 5 \cdot 10^{-6}$ s. The numerical simulations of the variables $\lambda(t)$ and w(t) are depicted.



Figure 4: Numerical convergence of the implicit Euler scheme applied to (6.3). The reference curve is taken for $h = 5.10^{-6}$ s. The error formula used is given in (6.6).



Figure 5: An RLCD passive circuit.

equivalently as $\mathbb{R} \ni i_5(t) \perp u(t) - u_c \in \{0\}$. The dynamics is given by:

(6.7)
$$\begin{cases} C\dot{e}(t) = -R^{-1}e(t) + i_4(t) - i_1(t) - i_5(t) \\ L\frac{di_1}{dt}(t) = e(t) + u_1(t) \\ 0 = e(t) + u(t) \\ 0 \le u_D(t) \perp i_4(t) \ge 0 \\ \mathbb{R} \ni i_5(t) \perp u(t) - u_c \in \{0\}. \end{cases}$$

The state vector is $\xi = (\xi_1, \xi_2, \xi_3)^\top = (e, i_1, i_5)^\top$, and $u_D(t) = e(t)$. The complementarity conditions can be written as $K = \mathbb{R}^+ \times \mathbb{R} \ni (u_D(t), i_5(t)) \perp (i_4(t), u(t) - u_c) \in \mathbb{R}^+ \times \{0\} = K^*$. Setting $\lambda_1 = i_4$ and

574
$$\lambda_2 = u(t) - u_c$$
, we obtain $P = \begin{pmatrix} C & 0 & 0 \\ 0 & L & 0 \\ 0 & 0 & 0 \end{pmatrix}$, $A = \begin{pmatrix} -R^{-1} & -1 & -1 \\ 1 & 0 & 0 \\ 1 & 0 & 0 \end{pmatrix}$, $B = \begin{pmatrix} 1 & 0 \\ 0 & 0 \\ 0 & 1 \end{pmatrix}$, $C = \begin{pmatrix} 1 & 0 & 0 \\ 0 & 0 & 1 \end{pmatrix}$, $D = 0$,

and the LDCS form in (2.1) is equal to:

(6.8)
$$\begin{cases} \dot{\xi}_{1}(t) = -\frac{1}{RC}\xi_{1}(t) - \frac{1}{C}\xi_{2}(t) - \frac{1}{C}\xi_{3}(t) + \frac{1}{C}\lambda_{1}(t) \\ \dot{\xi}_{2}(t) = \frac{1}{L}\xi_{1}(t) + \frac{u_{1}(t)}{L} \\ 0 = \xi_{1}(t) + \lambda_{2}(t) + u_{c} \\ 0 \le \xi_{1}(t) \pm \lambda_{1}(t) \ge 0 \\ \mathbb{R} \ni \xi_{3}(t) \pm \lambda_{2}(t) \in \{0\}. \end{cases}$$

576 Let
$$S = \begin{pmatrix} 0 & \frac{1}{L} & -\frac{1}{L} \\ -1 & 0 & -\frac{1}{R} \\ 0 & 0 & -C \end{pmatrix}, T = \begin{pmatrix} 0 & 0 & -\frac{1}{C} \\ 1 & 0 & 0 \\ -1 & 1 & 0 \end{pmatrix}, \text{then } SAT = \begin{pmatrix} 0 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 1 \end{pmatrix}, SPT = \begin{pmatrix} \frac{1}{L} & 0 & 0 \\ 0 & 0 & \frac{1}{C} \\ 0 & 0 & 0 \end{pmatrix}, CT = \begin{pmatrix} 0 & 0 & -\frac{1}{C} \\ -1 & 1 & 0 \end{pmatrix}$$

577 $SB = \begin{pmatrix} 0 & -\frac{1}{L} \\ 1 & -\frac{1}{R} \\ 0 & -C \end{pmatrix}.$ Using the state variable change $\zeta = T^{-1}\xi$ with $T^{-1} = \begin{pmatrix} 0 & 1 & 0 \\ 0 & 1 & 1 \\ -C & 0 & 0 \end{pmatrix}$, one obtains:

(6.9)
$$\begin{cases} \dot{\zeta}_{1}(t) = -(\lambda_{2}(t) + u_{c}(t)) + \frac{u_{1}(t)}{L} \\ \frac{1}{C} \dot{\zeta}_{3}(t) = \zeta_{2}(t) + \lambda_{1}(t) - \frac{1}{R}(\lambda_{2}(t) + u_{c}(t)) \\ 0 = -C(\lambda_{2}(t) + u_{c}(t)) + \zeta_{3}(t) \\ K^{\star} \ni \lambda(t) \perp w(t) = \begin{pmatrix} 0 \\ -1 \end{pmatrix} \zeta_{1}(t) + \begin{pmatrix} 0 \\ 1 \end{pmatrix} \zeta_{2}(t) + \begin{pmatrix} -\frac{1}{C} \\ 0 \end{pmatrix} \zeta_{3}(t) \in K \end{cases}$$

Let $x_1 = \zeta_1, x_2 = \frac{1}{C}\zeta_3, z = \zeta_2, B_1 = (0 - 1), B_2 = (1 - \frac{1}{R}), B_3 = (0 - 1), C_1 = \begin{pmatrix} 0 \\ -1 \end{pmatrix}, C_3 = \begin{pmatrix} 0 \\ 1 \end{pmatrix}, C_2 = \begin{pmatrix} -1 \\ 0 \end{pmatrix}, C_3 = \begin{pmatrix} 0 \\ -1 \end{pmatrix}, C_4 = \begin{pmatrix} 0 \\ -1 \end{pmatrix}, C_5 = \begin{pmatrix}$

(6.10)
$$\begin{cases} \dot{x}_1(t) = -\lambda_2(t) - u_c(t) + \frac{u_1(t)}{L} \\ \dot{x}_2(t) = z(t) + \lambda_1(t) - \frac{1}{R}(\lambda_2(t) + u_c(t)) \\ 0 = -\lambda_2(t) - u_c(t) + x_2(t) \\ K^{\star} \ni \lambda(t) \perp w(t) = \begin{pmatrix} 0 \\ -1 \end{pmatrix} x_1(t) + \begin{pmatrix} -1 \\ 0 \end{pmatrix} x_2(t) + \begin{pmatrix} 0 \\ 1 \end{pmatrix} z(t) \in K \end{cases}$$

In the simulations depicted in Figures 6a and 6b, it has been chosen $u_c(t) = \sin(10t) - 1 \le 0$, which complies with the system's intrinsic constraints. Then (when $u_1(t) = 0$) one recovers the standard form in (2.2), with a different class of sets *K*, however. One has $A_1 = 0$, $B_3^{T}X_3 = -C_3$ with $X_3 = 1$, and

 $X_1B_1 = C_1^{\top}$ (this follows necessarily from (D.2) and $Q_1 \leq 0$ and using [25, Lemma A.69]), with $X_1 = 1$ 583 (hence the quadruple (A_1, B_1, C_1, \hat{D}) is passive but it is not strictly passive and its associated transfer 584 function is not strongly SPR). Also $\hat{D} = -C_3B_2 - C_2B_3 = \begin{pmatrix} 0 & 1 \\ -1 & \frac{1}{2} \end{pmatrix} \ge 0$. Also $B_3^{\top}X_3B_3 = \begin{pmatrix} 0 & 0 \\ 0 & 1 \end{pmatrix} \ge 0$. 585 It is computed that $\bar{D}_h = \begin{pmatrix} 0 & h \\ -h & \frac{h}{R} + h^2 + 1 \end{pmatrix}$. The various sets for the set S in Proposition 3.2 are 586 calculated as: $K^* = \mathbb{R}^+ \times \{0\}$, $\operatorname{Ker}(\bar{D}_h + \bar{D}_h^{\mathrm{T}}) = \begin{pmatrix} 0 & 0 \\ 0 & 2(\frac{h}{R} + h^2 + 1) \end{pmatrix} = \mathbb{R} \times \{0\}$, $\{w \mid \bar{D}_h w \in K = \mathbb{R}^+ \times \mathbb{R}\} = \mathbb{R} \times \mathbb{R}^+$. So $S = \mathbb{R}^+ \times \{0\} \neq 0$ and item 2 of Proposition 3.2 has to be applied. It is calculated 587 588 that $\langle h\beta_k - \bar{D}_h\xi_0, v \rangle = \langle \begin{pmatrix} -h\xi_0^1 \\ -hx_{1,k} - x_{2,k} + h\xi_0^1 - (\frac{h}{2} + h^2 + 1)\xi_0^2 \end{pmatrix}$. Also $v \in K^*$ implies that $v_1 \ge 0$ and $v_2 = 0$, so it follows that $\langle h\beta_k - \bar{D}_h\xi_0, v \rangle + \psi_{K^*}(v) > 0$ is guaranteed, if $v_1 > 0$, by taking $\xi_0^1 < 0$. It is 589 590 inferred that item 2 in Proposition 3.2 is verified and the OSNSP has at least one solution. 591 Any two solutions λ_{k+1}^1 and λ_{k+1}^2 satisfy $\lambda_{2,k+1}^1 = \lambda_{2,k+1}^2$ for all $h \ge 0$, in view of the form of \overline{D}_h . Then both $x_{1,k+1}$ and $x_{2,k+1}$ are unique due to the form of B_1 and B_3 . For the convergence analysis, let us note 592 593 that this system does not satisfy Assumption 2, therefore let us use the material in section 5. Using the 594 Frobenius norm one finds that $h_3 = \sqrt{2h^2} + \left(\frac{h}{R} - 1\right)^2$, using the Hölder induced norm $|| \cdot ||_{2,2} = \sigma_{\max}(\cdot)$ 595 one finds a smaller $h_3 = \sqrt{h^2 + \left(\frac{h}{R} - 1\right)^2}$, and $0 \le h_3 \le 1$ for all $h \in [0, \frac{2R}{1+R^2}]$. The material in section 5 is 596 based on the LCP in (5.14), for which the existence of a minimal norm solution satisfying the inequalities 597 stated below (5.14) is guaranteed. However in this example $K = \mathbb{R}_+ \times \mathbb{R} \neq \mathbb{R}_+^2$. The following problem 598 has to be studied instead (it is noteworthy that the rest of the analysis remains unchanged as long as 599 the minimal norm solution can be characterised): 600 $K \ni B_3^\top X_3 B_3 \lambda_{k+1} + q_k^\lambda \perp \lambda_{k+1} \in K^\star,$ (6.11)(where q_k^{λ} is defined in (5.14)), *i.e.*, 601 $\mathbb{R}_{+} \times \{0\} \ni \begin{pmatrix} \lambda_{1,k+1} \\ \lambda_{2,k+1} \end{pmatrix} \perp \begin{pmatrix} 0 \\ \lambda_{2,k+1} \end{pmatrix} + q_{k}^{\lambda} \in \mathbb{R}_{+} \times \mathbb{R}.$

(6.12) $\mathbb{R}_{+} \times \{0\} \ni \begin{pmatrix} \lambda_{1,k+1} \\ \lambda_{2,k+1} \end{pmatrix} \perp \begin{pmatrix} 0 \\ \lambda_{2,k+1} \end{pmatrix} + q_{k}^{\lambda} \in \mathbb{R}_{+} \times \mathbb{R}.$ Therefore $\lambda_{2,k+1} = 0$, and $\lambda_{1,k+1} = 0$ is always a solution (notice that the external excitation is constrained

⁶⁰² Therefore $\lambda_{2,k+1} = 0$, and $\lambda_{1,k+1} = 0$ is always a solution (notice that the external excitation is constrained ⁶⁰³ to be non negative to comply with the system's constraints). This result agrees with the numerical ⁶⁰⁴ simulations in Figure 6b.

605 6.3 THIRD EXAMPLE

Let us consider now an academic toy-example directly in the form of (2.2), whose dynamics are given as:

(6.13)
$$\begin{cases} \dot{x}_{1}(t) = -2x_{1}(t) + 2\lambda_{1}(t) - \lambda_{2}(t) \\ \dot{x}_{2}(t) = z(t) \\ 0 = x_{2}(t) + \lambda_{1}(t) - 2\lambda_{2}(t) \\ 0 \le \begin{pmatrix} \lambda_{1}(t) \\ \lambda_{2}(t) \end{pmatrix} \perp \begin{pmatrix} 2 \\ -1 \end{pmatrix} x_{1}(t) + \begin{pmatrix} 1 \\ 3 \end{pmatrix} x_{2}(t) + \begin{pmatrix} -1 \\ 2 \end{pmatrix} z(t) + \begin{pmatrix} 2 & -1 \\ 0 & 1 \end{pmatrix} \begin{pmatrix} \lambda_{1}(t) \\ \lambda_{2}(t) \end{pmatrix} \ge 0 \end{cases}$$

It can be verified that the transfer matrix $H(s) = D + C(sP - A)^{-1}B$ is positive real [25, Definition 2.34], and the state-space realization is minimal (see section G). In fact H(s) is even strongly (or extended) positive real [25, Definition 2.78]. Thus [28, Theorem 5.4 (A') (B')] applies. It is noteworthy



Figure 6: Implicit Euler numerical simulation of (6.9) with $h = 5 \cdot 10^{-6}$ s. The solutions $x_1(t)$ are not displayed as they are independent of z(t), $\lambda_1(t)$. The variable $\lambda_2(t)$ trivially equals 0.

that since $D = \begin{pmatrix} 2 & -1 \\ 0 & 1 \end{pmatrix} > 0$, it follows that the multipliers λ are Lipschitz continuous functions of 611 the state (x_1, x_2, z) . Thus the LDCS is a DAE with Lipschitz continuous nonlinear and nonsmooth 612 right-hand side. From (6.13) it follows that: $A_1 = -2 < 0$, $C_1 = X_1 B_1^T = (2, -1)^T$ with $X_1 = 1 > 0$, $X_3 = I_2$, 613 $B_3 = (1 - 2) = -C_3^{\top}$, and $\hat{D} = \begin{pmatrix} 1 & 1 \\ -3 & 7 \end{pmatrix} > 0$ with $(\hat{D} + \hat{D}^{\mathrm{T}}) > 0$. It follows from [28, Theorem 5.4] 614 (A'), Theorem 4.6] that the quadruple (A_1, B_1, C_1, \hat{D}) is strictly passive. It can be checked that the set 615 $S = \{0\}$ in Proposition 3.2 for any h > 0, thus there exists a unique solution λ_{k+1} to the implicit Euler 616 discretization of (6.13). However $C_3 \neq 0$ and Assumption 2 does not hold, preventing the convergence 617 analysis in section 4 and the stability analysis in section 4.1 to apply. Consequently, the "equivalent" 618 system (5.1) from Section 5 is analysed: 619

(6.14)
$$\begin{cases} \begin{pmatrix} 1 & 0 & 0 \\ 0 & 1 & -2 \\ 0 & -2 & 4 \end{pmatrix} \begin{pmatrix} \dot{x}_1(t) \\ \dot{\lambda}_1(t) \\ \dot{\lambda}_2(t) \end{pmatrix} = \begin{pmatrix} -2 & 2 & -1 \\ -2 & -1 & -1 \\ 1 & 3 & -7 \end{pmatrix} \begin{pmatrix} x_1(t) \\ \lambda_1(t) \\ \lambda_2(t) \end{pmatrix} + \begin{pmatrix} 0 \\ w_1(t) \\ w_2(t) \end{pmatrix} \\ w(t) \in -\mathcal{N}_{\mathbb{R}^2}(\lambda(t)) \end{cases}$$

Let us check if Assumption 4 (c) is verified. Using the Frobenius norm, one can obtain that $0 \le h_3 = \sqrt{60h^2 + 25 - 66h} \le 1$ if $0 < h \le 0.7$. In the context of (6.14), the LCP (5.14) becomes:

$$0 \leq \begin{pmatrix} \lambda_{1,k+1} \\ \lambda_{2,k+1} \end{pmatrix} \perp \begin{pmatrix} 1 & -2 \\ -2 & 4 \end{pmatrix} \begin{pmatrix} \lambda_{1,k+1} \\ \lambda_{2,k+1} \end{pmatrix} + \begin{pmatrix} (h-1) & (h+3) \\ (2-3h) & (7h-4) \end{pmatrix} \begin{pmatrix} \lambda_{1,k} \\ \lambda_{2,k} \end{pmatrix} + \begin{pmatrix} 2hx_{1,k} \\ -hx_{1,k} \end{pmatrix} \geq 0$$
(6.15)
$$0 \leq \begin{pmatrix} \lambda_{1,k+1} \\ \lambda_{2,k+1} \end{pmatrix} \perp \begin{pmatrix} \lambda_{1,k+1} - 2\lambda_{2,k+1} \\ 4\lambda_{2,k+1} - 2\lambda_{1,k+1} \end{pmatrix} + \begin{pmatrix} q_{1,k} \\ q_{2,k} \end{pmatrix} \geq 0.$$

From Corollary 4.1, there exists a solution to (6.15). We can compute λ_{k+1} in each mode of the LCP where $\lambda_{k+1} \neq 0$:

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- 1. If $\lambda_{1,k+1} = 0$ and $\lambda_{2,k+1} \ge 0$ then $\|\lambda_{2,k+1}\| \le \frac{1}{2} \|q_{1,k}\|$.
- 625 2. If $\lambda_{2,k+1} = 0$ and $\lambda_{1,k+1} \ge 0$ then $\|\lambda_{1,k+1}\| \le \frac{1}{2} \|q_{2,k}\|$.
- 3. If $\lambda_{1,k+1} > 0$ and $\lambda_{2,k+1} > 0$ there is infinitely many solutions. Indeed, we obtain :

(6.16)
$$\begin{cases} \lambda_{1,k+1} = 2\lambda_{2,k+1} - q_{1,k} \\ 0 = 2q_{1,k} + q_{2,k} \end{cases}$$

From [16, Proposition 6.1.7], there exists a minimal solution $\lambda_{\min,k+1}$. Then, (6.16)'s particular solution $\lambda_{2,k+1} = -\lambda_{1,k+1}$ leads to $\|\lambda_{\min,k+1}\| \le (1/3)\|q_{k+1}\|$. It follows that $\eta = \frac{1}{2} < 1$ for all h > 0, and $\eta h_3 \le 1$ for all $0 < h \le 1.4$, and finally, the results from Section 5.2 hold if small enough time steps are chosen. Problem (6.14) is simulated using the discretization (6.15). The LCP (6.15) is solved using the numerical solver Lemke provided in the SICONOS Toolbox³ [1]. In Figures 7a and 7b the simulated variables (x_1, x_2, z) and (λ_1, λ_2) are displayed for time step $h = 5 \cdot 10^{-6}$ s. The initial conditions are $(x_{1,0} = 5, \lambda_{1,0} = 5, \lambda_{2,0} = 0)$ and the discrete variables $x_{2,k}$ and z_k are given by the discretization of (5.1) (c):

(6.17)
$$\begin{cases} x_{2,k} = -B_3 \lambda_k \\ z_k = -B_3 \frac{\lambda_{k+1} - \lambda_k}{h}, \end{cases}$$

with *h* the time step. Let us notice in Figure 7a that the variable *z* possesses a discontinuity. In Figure 8a we plot the error provided in (6.6), with respect to a reference solution taken for $h = 5.10^{-6}$ s. We confirm the order one convergence rate for h > 0 sufficiently small. Finally, let us also notice that the solutions of the schemes in (3.1) and in (5.5), are not exactly the same (recall that the continuous systems (2.2) and (5.1) are "equivalent" in the sense explained after (5.3), see also section H). This can

be observed in Figure 8b where the differences between the two schemes are shown for (x_1, x_2) .



Figure 7: Implicit Euler numerical simulation of (6.13) with $h = 5 \cdot 10^{-6}$ s

³https://nonsmooth.gricad-pages.univ-grenoble-alpes.fr/siconos/index.html



Figure 8: Numerical convergence studies, with error formulas (6.6), of schemes (3.1) and (5.5). The reference is computed with the scheme (5.5) and $h = 5 \cdot 10^{-6}$ s.

640 7 CONCLUSION

This article deals with the time-discretization of a class of singular linear complementarity systems 641 (descriptor variable linear complementarity systems) which satisfy a passivity constraint. The implicit 642 Euler method is analysed. First the one-step nonsmooth problem is studied, and conditions which 643 guarantee its well-posedness are given. Then the approximate solutions convergence is analysed, and 644 it is shown that under some reasonable assumptions, convergence holds. Our results are however 645 partial in the sense that we are not able to prove the convergence of all the variables. One academic 646 example and two examples of electrical circuits with ideal diodes (equivalently hydraulic circuits with 647 check valves) illustrate the developments, with numerical simulations. 648

649 APPENDIX A PRESERVATION OF POSITIVE DEFINITENESS WITH PERTURBATION

We give here an excerpt of [37, Theorem 2.11], and a corollary of it. Let us recall that for a given $M \in \mathbb{R}^{n \times n}$, $||M||_{2,2}$ is the induced matricial norm such that $||M||_{2,2} = \sigma_{\max}(M)$ (the largest singular value).

Theorem A.1. [37] Let $M \in \mathbb{R}^{n \times n}$ be a positive definite matrix. Then every matrix

$$A \in \{A \in \mathbb{R}^{n \times n} \mid \left\| \left(\frac{M + M^{\top}}{2} \right)^{-1} \right\|_{2,2} ||M - A||_{2,2} < 1 \}$$

⁶⁵³ *is positive definite.*

654 Corollary A.2. [21] Let D = P + N, where D, P and N are $n \times n$ real matrices, and P > 0, not necessarily

655 symmetric. If

.1)
$$||N||_{2,2} < \frac{1}{\|\left(\frac{P+P^{\top}}{2}\right)^{-1}\|_{2,2}}$$

656 then D > 0.

(A

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Consider the term h_1 in Assumptions 1 and 4. If $h < \frac{1}{\sigma_{\max}(A_1)}$ then $I_{n_1} + hA_1 > 0$ from Corollary A.2. If in addition $A_1 \leq (<) 0$ then $0 < I_{n_1} + hA_1 \leq (<) I_{n_1}$. Consider now $\bar{A}_1 = (I_{n_1} + hA_1)^\top (I_{n_1} + hA_1) =$ $I_{n_1} + h(A_1 + A_1^\top) + h^2 A_1^\top A_1 = I_{n_1} + h\tilde{A}_1$, with $\tilde{A}_1 = A_1 + A_1^\top + hA_1^\top A_1 = \tilde{A}_1^\top$. If $A_1 < 0 \implies A_1 + A_1^\top < 0$, then $\tilde{A}_1 < 0$ for h > 0 small enough, still using Corollary A.2. Thus $\bar{A}_1 < I_{n_1}$ for h > 0 small enough. Morover for h > 0 small enough it follows that $\bar{A}_1 > 0$ by Corollary A.2. It is deduced that there exists $h_{\max} > 0$ such that $0 < \bar{A}_1 < I_{n_1}$ for all $h \in (0, h_{\max})$. From [16, Lemma 8.4.1], $0 < \lambda_{\min}(\bar{A}_1)$ and $\lambda_{\max}(\bar{A}_1) < 1$. It is inferred that $\sigma_{\max}(I_{n_1} + hA_1) = ||I_{n_1} + hA_1||_{2,2} = h_1 < 1$.

664 APPENDIX B WELL-POSEDNESS OF VARIATIONAL INEQUALITIES

- ⁶⁶⁵ The next results use the notions of recession functions and cones, which we briefly introduce now (see
- [24, 53] for illustrating examples), [72, 53, 89]. Let $f : \mathbb{R}^n \to \mathbb{R} \cup \{+\infty\}$ be a proper convex and lower
- semi-continuous function, we denote by dom $(f) \stackrel{\Delta}{=} \{x \in \mathbb{R}^n | f(x) < +\infty\}$ the domain of the function
- 668 $f(\cdot)$. The Fenchel transform $f^{\star}(\cdot)$ of $f(\cdot)$ is the proper, convex and lower semicontinuous function
- 669 defined by

(B.1) (for all
$$z \in \mathbb{R}^n$$
) $| f^{\star}(z) = \sup_{x \in \operatorname{dom}(f)} \{ \langle x, z \rangle - f(x) \}.$

The subdifferential $\partial f(x)$ of $f(\cdot)$ at $x \in \mathbb{R}^n$ is defined by

$$\partial f(x) = \{ \omega \in \mathbb{R}^n | f(v) - f(x) \ge \langle \omega, v - x \rangle, \forall v \in \mathbb{R}^n \}.$$

- 670 We denote by $\text{Dom}(\partial f) \stackrel{\Delta}{=} \{x \in \mathbb{R}^n | \partial f(x) \neq \emptyset\}$ the domain of the subdifferential operator $\partial f : \mathbb{R}^n \to \mathbb{R}^n$
- ⁶⁷¹ \mathbb{R}^n . Recall that (see, *e.g.*, Theorem 2, Chapter 10, Section 3 in [11]): Dom(∂f) ⊂ dom(f).

Let x_0 be any element in the domain dom(f) of $f(\cdot)$, the recession function $f_{\infty}(\cdot)$ of $f(\cdot)$ is defined by

(for all
$$x \in \mathbb{R}^n$$
): $f_{\infty}(x) = \lim_{\lambda \to +\infty} \frac{1}{\lambda} f(x_0 + \lambda x)$.

⁶⁷² The function $f_{\infty} : \mathbb{R}^n \to \mathbb{R} \cup \{+\infty\}$ is a proper convex and lower semicontinuous function which ⁶⁷³ describes the asymptotic behavior of $f(\cdot)$.

Let $K \subset \mathbb{R}^n$ be a nonempty closed convex set. Let x_0 be any element in K. The recession cone of K is defined by [72] [89, Definition 1.11]:

$$K_{\infty} = \bigcap_{\lambda > 0} \frac{1}{\lambda} (K - x_0) = \{ u \in \mathbb{R}^n | x + \lambda u \in K, \forall \lambda \ge 0, \forall x \in K \}.$$

⁶⁷⁴ The set K_{∞} is a nonempty closed convex cone that is described in terms of the directions which

- recede from *K*. The indicator function of a set $K \subseteq \mathbb{R}^n$ is $\Psi_K(x) = 0$ if $x \in K$, $\Psi_K(x) = +\infty$ if $x \notin K$. ⁶⁷⁶ If *K* is closed non empty convex, we have $\partial \Psi_K(x) = \mathcal{N}_K(x)$, the so-called normal cone to *K* at *x*,
- defined as $N_K(x) = \{v \in \mathbb{R}^n \mid v^{\top}(s-x) \le 0 \text{ for all } s \in K\}$. When K is finitely represented, *i.e.*,
- 678 $K = \{x \in \mathbb{R}^n \mid k_i(x) \ge 0, 1 \le i \le m\}$, and if the functions $k_i(\cdot)$ satisfy some constraint qualification
- (like, independency, or extensions like the MFCQ), then $N_K(x)$ is generated by the outwards normals
- at the active constraints $k_i(x) = 0$, *i.e.*, $\mathcal{N}_K(x) = \{v \in \mathbb{R}^n \mid v = -\lambda_i \nabla k_i(x), k_i(x) = 0, \lambda_i \ge 0\}$.
- Let us here recall some important properties of the recession function and recession cone (see *e.g.*, [17,
- 682 Proposition 1.4.8]):
- ⁶⁸³ Proposition B.1. The following statements hold:
- a) Let $f_1 : \mathbb{R}^n \to \mathbb{R} \cup \{+\infty\}$ and $f_2 : \mathbb{R}^n \to \mathbb{R} \cup \{+\infty\}$ be two proper, convex and lower semicontinuous
- functions. Suppose that $f_1 + f_2$ is proper. Then for all $x \in \mathbb{R}^n$: $(f_1 + f_2)_{\infty}(x) = (f_1)_{\infty}(x) + (f_2)_{\infty}(x)$.

- b) Let $f : \mathbb{R}^n \to \mathbb{R} \cup \{+\infty\}$ be a proper, convex and lower semicontinuous function and let K be a nonempty
- closed convex set, such that $f + \Psi_K$ is proper (equivalently dom $(f) \cap K$ is non empty). Then for all $x \in \mathbb{R}^n$:

688 $(f + \Psi_K)_{\infty}(x) = f_{\infty}(x) + (\Psi_K)_{\infty}(x).$

- c) Let $K \subset \mathbb{R}^n$ be a nonempty, closed and convex set. Then for all $x \in \mathbb{R}^n$: $(\Psi_K)_{\infty}(x) = \Psi_{K_{\infty}}(x)$. Moreover for all $x \in K$ and $e \in K_{\infty}$: $x + e \in K$.
- 691 d) If $K \subseteq \mathbb{R}^n$ is a nonempty closed and convex cone, then $K_{\infty} = K$.

⁶⁹² e) Let $K = P(A, b) \stackrel{\Delta}{=} \{x \in \mathbb{R}^n | Ax \ge b\}$ for $A \in \mathbb{R}^{m \times n}$ and $b \in \mathbb{R}^m$. If $K \ne \emptyset$ then $K_{\infty} = P(A, 0) = \{x \in \mathbb{R}^n | Ax \ge 0\}$.

694 $f \in \mathbb{R}^n$ is a non-empty closed convex bounded set if and only if $K_{\infty} = \{0_n\}$.

⁶⁹⁵ g) Let $f : \mathbb{R}^n \to \mathbb{R} \cup \{+\infty\}$ be a proper, convex and lower semicontinuous function. Then $epi(f_{\infty}) =$ ⁶⁹⁶ $(epi(f))_{\infty}$.

- ⁶⁹⁷ Sets as in item e) are called \mathcal{H} -polyhedra, and there is an equivalence between sets P(A, 0) and finitely
- ⁶⁹⁸ generated convex cones [89, Theorem 1.3]. Let us now concatenate [5, Theorem 3, Corollaries 3 and 4].
- ⁶⁹⁹ They concern variational inequalities (VIs) of the form: Find $u \in \mathbb{R}^n$ such that

(B.2)
$$\langle \mathbf{M}u + \mathbf{q}, v - u \rangle + \varphi(v) - \varphi(u) \ge 0$$
, for all $v \in \mathbb{R}^{n}$

- where $\mathbf{M} \in \mathbb{R}^{n \times n}$ is a real matrix, $\mathbf{q} \in \mathbb{R}^n$ a vector and $\varphi : \mathbb{R}^n \to \mathbb{R} \cup \{+\infty\}$ a proper convex and
- lower semicontinuous function. The VI in (B.2) is equivalent to the inclusion $\mathbf{M}u + \mathbf{q} \in -\partial \varphi(u) \Leftrightarrow u \in$

702
$$(\mathbf{M} + \partial \varphi)^{-1}(-\mathbf{q}).$$

The problem in (B.2) is denoted as $VI(\mathbf{M}, \mathbf{q}, \varphi)$ in the next proposition. We also set:

(B.3)
$$\mathcal{K}(\mathbf{M},\varphi) = \{ x \in \mathbb{R}^n | \mathbf{M}x \in (\mathrm{dom}(\varphi_\infty))^{\star} \}.$$

- Note that $(\operatorname{dom}(\varphi_{\infty}))^{\star}$ is the dual cone of the domain of the recession function φ_{∞} while $(\operatorname{dom}(\varphi))_{\infty}$ (that we may denote also as $\operatorname{dom}(\varphi)_{\infty}$) is the recession cone of $\operatorname{dom}(\varphi)$.
- Proposition B.2. [5] Let $\varphi : \mathbb{R}^n \to \mathbb{R} \cup \{+\infty\}$ be a proper, convex and lower semicontinuous function with closed domain, $\mathbf{M} \in \mathbb{R}^{n \times n}$, and suppose that $\mathbf{M} \ge 0$ (not necessarily symmetric).
- *a)* If $(\operatorname{dom}(\varphi))_{\infty} \cap \operatorname{ker}\{\mathbf{M} + \mathbf{M}^{\top}\} \cap \mathcal{K}(\mathbf{M}, \varphi) = \{0\}$ then for each $\mathbf{q} \in \mathbb{R}^n$, problem $VI(\mathbf{M}, \mathbf{q}, \varphi)$ has at *least one solution.*
- ⁷¹⁰ b) Suppose that $(\operatorname{dom}(\varphi))_{\infty} \cap \operatorname{ker}\{\mathbf{M} + \mathbf{M}^{\top}\} \cap \mathcal{K}(\mathbf{M}, \varphi) \neq \{0\}$. If there exists $x_0 \in \operatorname{dom}(\varphi)$ such that

(B.4)
$$\langle \mathbf{q} - \mathbf{M}^{\top} x_0, v \rangle + \varphi_{\infty}(v) > 0, \ \forall v \in (dom(\varphi))_{\infty} \cap \ker{\{\mathbf{M} + \mathbf{M}^{\top}\}} \cap \mathcal{K}(\mathbf{M}, \varphi), \ v \neq 0,$$

- *Then problem* $VI(\mathbf{M}, \mathbf{q}, \varphi)$ *has at least one solution.*
- ⁷¹² b') If $\mathbf{M} = \mathbf{M}^{\top}$ then one can take $x_0 = 0$ in b).
- *c)* If u_1 and u_2 denote two solutions of problem $VI(\mathbf{M}, \mathbf{q}, \varphi)$ then $u_1 u_2 \in \text{ker}\{\mathbf{M} + \mathbf{M}^{\mathsf{T}}\}$.
- *d)* If $\mathbf{M} = \mathbf{M}^{\top}$ and u_1 and u_2 denote two solutions of problem $VI(\mathbf{M}, \mathbf{q}, \varphi)$, then $\langle \mathbf{q}, u_1 u_2 \rangle = \varphi(u_2) \varphi(u_1)$.
- ⁷¹⁶ *e)* If $\mathbf{M} = \mathbf{M}^{\top}$ and $\varphi(x + z) = \varphi(x)$ for all $x \in \text{dom}(\varphi)$ and $z \in \text{ker}\{\mathbf{M}\}$ and $\langle \mathbf{q}, e \rangle \neq 0$ for all ⁷¹⁷ $e \in \text{ker}\{\mathbf{M}\}, e \neq 0$, then problem $VI(\mathbf{M}, \mathbf{q}, \varphi)$ has at most one solution.
- *f)* If $\mathbf{M} = \mathbf{M}^{\top}$, then *u* is a solution of $VI(\mathbf{M}, \mathbf{q}, \varphi)$ if and only if it is a solution of the optimization problem min_{$x \in \mathbb{R}^n = \frac{1}{2}x^{\top}\mathbf{M}x + \langle \mathbf{q}, x \rangle + \varphi(x)$.}

- Notice that the function $\varphi(\cdot)$ will never be strictly convex in our case (it is an indicator function) so
- that the strict convexity argument of [5, Theorem 5] which applies when **M** is a P_0 -matrix never holds.
- ⁷²² The study of VIs as in (B.2) can be traced back to [76].

723 APPENDIX C SOME CONVEX ANALYSIS AND COMPLEMENTARITY THEORY TOOLS

- If $K \subset \mathbb{R}^n$ is a set, then $K^* = \{z \in \mathbb{R}^n | \langle z, x \rangle \ge 0 \text{ for all } x \in K\}$ is its dual cone. Its closure is denoted
- \bar{K} . Let *K* be a nonempty closed convex cone, then:

(c.1)
$$K^{\star} \ni x \perp y \in K \iff x \in -\mathcal{N}_{K}(y) \iff y \in -\mathcal{N}_{K^{\star}}(x).$$

⁷²⁶ Let $M = M^{\top} > 0$, *x* and *y* two vectors, then

(c.2)
$$M(x-y) \in -\mathcal{N}_K(x) \Leftrightarrow x = \operatorname{proj}_M[K; y] \Leftrightarrow x = \min_{z \in K} \frac{1}{2} (z-y)^\top M(z-y).$$

We note that this is a particular case of (B.2), so that Proposition B.2 can be considered as the charac-

terization of a generalized projection operator $VI(\mathbf{M}, \mathbf{q}, \varphi)$. The first equivalence is [15, Proposition 6.46].

729 **6.46**].

730 APPENDIX D PASSIVE DESCRIPTOR VARIABLE SYSTEMS

Dissipative systems have been formalized rigorously by J.C. Willems in his two seminal articles [87, 88].

Passive systems are a particular case of dissipative systems. A quadruple (A, B, C, D) is said passive if

⁷³³ the next LMI is satisfied:

(D.1)
$$\begin{pmatrix} -XA - A^{\top}X & -XB + C^{\top} \\ -B^{\top}X + C & D + D^{\top} \end{pmatrix} \ge 0, \quad X = X^{\top} \ge 0.$$

It is said *strictly* passive if the inequalities are satisfied with > 0. This is not to be confused with the strict *state* passivity [25, Definition 4.54], since strict passivity implies strict state passivity, but not the inverse in general. In fact strict passivity is directy related to strongly SPR transfer matricess [25, p.62 and section 3.12.2], while SPR transfer matrices with minimal realisations are strictly state passive [25, Theorem 4.73] [66], *i.e.*, $-XA - A^TX + \mu X \ge 0$ for some $\mu > 0$, and one may have $D + D^T \ge 0$ but not > 0.

Positive real transfer matrices [25, Definition 2.34] with a state space realization represented by quintuples (*P*, *A*, *B*, *C*, *D*) which are minimal (see Appendix G for conditions of minimality of a statespace realization for descriptor systems) can be rewritten as in (2.2) [28, 51], and they have the positive semidefinite storage function $V(x_1, x_2) = \frac{1}{2}(x_1^{\top}X_1x_1 + x_2^{\top}X_3x_2)$, with supply rate $\lambda^{\top}w$. Indeed, by differentiating $V(x_1(t), x_2(t))$ along the trajectories of (2.2), adding $\pm \lambda^{\top}w$ and completing the squares, one can compute that:

(D.2)
$$\dot{V}(t) = \frac{1}{2} (x_1^{\top}, \lambda^{\top}) \underbrace{\begin{pmatrix} X_1 A_1 + A_1^{\top} X_1 & X_1 B_1 - C_1^{\top} \\ B_1^{\top} X_1 - C_1 & -(\hat{D} + \hat{D}^{\top}) \end{pmatrix}}_{\stackrel{\underline{A}=Q_1}{\underline{A}}} \begin{pmatrix} x_1 \\ \lambda \end{pmatrix} + \lambda^{\top}(t) w(t),$$

where $X_1 = X_1^{\top} \ge 0$, $Q_1 \le 0$, and if the system is strongly SPR then $X_1 \ge 0$, $Q_1 < 0$ [28, Theorem 5.4]. The equality in (D.2) is an infinitesimal dissipation equality, showing that the system is passive [25, Definition 2.1, Theorem 4.53]. Such systems possess particular structural properties [30, 33] [25,

- section 3.8]. In the SSPR case it follows from (D.2) and the complementarity conditions that $x_1(\cdot)$ and
- $_{750}$ $x_2^{\top}X_3x_2(\cdot)$ are bounded for any bounded initial condition.
- Let us recall some definitions, where A^* denotes the conjugate transpose of A (see [25, Chapter 2] for details):
- Definition D.1. A transfer matrix $H(s) \in \mathbb{C}^{m \times m}$, $s \in \mathbb{C}$, is positive real (PR) if: 1) H(s) has no poles in
- 754 $\mathbf{Re}(H(s)) > 0, 2$ H(s) is real for positive real s, 3 $H(s) + H^{\star}(s) \ge 0$ for all $\mathbf{Re}(s) > 0$. Let H(s) be
- regular, then it is strictly positive real (SPR) if $H(s \epsilon)$ is PR for some $\epsilon > 0$. It is strongly SPR (SSPR)
- if 1) H(s) is analytic in $\operatorname{Re}(s) \ge 0$, 2) $\operatorname{Re}(H(j\omega)) \ge \delta > 0$ for all $\omega \in [-\infty, +\infty]$ and some real δ .

757 APPENDIX E GRONWALL-BACHPATTE INEQUALITY [13, 47]

Let us consider a positive real numbers sequence $\{x_k\}$ such that $x_k \le \alpha_k + \sum_{i=0}^{k-1} \beta_i x_i$, where $\{\alpha_k\}$ and $\{\beta_k\}$ are two real numbers sequences with $\{\beta_k\}$ positive. Then one has

$$x_k \leq \alpha_k + \sum_{i=0}^{k-1} \beta_i \alpha_i \ \Pi_{j=i+1}^k (1+\beta_i).$$

758 APPENDIX F CALCULATION OF eta_i

- The starting point is the term $\eta ||C_1|| ||B_1|| \sum_{i=1}^k h_1^{k-i} \sum_{j=0}^{i-1} (\eta h_3)^{i-1-j} ||x_{1,j+1}||$ in (5.19). Let us consider the factors of the terms $||x_{1,j+1}||$ separately:
- Factor of $||x_{1,1}||$ (all terms with j = 0): $h_1^{k-1} + h_1^{k-2}(\eta h_3) + h_1^{k-3}(\eta h_3)^2 + \ldots + h_1(\eta h_3)^{k-2} + (\eta h_3)^{k-1} = \sum_{i=1}^k h_1^{k-i}(\eta h_3)^{i-1}$.

• Factor of $||x_{1,2}||$ (all terms with j = 1): $h_1^{k-2} + h_1^{k-3}(\eta h_3) + h_1^{k-4}(\eta h_3)^2 + \ldots + h_1(\eta h_3)^{k-3} + (\eta h_3)^{k-2} = \sum_{i=2}^k h_1^{k-i}(\eta h_3)^{i-2}$.

- 765 ...
- Factor of $||x_{1,k-1}||$ (all terms with j = k 2): $h_1 + (\eta h_3) = \sum_{i=k-1}^k h_1^{k-i} (\eta h_3)^{i-k+1}$.
- Factor of $||x_{1,k}||$ (all terms with j = k): $1 = \sum_{i=k}^{k} h_1^{k-i} (\eta h_3)^{i-k}$.
- 768 Therefore one infers (5.22).

769 APPENDIX G MINIMALITY OF STATE SPACE REALIZATIONS

⁷⁷⁰ Let us report the results stated in [28, Theorem 3.1] [44, Theorems 2-2.1, 2-3.1, 2-6.3].

Theorem G.1. Let $H(s) = C(sP - A)^{-1}B + D$, $s \in \mathbb{C}$, be a rational transfer function where P and A are n × n matrices. Then $P\dot{x}(t) = Ax(t) + Bu(t)$, y(t) = Cx(t) + Du(t) is a minimal realization of H(s) if, and only if, the following conditions are satisfied:

- $rank(A sP \ B) = n \text{ for all } s \in \mathbb{C},$
- $rank(P \ B) = n$,
- $rank(A^{\top} sP^{\top} C^{\top}) = n$,
- $rank(P^{\top} C^{\top}) = n$,
- $A \operatorname{Ker}(P) \subseteq \operatorname{Im}(P)$.

779 APPENDIX H DIFFERENCES BETWEEN (3.1) AND IN (5.5)

Let us consider the discrete scheme (5.5) with E(t) = 0 and X_1 full rank, this is rewritten as:

(H.1)
$$\begin{cases} (a) \ x_{1,k+1} = x_{1,k} + hA_1x_{1,k} + hB_1\lambda_k \\ (b) \ B_3^{\top}X_3B_3\lambda_{k+1} = B_3^{\top}X_3B_3\lambda_k - hC_1x_{1,k} - h\hat{D}\lambda_k + w_{k+1} \\ (c) \ 0 \le \lambda_{k+1} \perp w_{k+1} = C_1x_{1,k+1} + C_2x_{2,k+1} + C_3z_{k+1} + D\lambda_{k+1} \ge 0 \\ (d) \ x_{2,k} = -B_3\lambda_k \\ (e) \ z_k = -B_2\lambda_k - B_3\frac{\lambda_{k+1}-\lambda_k}{k}. \end{cases}$$

This is the numerical scheme with which the examples have been simulated. Combining (H.1) (d) and (e), gives:

(H.2)
$$x_{2,k+1} = x_{2,k} + hz_k + hB_2\lambda_k.$$

⁷⁸³ Notice that this is the choice made in (6.17). The equations in (H.1) (a) and (H.2) are not equal to their

counterparts in (3.1), because λ_k and z_k are used in (H.1) (a) and (H.2) instead of λ_{k+1} and z_{k+1} in (3.1).

This explains that the numerical data obtained from both schemes, differ. This also explains why the convergence results obtained for both schemes, are not the same neither.

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