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ANALYSIS OF THE IMPLICIT EULER TIME-DISCRETIZATION OF SEMI-EXPLICIT DIFFERENTIAL-ALGEBRAIC LINEAR COMPLEMENTARITY SYSTEMS*

BERNARD BROGLIATO †

Abstract. This article is largely concerned with the time-discretization of differential-algebraic equations (DAE) with complementarity constraints, which we name differential algebraic linear complementarity systems (DALCS). Specifically, the Euler implicit discretization of DALCS is analysed: the one-step non-smooth problem (OSNSP), that is a generalized equation, is shown to be wellposed under some conditions, then the convergence of the discretized solutions is studied, and the existence of solutions to the continuous-time system is shown as a consequence. Passivity of some operators is pivotal to the analysis. Examples from circuits, mechanics and switching DAE illustrate the applicability of the developments.

13 **Key words.** differential-algebraic system, descriptor-variable system, linear complementarity 14 system, recession function, discrete-time, well-posedness, passive system, Lur'e equations.

15 **AMS subject classifications.** 65K15, 90C33, 90C90, 93B52, 93D30, 93C30.

1. Introduction. The analysis of non-smooth dynamical systems with multi-16 valued right-hand sides satisfying maximal monotone properties has witnessed a large 17number of contributions [19]. Within this class one finds linear and nonlinear com-18 19plementarity dynamical systems, projected systems, differential variational inequali-20 ties, differential inclusions with maximal monotone right-hand sides, Moreau's sweeping process, some switching dynamics, etc. In this article we focus on singular, or 21 differential-algebraic linear complementarity systems (DALCS), which may be viewed 22 either as an extension of classical differential-algebraic equations (DAEs), or of "clas-2324 sical" linear complementarity systems. As will become clear later, DALCS represent 25an extension of the set-valued Lur'e systems studied, e.q., in [39, 40]. The associated canonical form is motivated by applications in chemistry [53, 52, 62], switching DAEs 26 analysis [48], circuits with nonsmooth electronic components [2, section 3.5, chapter 4] 27 [22], optimization-constrained differential equations [19, Example 2] with applications 28 29in atmospheric aerosol particles [38], and Lagrangian systems with both bilateral and unilateral constraints [17] (in spite of the fact that the results presented in this article 30 do not apply to all of these systems). 31 The study of singular nonsmooth dynamical systems has not received a lot of attention yet. Singularly perturbed differential inclusions have been tackled in [26, 56, 32, 58]. 33

The conditions on the set-valued mappings that are imposed (compactness [56, 58], boundedness and Hausdorff Lipschitz continuity [32], one-sided Lipschitz condition [26]) are not satisfied by complementarity conditions (or by normal cones to convex sets). Hence these results do not apply to DALCS. The well-posedness of DAEs in semi-explicit form with nonsmooth constraints (which are related to the class of systems studied in this article, see Remark 2.1 below) has been analysed in [53, 54, 52]. Passive singular set-valued systems are studied in [22], who derived well-posedness

41 relying on a special Weierstrass' representation. However the time-discretisation is

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- 42 not tackled in the above references.
- 43 Our goal is to analyse the well-posedness of the one-step nonsmooth problem (OS-
- 44 NSP) obtained after an implict (backward) Euler discretization, and the convergence
- 45 of the discrete-time solutions towards a continuous-time limit. The second step is use-
- 46 ful not only for proving the existence of solutions to the continuous-time system, but
- 47 also, perhaps most importantly, to show that the proposed time-discretization makes
- 48 sense for numerical simulations. Passivity properties of suitable operators are pivotal
- 49 throughout the article. It is noteworthy that implicit Euler methods are well-known
- 50 in DAEs [30, 14, 33, 42, 12, 35, 61, 7], differential variational inequalities [46, 24],
- and complementarity systems fields [41, 21, 59, 34], as well as in the field of maximal monotone differential inclusions [8, 47] (implicit Euler method is often called
- proximal algorithm, while Euler method refers to the explicit algorithm [47]), and Moreau's sweeping processes with the catching-up algorithm [19, 44, 45]. However
- 54 Moreau's sweeping processes with the catching-up algor 55 their application to DALCS has not been studied yet.
- 56 The article is organised as follows: Section 2 introduces the general dynamics of
- 57 the nonsmooth singular systems, section 3 presents their time-discretization, section
- ⁵⁸ 4 studies the well-posedness of the OSNSP, section 5 is dedicated to convergence
- analyses: results in section 5.1 are based on the properties of the numerical scheme,
- 60 while section 5.2 uses properties of the system. Section 6 is dedicated to illustrating
- 61 examples and simulations, from circuits, mechanics and switching DAEs. Conclusions
- ⁶² are drawn in section 7. The appendix is dedicated to recall various mathematical tools.
- Notation and definitions: for any vector $x \in \mathbb{R}^n$ and any matrix $M \in \mathbb{R}^{m \times n}$, 63 ||x|| is the Euclidean norm and ||M|| is the Frobenius norm, which are compatible 64 norms [10, Proposition 9.3.5], *i.e.*, $||Mx|| \leq ||M|| ||x||$. Let $M \in \mathbb{R}^{n \times m}$, then Im(M)65 is its range, Ker(M) is its null space. We use $\langle x, y \rangle = x^{\top}y$, so $\langle x, x \rangle = ||x||^2$. Positive 66 definite matrix: $M \succ 0$ if $x^{\top}Mx > 0$ for all $x \neq 0$, positive semidefinite matrix: 67 $M \succeq 0$ if $x^{\top} M x \ge 0$ for all x (such M is not necessarily symmetric). The maximum 68 singular value is denoted as $\sigma_{\max}(M)$, and the minimum eigenvalue as $\lambda_{\min}(M)$. The 69 ith row of M is denoted as $M_{i\bullet}$. The $n \times n$ identity matrix is denoted I_n . A set-valued 70 mapping $A: \mathbb{R}^n \rightrightarrows \mathbb{R}^n$ is said monotone if for all $x_1, x_2, y_1 \in A(x_1), y_2 \in A(x_2)$, one 71 has $\langle x_1 - x_2, y_1 - y_2 \rangle \geq 0$. It is maximal monotone if its graph cannot be enlarged 72without destroying monotonicity. It is η -strongly monotone if there exists $\eta > 0$ such 73 that $\langle x_1 - x_2, y_1 - y_2 \rangle \ge \eta ||x_1 - x_2||^2$. See also Appendix A for further definitions. 74
- **2.** The class of nonsmooth singular dynamical systems. Let us consider the following differential algebraic linear complementarity system (DALCS):

77 (2.1)
$$\begin{cases} (a) \ P_I \dot{x}(t) = A_I x(t) + B_I z(t) + \bar{B}_I \lambda(t) + E_I(t) \\ (b) \ 0 = A_{II} x(t) + B_{II} z(t) + \bar{B}_{II} \lambda(t) + E_{II}(t) \\ (c) \ K^* \ni \lambda(t) \perp w(t) = C_I x(t) + C_{II} z(t) + D\lambda(t) + F(t) \in K, \end{cases}$$

with $x(t) \in \mathbb{R}^{n_1}$, $z(t) \in \mathbb{R}^{n_2}$, $\lambda(t) \in \mathbb{R}^m$, $w(t) \in \mathbb{R}^m$, $P_I \in \mathbb{R}^{n_1 \times n_1}$ is full-rank, $K \subseteq \mathbb{R}^m$ is a closed nonempty convex cone, K^* is its dual cone, $\bar{B}_{II} \in \mathbb{R}^{p \times m}$, $A_{II} \in \mathbb{R}^{p \times n_1}$, $B_{II} \in \mathbb{R}^{p \times n_2}$, $C_I \in \mathbb{R}^{m \times n_1}$, $C_{II} \in \mathbb{R}^{m \times n_2}$, $D \in \mathbb{R}^{m \times m}$, $B_I \in \mathbb{R}^{n_1 \times n_2}$, $A_I \in \mathbb{R}^{n_1 \times n_1}$, $\bar{B}_I \in \mathbb{R}^{n_1 \times m}$, and $E_I(\cdot)$, $E_{II}(\cdot)$, $F(\cdot)$ are bounded, Lipschitz continuous functions of time. The form in (2.1) is that of a semi-explicit DALCS. It corresponds to a kind of extension of the impulse-free Weierstrass form [20], where the transformed state matrix has non zero off-diagonal terms A_{II} and B_I .

Remark 2.1. Assume that $K = \mathbb{R}^m_+$. The complementarity constraint (2.1) (c) is equivalently rewritten as $\phi(\lambda, w) = 0$ where $\phi(\cdot)$ is a *C*-function [3, 29]. Therefore

(2.1) (b) (c) can be rewritten equivalently as $\Phi(x, z, \lambda, t) = 0$ for some function $\Phi(\cdot)$, 87 usually nondifferentiable. The system (2.1) thus possesses the form of a semi-explicit 88 DAE as studied in [53, 52], assuming piecewise differentiable functions $\Phi(\cdot)$, with x 89 the differential state, z and λ the algebraic state variables. The so-called regularity 90 conditions (which are an extension of index 1 conditions for DAEs) used in [53, 54, 91 52, 51 may hold in particular cases for (2.1). Examples in section 6.3 illustrate the 92

links betwen both approaches in [53, 54, 52, 51] and in this article. 93

Consider (2.1), then the initial data $x(0) = x_0$ and $z(0) = z_0$ are compatible with the 94system's constraints if and only if: 95

96 (2.2)
$$0 \in A_{II}x_0 + B_{II}z_0 + E_{II}(0) - \bar{B}_{II}(D + \mathcal{N}_{K^*})^{-1}(C_Ix_0 + C_{II}z_0 + F(0)) \\ \lambda(0) \in -(D + \mathcal{N}_{K^*})^{-1}(C_Ix_0 + C_{II}z_0 + F(0)).$$

This implies that $C_I x_0 + C_{II} z_0 + F(0) \in \text{Im}(D + \mathcal{N}_{K^*})$ and $A_{II} x_0 + B_{II} z_0 + E_{II}(0) \in$ B_{II} Dom $(D + \mathcal{N}_{K^*})$. 98

3. Time-discretization of DALCS with implicit Euler methods. Let us 99 consider the time interval $[0,T], T > 0, h = \frac{T}{N}, N \in \mathbb{N}_+, t_k = kh, k \in \{0,1,\ldots,N\},$ 100 $t_0 = 0, t_N = T$. For a continuous function $f(\cdot), f_k \stackrel{\Delta}{=} f(t_k)$. The implicit Euler 101

102 discretization of
$$(2.1)$$
 reads as:

103 (3.1)
$$\begin{cases} (a) \ P_I(x_{k+1} - x_k) = hA_I x_{k+1} + hB_I z_{k+1} + hB_I \lambda_{k+1} + hE_{I,k} \\ (b) \ A_{II} x_{k+1} + B_{II} z_{k+1} + \bar{B}_{II} \lambda_{k+1} + E_{II,k} = 0 \\ (c) \ K^* \ni \lambda_{k+1} \perp C_I x_{k+1} + C_{II} z_{k+1} + D\lambda_{k+1} + F_k \in K. \end{cases}$$

The problem in (3.1) (b) (c) is a mixed LCP (MLCP). Let us now state a basic 104 105assumption.

106 ASSUMPTION 1. The matrix
$$P_I^h \stackrel{\Delta}{=} P_I - hA_I$$
 is full-rank.

In case $P_I \succ 0$, [25, Theorem 2.11] can be used to set conditions on hA_I so that 107 $P_I^h \succ 0$: there always exists h_{\max} such that the assumption holds for all $h \in [0, h_{\max})$ 108 (with $h_{\max} = +\infty$ if $A_I \preccurlyeq 0$). Let us define the following matrices and vectors: 1. $\tilde{A} = h A_{II} (P_I^h)^{-1} B_I + B_{II},$ 109

110

111 2.
$$B = hA_{II}(P_I^h)^{-1}\bar{B}_I + \bar{B}_{II},$$

112 3.
$$E_k = E_{II,k} + hA_{II}(P_I^h)^{-1}E_{I,k} + A_{II}(P_I^h)^{-1}P_Ix_k,$$

- 4. $\hat{C} = hC_I(P_I^h)^{-1}B_I + C_{II},$ 5. $\tilde{D} = hC_I(P_I^h)^{-1}\bar{B}_I + D,$ 113
- 114

115 6.
$$F_k = F_k + hC_I(P_I^h)^{-1}E_{I,k} + C_I(P_I^h)^{-1}P_Ix_k.$$

116Then we deduce from (3.1) the difference equations:

117 (3.2)
$$\begin{cases} (a) \ x_{k+1} = (P_I^h)^{-1}(P_I x_k + hB_I z_{k+1} + h\bar{B}_I \lambda_{k+1} + hE_{I,k}) \\ (b) \ \tilde{A} z_{k+1} + \tilde{B} \lambda_{k+1} + \tilde{E}_k = 0 \\ (c) \ K^* \ni \lambda_{k+1} \perp \tilde{C} z_{k+1} + \tilde{D} \lambda_{k+1} + \tilde{F}_k \in K, \end{cases}$$

where the last two lines of (3.2) are a MLCP with unknowns z_{k+1} and λ_{k+1} . It is 118therefore possible to study the well-posedness of this MLCP independently of the 119

first line in (3.2). Letting $\zeta_{k+1} \stackrel{\Delta}{=} -z_{k+1}$, this MLCP can be rewritten equivalently as 120 follows: 121

122 (3.3)
$$\begin{cases} \lambda_{k+1} \in (\mathcal{N}_{K^{\star}} + \tilde{D})^{-1}(-\tilde{C}z_{k+1} - \tilde{F}_k), \\ -\tilde{A}\zeta_{k+1} + \tilde{E}_k \in -\tilde{B} \ (\mathcal{N}_{K^{\star}} + \tilde{D})^{-1}(\tilde{C}\zeta_{k+1} - \tilde{F}_k), \\ x_{k+1} = (P_I^h)^{-1}(P_I x_k + hB_I z_{k+1} + h\bar{B}_I \lambda_{k+1} + hE_{I,k}). \end{cases}$$

Also the following calculations stem from the complementarity conditions in (3.2): 123 (3.4)

$$\lambda_{k+1} \in -\mathcal{N}_K(Cz_{k+1} + D\lambda_{k+1} + F_k) \quad \Leftrightarrow -\lambda_{k+1} \in \mathcal{N}_{\Gamma_k}(Cz_{k+1} + D\lambda_{k+1})$$

$$\Leftrightarrow \quad \partial \sigma_{\Gamma_k}(-\lambda_{k+1}) \ni \tilde{C}z_{k+1} + \tilde{D}\lambda_{k+1} \qquad \Leftrightarrow (\partial \sigma_{\Gamma_k} + \tilde{D})(-\lambda_{k+1}) \ni \tilde{C}z_{k+1}$$

$$\Leftrightarrow \quad \lambda_{k+1} \in -(\partial \sigma_{\Gamma_k} + \tilde{D})^{-1}(\tilde{C}z_{k+1}).$$

where the inversion of set-valued mappings is used in the second equivalence, and 125where $\sigma_{\Gamma_k}(\cdot)$ is the support function of the set $\Gamma_k(t_k, x_k) \stackrel{\Delta}{=} K - \tilde{F}_k(t_k, x_k)$ (which is a closed convex set, but not a cone as long as $\tilde{F}_k \neq 0$ which is the generic case). 126127Therefore equivalently we can rewrite (3.3) as: 128

(3.5)
$$\begin{cases} (a) \ \lambda_{k+1} \in -(\partial \sigma_{\Gamma_k} + \tilde{D})^{-1}(\tilde{C}z_{k+1}) \\ (b) \ -\tilde{A}z_{k+1} + \tilde{B}(\partial \sigma_{\Gamma_k} + \tilde{D})^{-1}(\tilde{C}z_{k+1}) \ni \tilde{E}_k \\ (c) \ x_{k+1} = (P_I^h)^{-1}(P_I x_k + hB_I z_{k+1} + h\bar{B}_I \lambda_{k+1} + hE_{I,k}). \end{cases}$$

The second line in (3.3) is a generalized equation (GE) with unknown ζ_{k+1} . Its well-130 posedness depends on the properties of the operator $\mathcal{F} : \zeta \mapsto -A\zeta + B (\mathcal{N}_{K^{\star}} +$ 131 $\tilde{D})^{-1}(\tilde{C}\zeta - \tilde{F}_k)$. Similarly for (3.5) with the operator $\mathcal{G}: \zeta \mapsto -\tilde{A}\zeta + \tilde{B}(\partial\sigma_{\Gamma_k} + \tilde{D})$ 132 $(\hat{D})^{-1}(\hat{C}\zeta)$, and (3.5) (b) is: $\mathcal{G}(z_{k+1}) \ni \hat{E}_k$. Various results have been published to 133 characterize such operators [16, 23, 6] and to guarantee that they are single-valued 134Lipschitz continuous [16, Propositions 1, 2, 3, Corollary 1], or maximal monotone 135(possibly set-valued) [23, 6]. The difficulty in the analysis of DALCS is that they 136involve the composition of two operators as seen in the OSNSP for advancing the 137 scheme from step k to step k + 1: 138(3.6)

139
$$P_I^h x_{k+1} \in P_I x_k + h B_I \mathcal{G}^{-1}(\tilde{E}_k(x_k)) - h \bar{B}_I (\partial \sigma_{\Gamma_k} + \tilde{D})^{-1} \left(\tilde{C} \mathcal{G}^{-1}(\tilde{E}_k(x_k)) \right) + h E_{I,k}$$

Clearly the scheme can be advanced with a unique solution if both λ_{k+1} and z_{k+1} 140are uniquely defined (the "if and only if" holds with $\bar{B}_I \lambda_{k+1}$ and $B_I z_{k+1}$ uniquely 141defined). Thus next we are going to characterize the GEs in (3.3) and (3.5) in order 142 143to analyse the scheme (3.6).

Two main classes of results are provided in the sequel of the article: first the well-144 posedness (existence and uniqueness of solutions) for the OSNSP associated with the 145time-discretizations in (3.1) (Propositions 4.3 and 4.4, Corollary 4.6) and (3.2) (Propo-146147sitions 4.7 and 4.8); second, the convergence of the discrete-time solutions is analysed, and the limit functions are shown to be solutions of the DALCS (Propositions 5.2 and 1481495.9).

4. OSNSP well-posedness. The approach chosen in this article consists of 150considering the system as an ODE (2.1) (a) with a nonsmooth constraint (2.1) (b) 151(c): either (3.1) (b) (c) that is a mixed LCP with unknowns z_{k+1} , λ_{k+1} and parameter 152 x_{k+1} (section 4.2), or (3.2) (b) (c) (equivalently (3.5) (a) (b)) that is a mixed LCP 153with z_{k+1} , λ_{k+1} (section 4.1). As will become clear later, the results rely heavily on 154the positive (semi) definiteness of B_{II} , and on specific passivity couplings between 155z in the complementarity constraint (matrix C_{II}) and λ in the algebraic constraint 156(matrix \bar{B}_{II}). When $\tilde{D} = 0$ in (3.2) (c), let us state the following: 157

ASSUMPTION 2. There exists a $p \times p$ matrix $\tilde{X} = \tilde{X}^{\top} \succ 0$ such that $\tilde{X}\tilde{B} = \tilde{C}^{\top}$. 158

- The alternative option when D = 0 in (3.1) (c): 159
- ASSUMPTION 3. There exists a $p \times p$ matrix $X = X^{\top} \succ 0$ such that $X\bar{B}_{II} = C_{II}^{\top}$. 160 161

Let us notice that Assumption 2 implies that $\tilde{X}\bar{B}_{II} = C_{II}^{\top} + \mathcal{O}(h)$, while Assumption 162 3 implies that $X\tilde{B} = \tilde{C}^{\top} + \mathcal{O}(h)$. The influence of the $\mathcal{O}(h)$ terms may not be easy 163to characterize precisely for finite h > 0. Assumption 3 is a model assumption, while 164 Assumption 2 is a property of the numerical method. One sees that Assumption 3 165corresponds to a passivity constraint on the triple (B_{II}, B_{II}, C_{II}) [18]. Assumption 166 2 imposes an input/output passivity-like constraint on the triple (A, B, C), which is 167 an $\mathcal{O}(h)$ approximation of $(B_{II}, \overline{B}_{II}, C_{II})$. When $D \neq 0$ and $D \neq 0$, extensions of 168 Assumptions 2 and 3 are as follows: 169

170 ASSUMPTION 4. The quadruple $(\tilde{A}, \tilde{B}, \tilde{C}, \tilde{D})$ is passive.

171 and the counterpart of Assumption 3:

172 ASSUMPTION 5. The quadruple $(B_{II}, B_{II}, C_{II}, D)$ is passive.

Passivity is defined in Appendix C. Notice that the above assumptions imply that $n_2 = p$. Assumption 4 refers to the numerical method while Assumption 5 is a model property. Assumptions 2 and 4 refer to (3.2) (b) (c) or (3.5) (a) (b), while Assumptions 3 and 5 refer to (3.1) (b) (c). Let us present a result about the relationships between Assumptions 4 and 5, *i.e.*, when does the plant's passivity guarantees the numerical scheme's passivity.

179 LEMMA 4.1. The following statements are true:

- 180 1. Let $(B_{II}, \overline{B}_{II}, C_{II}, D)$ be strictly passive. Then $(\tilde{A}, \tilde{B}, \tilde{C}, \tilde{D})$ is strictly passive 181 for small enough h > 0.
- 182 2. Let $(B_{II}, \bar{B}_{II}, C_{II}, D)$ be strictly state passive. Then $(\tilde{A}, \tilde{B}, \tilde{C}, \tilde{D})$ is strictly 183 state passive if: a) h > 0 is small enough, b) $D + D^{\top} = 0$ and $C_I(P_I^h)^{-1}\bar{B}_I +$ 184 $(C_I(P_I^h)^{-1}\bar{B}_I)^{\top} \succeq 0, c) XA_{II}(P_I^h)^{-1}\bar{B}_I - B_I^{\top}(P_I^h)^{-\top}C_I^{\top} = 0$, where X is a 185 solution of the passivity LMI for $(B_{II}, \bar{B}_{II}, C_{II}, D)$.

186 3. Let h > 0 and $(B_{II}, \overline{B}_{II}, C_{II}, D)$ be passive. Then (A, B, C, D) is passive if: 187 b) and c) in item 2 hold, and a) $XA_{II}(P_I^h)^{-1}B_I + B_I^{\top}(P_I^h)^{-\top}A_{II}^{\top}X \preccurlyeq 0$, 188 where X is a solution of the passivity LMI for $(B_{II}, \overline{B}_{II}, C_{II}, D)$.

189 **Proof:**

190 1. The strict passivity of $(B_{II}, \overline{B}_{II}, C_{II}, D)$ means that there exists $X = X^{\top} \succ 0$ 191 such that:

192 (4.1)
$$M \stackrel{\Delta}{=} \begin{pmatrix} -XB_{II} - B_{II}^{\top}X & -X\bar{B}_{II} + C_{II}^{\top} \\ -\bar{B}_{II}^{\top}X + C_{II} & D + D^{\top} \end{pmatrix} \succ 0.$$

193 In particular this implies that $D + D^{\top} \succ 0$ [18, Theorem A.65] and the 194 associated transfer function is strongly SPR [18, Remark 3.19, Section 3.12.2]. 195 Now we have $\tilde{A} = B_{II} + \mathcal{O}(h), \ \tilde{B} = \bar{B}_{II} + \mathcal{O}(h), \ \tilde{C} = C_{II} + \mathcal{O}(h), \ \tilde{D} =$ 196 $D + \mathcal{O}(h)$. Therefore

197 (4.2)
$$M = \underbrace{\begin{pmatrix} -X\tilde{A} - \tilde{A}^{\top}X & -X\tilde{B} + \tilde{C}^{\top} \\ -\tilde{B}^{\top}X + \tilde{C}^{\top} & \tilde{D} + \tilde{D}^{\top} \end{pmatrix}}_{\triangleq_{\tilde{M}}} + N(h)$$

198 where N(h) stands for a matrix with suitable dimensions and which vanishes 199 as $h \to 0$ (in fact N(h) is $\mathcal{O}(h)$). Thus $\tilde{M} = M - N(h)$. Applying [25, 200 Theorem 2.11], it follows that $\tilde{M} \succ 0$ provided that $\sigma_{\max}(N(h)) < \lambda_{\min}^2(M)$, 201 an inequality which can always be verified if h > 0 is small enough. Under 202 this condition, it follows that the strict passivity of $(B_{II}, \bar{B}_{II}, C_{II}, D)$ implies 203 the strict passivity of $(\tilde{A}, \tilde{B}, \tilde{C}, \tilde{D}) (\Rightarrow \tilde{D} \succ 0)$.

204		2.	The result is proved by developing M in (4.2), then using [10, Proposition
205			8.2.4, Fact 8.21.22], or [18, Lemma A.69], as well as the fact that $-XB_{II}$ –
206			$B_{II}^{\top}X \succ 0$ from the strict state passivity.
207		3.	The proof is similar to that of item 2.
	_		

²⁰⁸

The converse implication in item 1 does not hold, since one may have $M \succ 0$ for some 209 h > 0, but $M \not\geq 0$. It is possible to state necessary and sufficient conditions, relying 210 on [10, Proposition 8.2.4, Fact 8.21.22] or [18, Lemma A.70]. However the obtained 211 conditions are cumbersome and hardly usable. Lemma 4.1 can be used to determine 212213 the conditions under which Assumption 4 is verified in the results presented in the next section. The conditions b) and c) in Lemma 4.1 item 2 hold if: b) $C_I = 0$ or $B_I = 0$, c) 214215 $(A_{II} = 0 \text{ or } B_I = 0)$ and $(B_I = 0 \text{ or } C_I = 0)$. Condition a) in Lemma 4.1 item 3 holds if $A_{II} = 0$ or $B_I = 0$. Let us remind that the solution $X = X^{\top} \succ 0$ of the Lyapunov equation $XB_{II} + B_{II}^{\top}X = -Q$, $Q = Q^{\top} \succ 0$, can be written $X = (-\frac{1}{2}Q + F)B_{II}^{-1}$, 216217where $F = -F^{\top}$ is unique [60, Theorems 1 and 2]. Let us denote the symmetric 218part of a square matrix M as sym(M) and its skew-symmetric part as skw(M). This 219 220 allows us to state the following necessary conditions.

221 COROLLARY 4.2. 1) Let $(B_{II}, \overline{B}_{II}, C_{II}, D)$ be strictly state passive with $XB_{II} + B_{II}^{\top}X = -Q, Q = Q^{\top} \succ 0$, and assume that $A_{II}(P_I^h)^{-1}\overline{B}_I$ is right-invertible $(\Rightarrow p \le m)$. Then $X = (-\frac{1}{2}Q + F)B_{II}^{-1}$, for some $F = -F^{\top}$ uniquely defined, and condition 224 c) item 2 in Lemma 4.1 holds only if

225 (4.3)
$$\operatorname{sym}(\frac{1}{2}QB_{II}^{-1}A_{II}(P_I^h)^{-1}\bar{B}_I + B_I^{\top}(P_I^h)^{-\top}C_I^{\top}) = F\operatorname{skw}(B_{II}^{-1}A_{II}(P_I^h)^{-1}\bar{B}_I)F,$$

226 for any pair (Q, F), and

227 (4.4)
$$\operatorname{Ker}(C_I(P_I^h)^{-1}B_I) = \{0\}.$$

228 2) Let $(B_{II}, \overline{B}_{II}, C_{II}, D)$ be passive with $XB_{II} + B_{II}^{\top}X = -Q$, $Q = Q^{\top} \succeq 0$ and B_{II} 229 invertible. Then $X = (-\frac{1}{2}Q + F)B_{II}^{-1}$, for some $F = -F^{\top}$, and condition a) item 3 230 in Lemma 4.1 holds only if

231 (4.5)
$$\operatorname{sym}(-QB_{II}^{-1}A_{II}(P_I^h)^{-1}B_I) + 2F\operatorname{skw}(B_{II}^{-1}A_{II}(P_I^h)^{-1}B_I)F \preccurlyeq 0,$$

232 for any pair (Q, F).

Proof: 1) Condition c) yields $(-\frac{1}{2}Q + F)B_{II}^{-1}A_{II}(P_I^h)^{-1}\overline{B}_I = B_I^{\top}(P_I^h)^{-\top}C_I^{\top}$, then (4.3) follows using $F = -F^{\top}$. Using the right-invertibility, (4.4) follows from [23, equ. (70)] and the fact that X is full rank. 2) From passivity any solution must be symmetric. The Lyapunov equation and the skew-symmetry of F yield (4.5). \boxtimes

We see that the right-invertibility condition (resp. condition (4.4)) in item 1 of the lemma, characterizes the couplings between x and λ in the index reduction (resp. between x and z in the relative degree reduction (*i.e.*, in \dot{w})).

4.1. Analysis of the OSNSP for (3.2). The next proposition analyses the OSNSP well-posedness, and applies to two main cases: $\tilde{D} = 0$, and $\tilde{D} \succ 0$. Let us remind that $p = n_2$ as a consequence of Assumptions 2, 3, 4 and 5, so that \tilde{A} and B_{II} are $p \times p$ matrices.

244 PROPOSITION 4.3. Let h > 0, $k \ge 0$ be given.

- 1. Let Assumption 4 hold. Suppose that either $\operatorname{Im}(C) \cap \operatorname{int}(\operatorname{Im}(\partial \sigma_{\Gamma_k} + D)) \neq \emptyset$ or \tilde{B} is square full rank. Suppose also that there is ζ_0 such that $\tilde{C}\zeta_0 \in$ Im $(\partial \sigma_{\Gamma_k} + \tilde{D})$. Then the operator $\mathcal{G} : \zeta \mapsto -\tilde{A}\zeta + \tilde{B}(\partial \sigma_{\Gamma_k} + \tilde{D})^{-1}(\tilde{C}\zeta)$ is maximal monotone with non empty domain.
 - 2. Let $\tilde{D} \succ 0$, then both mappings $(\mathcal{N}_{K^*} + \tilde{D})^{-1}(\cdot)$ and $(\partial \sigma_{\Gamma_k} + \tilde{D})^{-1}(\cdot)$ are well-defined, single-valued Lipschitz continuous and maximal monotone for each $h \ge 0$, $k \ge 0$ (maximal monotonicity holds for $\tilde{D} \succeq 0$). The operator $\mathcal{G}(\cdot)$ is single-valued and Lipschitz continuous also.
 - 3. Let $\tilde{D} = 0$, Assumption 2 hold true, $\tilde{R}^2 = \tilde{X}$, $\tilde{R} = \tilde{R}^\top \succ 0$, $\operatorname{Im}(\tilde{C}) \cap \Gamma_k \neq \emptyset$ and K = P(G, 0) (see Proposition A.1 e)) for some matrix $G \in \mathbb{R}^{m \times m}$. Let us assume that $\tilde{R}\tilde{A}\tilde{R}^{-1} \preccurlyeq 0$. Let us consider the set:

256 (4.6)
$$\tilde{\mathcal{S}} \stackrel{\Delta}{=} P(G\tilde{C}\tilde{R}^{-1}, 0) \\ \cap \{\xi \in \mathbb{R}^p | \tilde{R}\tilde{A}\tilde{R}^{-1}\xi = -\sum_{i=1}^m \lambda_i (G\tilde{C}\tilde{R}^{-1})_{i\bullet}^\top, \ \lambda_i \ge 0 \} \\ \cap \operatorname{Ker}(\tilde{R}\tilde{A}\tilde{R}^{-1} + \tilde{R}^{-1}\tilde{A}^\top \tilde{R}).$$

257 Define $\tilde{q}_k = -\tilde{R}\tilde{E}_k$, $\tilde{f}(\cdot) = \psi_{\Gamma_k} \circ \tilde{C}\tilde{R}^{-1}(\cdot)$, $\Gamma_k = K - \tilde{F}_k$, and consider the 258 $VI(-\tilde{R}\tilde{A}\tilde{R}^{-1}, \tilde{q}_k, \tilde{f})$: Find $\xi_{k+1} = \tilde{R}z_{k+1}$ such that

259 (4.7)
$$\langle -\tilde{R}\tilde{A}\tilde{R}^{-1}\xi_{k+1} + \tilde{q}_k, v - \xi_{k+1} \rangle + \tilde{f}(v) - \tilde{f}(\xi_{k+1}) \ge 0, \text{ for all } v \in \operatorname{dom}(\tilde{f}).$$

- 260 (a) If $\tilde{S} = \{0\}$, then the VI $(-\tilde{R}\tilde{A}\tilde{R}^{-1}, \tilde{q}_k, \tilde{f})$ has at least one solution.
- 261 (b) If $\tilde{S} \neq \{0\}$, and if there exists $\xi_0 \in \tilde{\Gamma}_k$ such that

(4.8)
$$\langle \tilde{q}_k - (-\tilde{R}\tilde{A}\tilde{R}^{-1})^{\top}\xi_0, v \rangle > 0 \text{ for all } v \in \tilde{\mathcal{S}}, \ v \neq 0,$$

where $\tilde{\Gamma}_k = \{\xi \in \mathbb{R}^{n_2} \mid G\tilde{C}\tilde{R}^{-1}\xi + G\tilde{F}_k \geq 0\} = \operatorname{dom}(\tilde{f})$, then the $VI(-\tilde{R}\tilde{A}\tilde{R}^{-1}, \tilde{q}_k, \tilde{f})$ has at least one solution.

- (c) If ξ_{k+1}^1 and ξ_{k+1}^2 are two solutions of the $VI(-\tilde{R}\tilde{A}\tilde{R}^{-1}, \tilde{q}_k, \tilde{f})$, then $\xi_{k+1}^1 \xi_{k+1}^2 \in \operatorname{Ker}(\tilde{R}\tilde{A}\tilde{R}^{-1} + \tilde{R}^{-1}\tilde{A}^{\top}\tilde{R})$.
- (d) Let $\tilde{R}\tilde{A}\tilde{R}^{-1}$ be symmetric. If ξ_{k+1}^1 and ξ_{k+1}^2 are two solutions of the $VI(-\tilde{R}\tilde{A}\tilde{R}^{-1}, \tilde{q}_k, \tilde{f})$, then $\langle \tilde{q}_k, \xi_{k+1}^1 \xi_{k+1}^2 \rangle = 0$.
- (e) Let R̃ÃR⁻¹ be symmetric. Then any solution of the VI(−R̃ÃR⁻¹, q̃_k, f̃) is also a solution to the problem: min_{ξ∈Γ̃k} −¹/₂ξ^TR̃ÃR⁻¹ξ + ⟨q̃_k,ξ⟩.

(f) Let
$$-\hat{R}\hat{A}\hat{R}^{-1} \succ 0$$
 and be symmetric. Then $\xi_{k+1} = \operatorname{proj}_{-\tilde{R}\tilde{A}\tilde{R}^{-1}}[\hat{\Gamma}_k; -\tilde{q}_k]$

272 **Proof:**

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- 2731. Follows from [6, Theorem 1], which is itself inspired from [23] and [16, The-274orem 2], and which states that the negative feedback interconnection of a275passive system with a maximal monotone operator, defines another maximal276monotone operator (under the stated basic conditions).
 - 2. Follows from [16, Proposition 1] (see Proposition B.1) and from [55, Lemma 1] (or [9, Corollary 24.4 (ii)]) since $\operatorname{Im}(\tilde{D}) = \mathbb{R}^m$ and $\operatorname{int}(\operatorname{Dom}(\mathcal{N}_{K^*})) = \operatorname{int}(K^*) \neq \emptyset$, and $\operatorname{int}(\operatorname{Dom}(\partial \sigma_{\Gamma_k})) \neq \emptyset$ because Γ_k is convex and $\operatorname{int}(\Gamma_k) \neq \emptyset$.
- 3. The proof follows from Propositions A.1 and A.2. First notice that the MLCP $\tilde{A}z_{k+1} + \tilde{B}\lambda_{k+1} + \tilde{E}_k = 0, \ K^* \ni \lambda_{k+1} \perp \tilde{C}z_{k+1} + \tilde{D}\lambda_{k+1} + \tilde{F}_k \in K$, can be equivalently rewritten as the VI $(-\tilde{R}\tilde{A}\tilde{R}^{-1}, \tilde{q}_k, \tilde{f})$ in (4.7). Indeed one obtains $-\tilde{R}\tilde{A}\tilde{R}^{-1}\xi_{k+1} - \tilde{R}\tilde{E}_k \in -\tilde{R}\tilde{B} \ \partial\psi_{\Gamma_k}(\tilde{C}\tilde{R}^{-1}\xi_{k+1})$. Using the chain rule of Convex Analysis [49, Theorem 23.9] and the fact that Γ_k is polyhedral, the MLCP can written equivalently as in (4.7). Let us now compute the set S in (4.6). Here dom(φ) = dom(\tilde{f}) = { $z \in \mathbb{R}^p$ | $\tilde{C}\tilde{R}^{-1}z \in \Gamma_k$ } =

287	$\{z \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}z + \tilde{F}_k \in K\} = \{z \in \mathbb{R}^p \mid G\tilde{C}\tilde{R}^{-1}z + G\tilde{F}_k \geq 0\}.$ Hence
288	$(\operatorname{dom}(\varphi))_{\infty} = P(G\tilde{C}\tilde{R}^{-1}, 0)$ from Proposition A.1 e). Now we have $\tilde{f}(\xi) =$
289	$\psi_{\Gamma_k}(\tilde{C}\tilde{R}^{-1}\xi) = \psi_{\tilde{\Gamma}_k}(\xi) \text{ with } \tilde{\Gamma}_k = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in $
290	$\mathbb{I}\!\!R^p \mid G\tilde{C}\tilde{R}^{-1}\xi + G\tilde{F}_k \ge 0\} = P(G\tilde{C}\tilde{R}^{-1}, -G\tilde{F}_k). \text{ Thus } \tilde{f}_{\infty}(\cdot) = \psi_{(\tilde{\Gamma}_k)_{\infty}}(\cdot)$
291	and dom $(\tilde{f}_{\infty}) = P(G\tilde{C}\tilde{R}^{-1}, 0)$. Now we have to calculate the dual cone
292	$P(G\tilde{C}\tilde{R}^{-1}, 0)^{\star}$. Using [31, Example 15], one finds that $P(G\tilde{C}\tilde{R}^{-1}, 0)^{\star} = \{w \in \mathbb{C}\}$
293	$\mathbb{R}^{n_2} \mid w = \sum_{i=1}^m \lambda_i (G\tilde{C}\tilde{R}^{-1})_{i\bullet}^\top, \ \lambda_i \ge 0\},$ where $(G\tilde{C}\tilde{R}^{-1})_{i\bullet}$ denotes the <i>i</i> th
294	row of the matrix $G\tilde{C}\tilde{R}^{-1} \in \mathbb{R}^{m \times n_2}$. Hence it follows that $\mathcal{K}(-\tilde{R}\tilde{A}\tilde{R}^{-1},\tilde{f}) =$
295	$\{\xi \in \mathbb{R}^{n_2} \mid \tilde{R}\tilde{A}\tilde{R}^{-1}\xi = -\sum_{i=1}^m \lambda_i (G\tilde{C}\tilde{R}^{-1})_{i\bullet}^{\top}, \lambda_i \geq 0\}.$ The third set in \mathcal{S}
296	is equal to $\operatorname{Ker}(\tilde{R}\tilde{A}\tilde{R}^{-1} + \tilde{R}^{-1}\tilde{A}^{\top}\tilde{R})$, from $\operatorname{Ker}(\mathbf{M} + \mathbf{M}^{\top})$ in Proposition A.2
297	a). Then the results in items (a) to (f) are a consequence of Proposition A.2.
298	where $\varphi_{\infty}(v) = 0$ in item (b). Item (f) follows from item e) in Proposition
299	A.2.
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Item 1 encompasses cases with $\tilde{D} \geq 0$. If item 2 holds, then λ_{k+1} is uniquely defined, provided that z_{k+1} exists (which is guaranteed if conditions of item 1 hold). From item 3 (c), it follows that z_{k+1} is unique if \tilde{A} has full rank (recall that B_{II} and \tilde{A} are square $p \times p$ matrices). However from (3.2) and (3.3), x_{k+1} is uniquely defined if and only if $\bar{B}_I(\mathcal{N}_{K^{\star}})^{-1}(-\tilde{C}z_{k+1}-\tilde{F}_k)$ is unique. This happens trivially when $\bar{B}_I = 0$ (if one thinks of (2.1) as the limit of a singularly perturbed system, this simply means that the multiplier enters only the z-dynamics). Let us now deal with the case $\tilde{D} \geq 0$.

ASSUMPTION 6. Let
$$h > 0$$
, and $\tilde{D} = \begin{pmatrix} D_{11} & 0 \\ 0 & 0 \end{pmatrix}$ with $0 \prec \tilde{D}_{11} = \tilde{D}_{11}^{\top} \in \mathbb{R}^{m_1 \times m_1}$.

Such specific structure occurs for instance when D has the same blockdiagonal structure and $\bar{B}_I = 0$, *i.e.*, the multiplier does not intervene in the *x*-dynamics. Let us split the multiplier as $\lambda = (\lambda_1^{\top}, \lambda_2^{\top})^{\top}, \lambda_1 \in \mathbb{R}^{m_1}, \lambda_2 \in \mathbb{R}^{m_2}, m_1 + m_2 = m$, and $K = K_1 \times K_2 \subset \mathbb{R}^{m_1} \times \mathbb{R}^{m_2}$ (hence $K^* = K_1^* \times K_2^*$), $\tilde{B} = (\tilde{B}_1 \quad \tilde{B}_2), \tilde{B}_1 \in \mathbb{R}^{p \times m_1}$, $\tilde{B}_2 \in \mathbb{R}^{p \times m_2}, \tilde{C} = \begin{pmatrix} \tilde{C}_1 \\ \tilde{C}_2 \end{pmatrix}, \tilde{C}_1 \in \mathbb{R}^{m_1 \times p}, \tilde{C}_2 \in \mathbb{R}^{m_2 \times p}$. If $\tilde{D}_{11} = 0$ then $\lambda = \lambda_2$, $\tilde{B} = \tilde{B}_2, \tilde{C} = \tilde{C}_2$. If $\tilde{D} = \tilde{D}_{11} \succ 0$ then item 2 in Proposition 4.3 applies.

Assumption 7. Let h > 0. There exists a $(p \times p)$ -matrix $\tilde{X}_2 = \tilde{X}_2^\top \succ 0$ such that $\tilde{X}_2 \tilde{B}_2 = \tilde{C}_2^\top$.

This assumption implies that there are couplings between the multiplier λ_{k+1} in the algebraic equation, and the algebraic variable z_{k+1} in the complementarity conditions, see (3.2). Define $\tilde{R}_2 = \tilde{R}_2^\top \succ 0$, $\tilde{R}_2^2 = \tilde{X}_2$, and $\xi_{k+1} = \tilde{R}_2 z_{k+1}$. The next result extends

321 Proposition 4.3 to the positive semidefinite case with blockdiagonal structure.

PROPOSITION 4.4. Let Assumptions 6 and 7 hold true, h > 0 and $k \ge 0$ be given. Define $\Gamma_{2,k} = K_2 - \tilde{F}_{2,k}$. Consider the VI: find $\xi_{k+1} \in \operatorname{dom}(\tilde{f}_2)$ such that (4.9)

$$324 \quad \tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1}\xi_{k+1} - \tilde{R}_{2}\tilde{E}_{k} + \underbrace{\tilde{R}_{2}\tilde{B}_{1}(\tilde{D}_{11} + \partial\psi_{K_{1}^{\star}})^{-1}(\tilde{C}_{1}\tilde{R}_{2}^{-1}\xi_{k+1} + \tilde{F}_{1,k})}_{\triangleq \tilde{G}_{2}(t_{k},\xi_{k+1})} \in -\partial\tilde{f}_{2}(\xi_{k+1})$$

325 where $\tilde{f}_2(\cdot) = \psi_{\Gamma_{2,h}} \circ \tilde{C}_2 \tilde{R}_2^{-1}(\cdot)$.

1. Suppose that
$$\operatorname{Im}(\tilde{C}_2) \cap \Gamma_{2,k} \neq \emptyset$$
. The MLCP $\tilde{A}z_{k+1} + \tilde{B}\lambda_{k+1} + \tilde{E}_k = 0$,
 $K^* \ni \lambda_{k+1} \perp \tilde{C}z_{k+1} + \tilde{D}\lambda_{k+1} + \tilde{F}_k \in K$, can be equivalently rewritten as the

VI in (4.9).

- 2. Assume that the set $\tilde{\Gamma}_{2,k} = \{\xi \in \mathbb{R}^p \mid \tilde{C}_2 \tilde{R}_2^{-1} \xi \in \Gamma_{2,k}\}$ is compact. Then the set of solutions to the VI in (4.9) is compact and nonempty.
 - 3. Assume that for each t_k , the following well-defined, single-valued and Lipschitz continuous mapping, is maximal monotone:

$$\tilde{G}_2: \xi_{k+1} \mapsto \tilde{R}_2 \tilde{B}_1 (\tilde{D}_{11} + \partial \psi_{K_1^\star})^{-1} (\tilde{C}_1 \tilde{R}_2^{-1} \xi_{k+1} - \tilde{F}_{1,k})$$

331 If $-\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} \succeq 0$, $\tilde{\Gamma}_{2,k} \cap \operatorname{Int}(\operatorname{Dom}(\tilde{G}_2)) \neq \emptyset$ or $\operatorname{Int}(\tilde{\Gamma}_{2,k}) \cap \operatorname{Dom}(\tilde{G}_2) \neq \emptyset$, then 332 the set of solutions to the VI in (4.9) is closed and convex. If $-\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} \succ 0$, 333 the VI in (4.9) has a unique solution.

334 **Proof:**

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1. The complementarity conditions in (3.2) (c) are rewritten as:

$$(4.10) \quad \begin{pmatrix} K_1^{\star} \\ K_2^{\star} \end{pmatrix} \ni \begin{pmatrix} \lambda_{1,k+1} \\ \lambda_{2,k+1} \end{pmatrix} \perp \begin{pmatrix} \tilde{C}_1 \\ \tilde{C}_2 \end{pmatrix} z_{k+1} + \begin{pmatrix} \tilde{D}_{11}\lambda_{1,k+1} \\ 0 \end{pmatrix} + \begin{pmatrix} \tilde{F}_{1,k} \\ \tilde{F}_{2,k} \end{pmatrix} \in \begin{pmatrix} K_1 \\ K_2 \end{pmatrix}$$

while the equality constraint in (3.2) (b) is rewritten as:

338 (4.11)
$$\tilde{A}z_{k+1} + \tilde{B}_1\lambda_{1,k+1} + \tilde{B}_2\lambda_{2,k+1} + \tilde{E}_k = 0$$

339 From (4.10) it is deduced that:

(4.12)
$$\begin{cases} \lambda_{1,k+1} = -(\mathcal{N}_{K_1^*} + \tilde{D}_{11})^{-1}(\tilde{C}_1 z_{k+1} + \tilde{F}_{1,k}) \\ \lambda_{2,k+1} \in -\mathcal{N}_{K_2 - \tilde{F}_{2,k}}(\tilde{C}_2 z_{k+1}). \end{cases}$$

341 Using (4.11) it is inferred:

which is (4.9) since $\tilde{R}_2 \tilde{B}_2 = \tilde{R}_2^{-1} \tilde{X}_2 \tilde{B}_2 = \tilde{R}_2^{-1} \tilde{C}_2^{\top}$ and applying the chain rule of Convex Analysis for polyhedral functions.

- 2. This follows from [29, Corollary 2.2.5] and since the map $(D_{11} + \partial \psi_{K_1^*})^{-1}(\cdot)$ is single-valued, well defined and Lipschitz continuous (see Proposition B.1).
- 347 3. The first assertion is a consequence of [9, Corollary 24.4] (for the maximal 348 monotonicity of the mapping $(-\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} + \tilde{G}_2 + \partial \tilde{f}_2)(\cdot)$, and of [9, Proposition 349 23.39]. The second assertion follows from [29, Theorem 2.3.3], noting that the 350 mapping in the left-hand side of (4.9) is maximal monotone [9, Corollary 24.4 351 (ii)] and strongly monotone since $-\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} \succ 0$. This can also be proved 352 using [9, Corollary 23.37]. Single-valuedness and Lipschitz continuity stem 353 from $\tilde{D}_{11} \succ 0$.

Further features of the solutions to the VI in (4.9) will be studied in section 5.1.3. The next lemma states a classical condition that guarantees the maximal monotonicity of Proposition 4.4 item 3. The first condition in item 1 is a constraint qualification which guarantees that the function $\tilde{f}_2(\cdot)$ has a nonempty domain. It is usually seen as a technical condition, however in some applications it may not be satisfied. Item 1 states conditions under which the VI in (4.9) can be constructed. Item 3 uses some results about the zeroes of maximal monotone mappings, or of the sum of these.

LEMMA 4.5. Let $\tilde{X}_2\tilde{B}_1 = \tilde{C}_1^{\top}$ for \tilde{X}_2 in Assumption 7, then for each fixed t_k , the mapping $\tilde{G}_2 : \xi_{k+1} \mapsto \tilde{R}_2\tilde{B}_1(\tilde{D}_{11} + \partial \psi_{K_1^{\star}})^{-1}(\tilde{C}_1\tilde{R}_2^{-1}\xi_{k+1} + \tilde{F}_{1,k})$ is maximal monotone.

Proof: Indeed $\tilde{X}_2\tilde{B}_1 = \tilde{C}_1^{\top} \Rightarrow \tilde{R}_2\tilde{B}_1 = \tilde{R}_2^{-1}\tilde{C}_1^{\top} = (\tilde{C}_1\tilde{R}_2^{-1})^{\top}$. The mapping $(\tilde{D}_{11} + \partial\psi_{K_1^*})^{-1}(\cdot)$ is maximal monotone, and the following operation: $(\tilde{C}_1\tilde{R}_2^{-1})^{\top}(\tilde{D}_{11} + \partial\psi_{K_1^*})^{-1}(\tilde{C}_1\tilde{R}_2^{-1} \cdot + \tilde{F}_{1,k})$ preserves the maximal monotonicity [50, Exercise 12.4].

Propositions 4.3 and 4.4 provide some informations about the existence and uniqueness of the unknown z_{k+1} , and characterize some mappings. It is noteworthy that other well-posedness results for VIs are available in the literature. For instance if the set $\tilde{\Gamma}_k$ is compact, [29, Corollary 2.2.5] guarantees existence of a solution to VI $(-\tilde{R}\tilde{A}\tilde{R}^{-1}, \tilde{q}_k, \tilde{f})$ with a compact set of solutions. Also when $\tilde{\Gamma}_k = I\!\!R^m_+$ one may rely on P-matrices and LCPs properties, when it is a closed convex cone one may use

- 372 copositive matrices [29, Theorem 2.5.10].
- 373 The next corollary completes the existence and uniqueness results for the OSNSP.
- 374 COROLLARY 4.6. The next results hold true:
- 1. (set of solutions z_{k+1}) Let Assumption 4 and the conditions of item 1 in Proposition 4.3 hold. Then the set of solutions z_{k+1} in (3.5) (or the set of solutions ζ_{k+1} in (3.3)), is closed and convex. If in addition the conditions of item 3 (a) or (b) in Proposition 4.3 hold, then the set of solutions is non empty. If in addition $\tilde{A} \prec 0$, then z_{k+1} is unique.
 - (set of solutions λ_{k+1}) Let C Dom(G) ⊆ Im(∂σ_{Γk} + D), and the conditions of item 1 in Proposition 4.3 hold true. Then the set of solutions λ_{k+1} in (3.5) (a) is non empty for each z_{k+1}.
- 383 3. Let Assumption 4 hold with strict passivity, and the conditions of item 1 in 384 Proposition 4.3 be true. Then the GE in (3.5) (or the GE in (3.3)) has 385 a unique solution (λ_{k+1}, z_{k+1}) for any \tilde{E}_k . Consequently x_{k+1} in (3.2) is 386 unique.
 - 4. Let $\tilde{D} = 0$ and $\mathcal{G}(\cdot)$ be maximal monotone. Then if the set $\tilde{\Gamma}_k = \{\zeta \in \mathbb{R}^{n_2} \mid G\tilde{C}\tilde{R}^{-1}\zeta + G\tilde{F}_k \geq 0\}$ is bounded, the GE with unknown z_{k+1} in (3.5) (or in (3.3)) has at least one solution and the set of solutions is compact.

390 Proof:

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- 1. Follows from [9, Proposition 23.39] by the maximal monotonicity of $\mathcal{G}(\cdot)$. If $-\tilde{A} \succ 0$ then $\mathcal{G}(\cdot)$ is strongly monotone (however it may not be Lipschitz continuous and this property is used later in section 5.2), and so is $\mathcal{G}(\cdot) + \tilde{E}_k$. Hence the generalized equation $\mathcal{G}(z_{k+1}) + \tilde{E}_k \ni 0$ has a unique solution [9, Corollary 23.37].
 - 2. Let us consider $\lambda_{k+1} \in -(\partial \sigma_{\Gamma_k} + \tilde{D})^{-1}(\tilde{C}z_{k+1})$. We know that any solution $z_{k+1} \in \mathcal{G}^{-1}(\tilde{E}_k) \subseteq \operatorname{Im}(\mathcal{G}^{-1}) = \operatorname{Dom}(\mathcal{G})$. Since $\operatorname{Dom}(\partial \sigma_{\Gamma_k} + \tilde{D})^{-1} = \operatorname{Im}(\partial \sigma_{\Gamma_k} + \tilde{D})$ the result follows.
 - 3. Strict passivity implies that $\tilde{D} \succ 0$ and \tilde{A} is a Hurwitz matrix, therefore from a slight modification of [6, Theorem 1] it follows that the operator $\mathcal{G}(\cdot)$ is strongly, hence strictly monotone (indeed there exists a real $\alpha > 0$ small enough such that $(\tilde{A} + \alpha I_{n_2}, \tilde{B}, \tilde{C}, \tilde{D})$ is passive). Using [9, Proposition 23.35] shows that z_{k+1} is unique, and item 2 in Proposition 4.3 shows that λ_{k+1} is unique. The conclusion follows.
- 405 4. Under these conditions $Dom(\mathcal{G})$ is bounded (compact). Hence from [9, Propo-406 sition 23.36 (iii)] existence of a solution holds true. From [29, Corollary 2.2.5] 407 the compactness of the set of solutions follows. This could also be inferred 408 from item (e) in Proposition 4.3 in case $\tilde{R}\tilde{A}\tilde{R}^{-1}$ is symmetric.

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Item 3 in Corollary 4.6 guarantees the uniqueness of solutions to the OSNSP in (3.6), 410in which case the \in can be replaced by equality. Concerning item 2 in Corollary 4.6, 411 an important step is to characterize both $\text{Dom}(\mathcal{G})$ and $\text{Im}(\partial \sigma_{\Gamma_k} + \tilde{D})$. If $\tilde{D} = \tilde{D}^{\top} \succeq 0$, 412 one can apply [15, Theorem 3], and infer that $\operatorname{Im}(\partial \sigma_{\Gamma_k} + \tilde{D}) \simeq \operatorname{Im}(\partial \sigma_{\Gamma_k}) + \operatorname{Im}(\tilde{D}) =$ 413 $\operatorname{Dom}(\partial \psi_{\Gamma_k}) + \operatorname{Im}(\tilde{D}) \subset \operatorname{dom}(\psi_{\Gamma_k}) + \operatorname{Im}(\tilde{D}) = \Gamma_k + \operatorname{Im}(\tilde{D})$ (see [15] for the exact 414 meaning of \simeq in this context). Similar calculations may be done for in item 1 of 415Proposition 4.3. It is inferred that the basic condition in Proposition 4.3 item 1, 416 *i.e.*, $\operatorname{Im}(C) \cap \operatorname{int}(\operatorname{Im}(\partial \sigma_{\Gamma_k} + D)) \neq \emptyset$, or in Corollary 4.6 item 2, *i.e.*, $C \operatorname{Dom}(\mathcal{G}) \subseteq \mathcal{I}$ 417 $\operatorname{Im}(\partial \sigma_{\Gamma_k} + D)$, always are verified when $D \succ 0$. For the sake of briefness the conditions 418 for the existence and uniqueness of λ_{k+1} in (4.12) are not stated, similar tools as for 419 Corollary 4.6 can be used. 420

421 **4.2.** Analysis of the OSNSP for (3.1). It is also interesting to investigate 422 the case when Assumption 3 (resp. Assumption 5) is satisfied instead of Assumption 423 2 (resp. Assumption 4). The operator $\mathcal{G}(\cdot)$ is now defined as $\mathcal{G} : \zeta \mapsto -B_{II}\zeta +$ 424 $(\partial \sigma_{\Gamma(t_k,x_{k+1})} + D)^{-1}(C_{II}\zeta)$, where $\Gamma(t_k,x_{k+1}) = \{\xi \in \mathbb{R}^m \mid \xi + C_I x_{k+1} + F_k \in$ 425 $K\}$. Consequently the results of Proposition 4.3 item 1, Corollary 4.6 item 3, with 426 Assumption 5 plus strict passivity, and Proposition 4.3 item 2 when $D \succ 0$, hold. For 427 the sake of briefness we do not develop all the arguments once again.

428 Let us focus now on the case $D \geq 0$. From Assumption 3, D = 0 and (2.1) one infers:

429 (4.14)
$$-RB_{II}R^{-1}\xi(t) - RA_{II}x(t) - RE_{II}(t) \in -\partial\psi_{\Phi(t,x(t))}(\xi(t)),$$

430 where $\xi = -Rz$, $R^2 = X$, $R = R^{\top} \succ 0$, $\Phi(t, x) = \{\xi \in \mathbb{R}^p \mid C_{II}R^{-1}\xi + C_Ix + F(t) \in K\}$, and the set-valued term is obtained using the chain rule, under the Convex 432 Analysis chain rule qualification condition [49, Theorem 23.9]. The first step is to 433 analyse the existence and uniqueness of solutions to the VI:

434 (4.15)
$$-RB_{II}R^{-1}\xi_{k+1} - RA_{II}x_{k+1} - RE_{II,k} \in -\partial\psi_{\Phi(t_k,x_{k+1})}(\xi_{k+1}),$$

which is the counterpart of (4.7), relying on Proposition A.2. The unknown is ξ_{k+1} and the other terms and variables are considered as data for this VI. Similar results to those in item 3 of Proposition 4.3 can be obtained.

438 PROPOSITION 4.7. Let D = 0, Assumption 3 be true, h > 0 be given, assume that 439 Im $(C_{II}) \cap (K - C_I x_{k+1} + F_k) \neq \emptyset$ and that K is polyhedral. Consider the discrete-time 440 system in (3.1), and the associated $VI(-RB_{II}R^{-1}, -RA_{II}x_{k+1} - RE_{II,k}, \psi_{\Phi(t_k, x_{k+1})})$ 441 in (4.15). Assume that $-RB_{II}R^{-1} \geq 0$. Consider the set: (4.16)

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$$\mathcal{S} \stackrel{\Delta}{=} P(GC_{II}R^{-1}, 0) \cap \{\xi \in \mathbb{R}^{n_2} | RB_{II}R^{-1}\xi = \sum_{i=1}^m \lambda_i (GC_{II}R^{-1})_{i\bullet}^\top, \ \lambda_i \ge 0\}$$
$$\cap \operatorname{Ker}(RB_{II}R^{-1} + R^{-1}B_{II}^\top R).$$

443 *Then:*

444	1. If $S = \{0\}$ the VI $(-RB_{II}R^{-1}, -RA_{II}x_{k+1} - RE_{II,k}, \psi_{\Phi(t_k, x_{k+1})})$ has a unique
445	solution.

446 2. If $S \neq \{0\}$, and if there exists $\xi_0 \in \Phi(t_k, x_{k+1})$ such that $\langle -RA_{II}x_{k+1} - RE_{II,k} + RB_{II}R^{-1}\xi_0, v \rangle > 0$ for all $v \neq 0$ and $v \in S$, then it follows that the 448 $VI(-RB_{II}R^{-1}, -RA_{II}x_{k+1} - \tilde{R}E_{II,k}, \psi_{\Phi(t_k, x_{k+1})})$ has a unique solution.

Proof: The set S in (4.16) is constructed as follows. The chain rule holds from 449 [49, Theorem 23.9] and the fact that $\psi_{K-C_I x_{k+1}+F_k}(\cdot)$ is polyhedral. Hence the 450VI in (4.15) is equivalent to (3.1) (b) (c). Here $\varphi = \psi_{\Phi(t_k, x_{k+1})}$, hence dom(φ) = 451 $\Phi(t_k, x_{k+1})$ and $(\operatorname{dom}(\varphi))_{\infty}$. Assume that K = P(G, 0) (see Proposition A.1 (e)), then $(\operatorname{dom}(\varphi))_{\infty} = \{\xi \in \mathbb{R}^p \mid GC_{II}R^{-1}\xi \geq 0\} = P(GC_{II}R^{-1}, 0)$. Further, $\varphi_{\infty} = P(GC_{II}R^{-1}, 0)$. 452453 $\begin{aligned} \psi_{(\Phi(t_k, x_{k+1}))_{\infty}} &= \psi_{P(GC_{II}R^{-1}, 0)}, \text{ hence } \operatorname{dom}(\varphi_{\infty}) = P(GC_{II}R^{-1}, 0) \text{ and } \operatorname{dom}(\varphi_{\infty})^{\star} = \\ (P(GC_{II}R^{-1}, 0))^{\star} &= \{ w \in \mathbb{R}^{n_2} \mid w = -\sum_{i=1}^m \lambda_i (GC_{II}R^{-1})_{i \bullet}^{\top}, \ \lambda_i \geq 0 \}, \text{ where } [31, 1] \end{aligned}$ 454455Example 15] was used as in the proof of Proposition 4.3. Thus $\mathcal{K}(-RB_{II}R^{-1},\varphi) =$ 456 $\{\xi \in \mathbb{R}^{n_2} \mid -RB_{II}R^{-1}\xi = -\sum_{i=1}^m \lambda_i (GC_{II}R^{-1})_{i\bullet}^{\top}, \lambda_i \geq 0\}.$ Finally in Proposition A.2: Ker $(\mathbf{M} + \mathbf{M}^{\top}) = \text{Ker}(RB_{II}R^{-1} + (RB_{II}R^{-1})^{\top})$ and $R = R^{\top}$. Items 1 and 2 457458459are a direct consequence of Proposition A.2 (a) and (b). \boxtimes

For the sake of briefness we do not state similar results as in Proposition 4.3 (c) to 460(f). To enlarge the analysis let us now state the counterpart of Assumption 6: 461

462 ASSUMPTION 8. Let
$$D = \begin{pmatrix} D_{11} & 0 \\ 0 & 0 \end{pmatrix}$$
 with $0 \prec D_{11} = D_{11}^{\top} \in \mathbb{R}^{m_1 \times m_1}$.

Similarly to Assumption 7, with the same splitting of the multiplier λ , the matrices 463 B_{II} and C_{II} , and of the cone K (see the comments after Assumption 6, in particular 464if $D_{11} = 0$ the case D = 0 is recovered with $G_2(\cdot, \cdot)$ in (4.17) which vanishes, and the 465set $\Phi_2 = \Phi$), let us state the following: 466

ASSUMPTION 9. There exists a $(p \times p)$ -matrix $X_2 = X_2^{\top} \succ 0$ such that $X_2 \bar{B}_{II}^2 =$ 467 $(C_{II}^2)^+$. 468

Let us recall that p is the number of equality constraints in (2.1). Let us define 469 $R_2 = R_2^{\top} \succ 0, R_2^2 = X_2, \xi_{k+1} = R_2 z_{k+1}, \Phi_2(t_k, x_{k+1}) = \{\xi \in \mathbb{R}^p \mid C_{II}^2 R_2^{-1} \xi + C_I^2 x_{k+1} + F_{2,k} \in K_2\}.$ The next proposition is the counterpart of Proposition 4.4. It 470 471states conditions that guarantee that z_{k+1} can be calculated from (3.1) (b) (c). 472

PROPOSITION 4.8. Let h > 0, and Assumptions 8 and 9 hold true. Consider the 473 VI: find $\xi_{k+1} \in \text{dom}(f_2)$ such that 474

475 (4.17)
$$-R_2 B_{II} R_2^{-1} \xi_{k+1} - R_2 A_{II} x_{k+1} - R_2 E_{II,k} + \underbrace{R_2 \bar{B}_{II}^1 (D_{11} + \partial \psi_{K_1^\star})^{-1} (C_{II}^1 R_2^{-1} \xi_{k+1} - F_{1,k})}_{\bullet} \in -\partial f_2(\xi_{k+1})$$

$$\stackrel{\Delta}{=} G_2(t_k,\xi_{k+1})$$

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with $f_2(\cdot) = (\psi_{K_2 - C_I^2 x_{k+1} - F_{2,k}} \circ C_{II}^2 R_2^{-1})(\cdot) = \psi_{\Phi_2(t_k, x_{k+1})}(\cdot).$ 1. Assume that $\operatorname{Im}(C_{II}^2) \cap (K_2 - C_I^2 x_{k+1} - F_{2,k}) \neq \emptyset$, then the VI in (4.17) is 477 equivalent to the MLCP in (3.1) (b) (c). 478

2. Suppose that the set $\Phi_2(t_k, x_{k+1})$ is compact and nonempty, then the set of 479solutions to the VI in (4.17) is compact and nonempty. 480

3. Assume that for each fixed t_k , the mapping $G_2(\cdot)$ is maximal monotone, and 481 that $-R_2B_{II}R_2^{-1} \succ 0$, then the VI in (4.17) has a unique solution. 482

4. Assume that $-R_2B_{II}R_2^{-1} \geq 0$, $G_2(\cdot)$ is maximal monotone for each fixed t_k , 483 and $\Phi_2(t_k, x_{k+1}) \cap \operatorname{Int}(\operatorname{Dom}(G_2)) \neq \emptyset$ or $\operatorname{Int}(\Phi_2(t_k, x_{k+1})) \cap \operatorname{Dom}(G_2) \neq \emptyset$. 484Then the set of solutions to the VI in (4.17) is closed and convex. 485

Proof: The proof of the first three items follows closely that of Proposition 4.4 and it 486is therefore omitted. Item 4: the VI is rewritten as $(-R_2B_{II}R_2^{-1} + G_2 + \partial f_2)(\xi_{k+1}) - (\xi_{k+1}) - ($ 487

$$\triangleq H$$

 $R_2A_{II}x_{k+1} - R_2E_{II,k} \ge 0$, and the assumption guarantees that $H_2(\cdot)$ is maximal 488 monotone [9, Corollary 24.4]. The result follows from [9, Proposition 23.39]. \boxtimes 489

The same as Lemma 4.5 can be stated to guarantee the maximal monotonicity in 490 491 items 3 and 4. It is noteworthy that the main discrepancy between Proposition 4.4 and Proposition 4.8, is that the solution ξ_{k+1} depends on x_{k+1} in the latter, 492 while it depends on x_k in the former. This is a difficulty with the application of 493Proposition 4.8 in some cases, for instance when $C_{II} = 0$ (examples in sections 6.2) 494 and 6.4 will demonstrate that this can occur). Actually, the condition $\operatorname{Im}(C_{II}^2) \cap (K_2 - K_2)$ 495 $C_I^2 x_{k+1} - F_{2,k} \neq \emptyset$ (that is a form of constraint qualification), or its counterparts in 496 Propositions 4.3, 4.4, 4.7, happens to be crucial in the setting of this work. It states 497 that the algebraic variable z has to enter the complementarity constraints through 498a C_{II} matrix which has a big enough range. This may occur in some systems, see 499section 6.5. 500

5. Convergence results. In this section it is assumed that the conditions for existence of solutions in Propositions 4.3, 4.4, 4.7, or 4.8, are fulfilled, so that sequences can be calculated. The section is divided into two main parts: in sections 5.1.1, 5.1.2 and 5.1.3, convergence is analysed from the properties of the quadruple $(\tilde{A}, \tilde{B}, \tilde{C}, \tilde{D})$ which characterize the MLCP in (3.2) (b) (c). In section 5.2, the analysis is led with the properties of $(B_{II}, \bar{B}_{II}, C_{II}, D)$ which characterize the MLCP in (3.1) (b) (c). The relationships between both approaches are explained.

508 5.1. Convergence results from (3.2). In this section the numerical scheme 509 in (3.2) and the results in section 4.1 are used. First let us notice that z_{k+1} and 510 λ_{k+1} are solutions of the generalized equation in (3.5) (a) (b). Thus if existence of 511 solutions hold, they both are functions of \tilde{E}_k (in addition to being functions of \tilde{F}_k 512 through the set Γ_k). Under some conditions which are recalled below, like item 2 of 513 Proposition 4.3, they can be single-valued Lipschitz continuous functions of \tilde{E}_k . Then 514 the right-hand side of the first line in (3.2) may be written compactly as

515 (5.1)
$$P_I^h x_{k+1} = P_I x_k + h \bar{E}_k(t_k, x_k)$$

for some $\overline{E}_k(t_k, x_k)$ which is a linear function of $E_{I,k}$, $E_{II,k}$, λ_{k+1} and z_{k+1} . The time-dependence of \overline{E}_k comes from $E_{I,k}$, $E_{II,k}$ and F_k , while the state dependence comes from the third terms in \tilde{E}_k and in \tilde{F}_k (which enters $\Gamma_k(t_k, x_k) = K - \tilde{F}_k$). From (5.1) it can be calculated and checked by induction that:

520 (5.2)
$$x_k = \left((P_I^h)^{-1} P_I \right)^k x_0 + h \sum_{j=0}^{k-1} \left((P_I^h)^{-1} P_I \right)^{k-1-j} (P_I^h)^{-1} \bar{E}_j(t_j, x_j), \ k \ge 1,$$

521 Since $P_I^h = P_I - hA_I$, it follows that for h > 0 small enough, $(P_I^h)^{-1}P_I = I_p + \mathcal{O}(h)$ 522 [10, Proposition 9.4.13] hence $(P_I^h)^{-1} = P_I^{-1} + \mathcal{O}(h)$. Thus

$$x_{k} = x_{0} + \mathcal{O}(h) + h \sum_{j=0}^{k-1} (I_{p} + \mathcal{O}(h))^{k-1-j} (P_{I}^{-1} + \mathcal{O}(h)) \bar{E}_{j}(t_{j}, x_{j})$$

$$= x_{0} + \mathcal{O}(h) + h \sum_{j=0}^{k-1} (I_{p} + \mathcal{O}(h)) (P_{I}^{-1} + \mathcal{O}(h)) \bar{E}_{j}(t_{j}, x_{j})$$

$$= x_{0} + \mathcal{O}(h) + h \sum_{j=0}^{k-1} (P_{I}^{-1} + \mathcal{O}(h)) \bar{E}_{j}(t_{j}, x_{j})$$

$$= x_{0} + \mathcal{O}(h) + h (P_{I}^{-1} + \mathcal{O}(h)) \sum_{j=0}^{k-1} \bar{E}_{j}(t_{j}, x_{j}),$$

524 for all $k \ge 1$. It is inferred that:

(5.3)

525 (5.4)
$$\dot{x}_k \stackrel{\Delta}{=} \frac{x_k - x_{k-1}}{h} = \frac{\mathcal{O}(h)}{h} + (P_I^{-1} + \mathcal{O}(h))\bar{E}(t_{k-1}, x_{k-1}) \\ = \mathcal{O}(1) + (P_I^{-1} + \mathcal{O}(h))\bar{E}(t_{k-1}, x_{k-1}).$$

Our aim is now to prove that both $\{x_k\}$ and $\{\dot{x}_k\}$ are uniformly bounded sequences, which will allow us to conclude about the convergence of the sequence $\{x^h(\cdot)\}$ (see the definition in (5.7)). For this it is necessary and sufficient that $\bar{E}_k(t_k, x_k)$ be uniformly bounded for all $k \ge 0$, since P_I is full rank. It follows from (3.5) (a) and (b), that there are two operators which are crucial to characterize the right-hand side in (3.6).

5.1.1. Analysis of the operator $\mathcal{G}(\cdot)$. Notice that the mapping $\mathcal{Z}: \zeta \mapsto \eta =$ 531 $(\partial \sigma_{\Gamma_k} + \tilde{D})^{-1}(\zeta)$ is under the conditions of Proposition 4.3 item 2, a well-defined 532 $\mathbb{R}^m \to \mathbb{R}^m$ single-valued Lipschitz continuous mapping which is parameterized by t_k and x_k through the set $\Gamma_k(t_k, x_k)$. Assume that x_k is bounded, then redoing 534the proof of Proposition B.1, it follows that the Lipschitz constant of $\mathcal{Z}(\cdot)$ does not depend on $\partial \sigma_{\Gamma_k}(\cdot)$ but depends only on \tilde{D} , being equal to $\frac{2}{\lambda_{\min}(\tilde{D}+\tilde{D}^{\top})}$. In turn 536 $\lambda_{\min}(\tilde{D}+\tilde{D}^{\top}) > 0$ uniformly in h > 0 only if $D \succ 0$. Thus the operator \mathcal{G} : 537 $\zeta \mapsto -\tilde{A}\zeta + \tilde{B}(\partial \sigma_{\Gamma_k} + \tilde{D})^{-1}(\tilde{C}\zeta)$ is also single-valued Lipschitz continuous (being the 538 composition of Lipschitz continuous mappings) with a Lipschitz constant that does 539 not depend on x_k , and which is bounded for bounded h. Let us now assume the 540following which reinforces Assumption 4 and guarantees that $D \succ 0$: 541

542 ASSUMPTION 10. The system $(\tilde{A}, \tilde{B}, \tilde{C}, \tilde{D})$ is strictly passive, uniformly in $h \ge 0$. 543

Then if the basic conditions stated in item 1 of Proposition 4.3 are verified in addition to Assumption 10, the operator $\mathcal{G}(\cdot)$ is strongly monotone for some constant $\epsilon > 0$, and we know from Corollary 4.6 item 3, that x_{k+1} is unique. As stated in Proposition 5464.3 item 1, maximal monotonicity follows from [6, Theorem 1]. Strong monotonicity 547548follows from the fact that under Assumption 10, the quadruple $(A + \epsilon I_{n_2}, B, C, D)$ is passive for a small enough $\epsilon > 0$, so that $(\mathcal{G} - \epsilon I_{n_2})(\cdot)$ is maximal monotone for 549such ϵ . Thus if Assumption 10 and the basic conditions in item 1 of Proposition 5504.3 are true, $\mathcal{G}: \mathbb{R}^p \to \mathbb{R}^p$ is well-defined and single-valued, Lipschitz continuous 551and strongly monotone. Consequently it is inferred from [50, Proposition 12.54] that under the above set of conditions, $\mathcal{G}^{-1}(\cdot)$ is single-valued, well defined and globally 553Lipschitz, and the solution z_{k+1} of the generalized equation $E_k(t_k, x_k) \in \mathcal{G}(z_{k+1}) \Leftrightarrow$ 554 $z_{k+1} = \mathcal{G}^{-1}(\tilde{E}_k(t_k, x_k))$ in (3.5) is Lipschitz continuous in x_k and t_k . Thus λ_{k+1} is 555also Lipschitz continuous in x_k , and we conclude that the same applies to x_{k+1} in 556(3.2).

Remark 5.1. Assumption 10 is stringent. Actually the strict state passivity of $(\tilde{A}, \tilde{B}, \tilde{C}, \tilde{D})$ is sufficient to guarantee that $\mathcal{G}(\cdot)$ is strongly monotone. However in the body of sections 5.1.1 and 5.1.2 we need also that λ_{k+1} be a Lipschitz continuous function of its argument, a sufficient condition being that $\tilde{D} \succ 0$, which in turn implies that strict state passivity and strict passivity are identical. Thus strict passivity could be relaxed to strict state passivity if it is assured that $\bar{B}_I(\partial\sigma_{\Gamma_k}+\tilde{D})^{-1}(\tilde{C}\cdot)$ is Lipschitz continuous. This may occur in certain cases (the trivial case is $\bar{B}_I = 0$).

565 **5.1.2. Boundedness and convergence analysis.** Under the above conditions, 566 the term $\bar{E}_k(t_k, x_k)$ and consequently $x_{k+1}(t_k, x_k)$ in (3.2) and $\dot{x}_k(t_k, x_k)$ in (5.4) are 567 single-valued and Lipschitz continuous in x_k and t_k (provided $E_I(\cdot)$, $E_{II}(\cdot)$ and $F(\cdot)$ 568 are Lipschitz continuous functions of t, which is a mild requirement on the data), 569 since \tilde{E}_k and \tilde{F}_k are linear in the argument x_k . We have $||\bar{E}(t_k^1, x_k^1) - \bar{E}(t_k^2, x_k^2)|| \leq$ 570 $L_t|t_k^1 - t_k^2| + L_x||x_k^1 - x_k^2||$ for any $t_k^1, t_k^2, x_k^1, x_k^2$, and for some $L_t > 0$ and $L_x > 0$. Thus 571 $||\bar{E}(t_k, x_k^1) - \bar{E}(t_k, x_k^2)|| \leq L_x||x_k^1 - x_k^2||$ and $||\bar{E}(t_k, x_k) - \bar{E}(t_k, 0)|| \leq L_x||x_k||$. Notice 572 that due to the Lipschitz continuity both z_{k+1} and λ_{k+1} are uniformly bounded in

- 573 time for fixed x_k (provided $E_I(\cdot)$, $E_{II}(\cdot)$ and $F(\cdot)$ are uniformly bounded, which is
- a mild requirement on the data). Therefore $\bar{E}(t_k, 0)$ is uniformly bounded also, and
- 575 $||\bar{E}(t_k, x_k)|| \le ||\bar{E}(t_k, 0)|| + L_x ||x_k|| \le m + L_x ||x_k||$. Now using (5.2) we get:
- (5.5) 576 $||x_k|| \le \sigma_{\max}^k((P_I^h)^{-1}P_I)||x_0|| + h\sigma_{\min}(P_I^h)\sum_{i=1}^{k-1}\sigma_{\max}^{k-j-1}(P_I(P_I^h)^{-1})(m+L_x||x_{j-1}||),$
- where we used the Hölder induced matrix norm $||\cdot||_{2,2} = \sigma_{\max}(\cdot)$ and [10, Proposition 9.4.9, Corollary 9.4.12]. Let us make the following:
- 579 ASSUMPTION 11. One has $\sigma_{\max}((P_I^h)^{-1}P_I) \leq 1$ for all $h \geq 0$.

Assumption 11 is a technical assumption which is satisfied if for instance the $n_1 \times n_1$ 580matrix $A_I \prec 0$ and $P_I \succ 0$. Indeed this implies that $P_I^h \succ P_I$, thus $(P_I^h)^{-1} P_I \prec I_{n_1}$. 581Notice that under Assumption 11 one no longer needs h > 0 small enough to obtain 582 the developments in (5.4), because the spectral radius of $(P_I^h)^{-1}P_I$ is smaller than 583one. It is also noteworthy that changing $A_I x_{k+1}$ to $A_I x_k$ in the right-hand side of 584(3.1) (a), allows one to dispense with Assumption 11, with a slight modification of the 585matrices A, B, C, D, E_k and F_k (where $(P_I^h)^{-1}$ is replaced by P_I^{-1} in all of them). 586 Consequently from (5.5) one infers that: 587

588 (5.6)
$$||x_k|| \le ||x_0|| + h\sigma_{\min}(P_I^h) \sum_{j=1}^k (m + L_x ||x_{j-1}||),$$

Assume now that $x_0, x_1, ..., x_{k-1}$ are uniformly bounded by some x_{\max} , then from (5.6) one has $||x_k|| \le x_{\max} + h\sigma_{\min}(P_I^h)k(m+L_xx_{\max}) \le x_{\max} + h\sigma_{\min}(P_I^h)n(m+L_xx_{\max}) =$ $x_{\max} + h\sigma_{\min}(P_I^h)\frac{T}{h}(m+L_xx_{\max}) = x_{\max} + \sigma_{\min}(P_I^h)T(m+L_xx_{\max})$. Therefore x_k is also uniformly bounded (its upperbound depends only on T which is fixed), and we conclude that all $x_k, 1 \le k \le n$, are bounded provided x_0 is bounded. Then it is inferred from (5.4) that \dot{x}_k is also uniformly bounded. Let us define the sequences of piecewise-linear approximations $\{x^h(\cdot)\}$ and $\{\dot{x}^h(\cdot)\}$ as follows:

596 (5.7)
$$\begin{cases} x^{h}(t) = x_{k+1} + \frac{t_{k+1}-t}{h}(x_{k} - x_{k+1}) \\ \dot{x}^{h}(t) = \frac{x_{k+1}-x_{k}}{h} \end{cases} \quad t \in [t_{k}, t_{k+1}).$$

We have shown that $x^{h}(\cdot)$ and $\dot{x}^{h}(\cdot)$ are uniformly bounded on [0,T] and for any h > 0, hence $x^{h}(\cdot)$ has a Lipschitz constant that does not depend on h, consequently 598the family $\{x^h\}_{h\geq 0}$ is equi-Lipschitz continuous. Thus from Arzela-Ascoli Theorem 599 $\{x^{h}(\cdot)\}$ stays in a compact subset of $\mathcal{C}^{0}([0,T];\mathbb{R}^{n_{1}})$ and there exist subsequences 600 which converge uniformly in $\mathcal{C}^0([0,T];\mathbb{R}^{n_1})$ towards a Lipschitz continuous limit $x(\cdot)$. 601 Let $\tilde{E}^{h}(\cdot)$ and $\tilde{F}^{h}(\cdot)$ be defined in a similar way as $x^{h}(\cdot)$ in (5.7). Due to the Lip-602 schitz continuity of $E_I(\cdot)$, $E_{II}(\cdot)$, and $F(\cdot)$, the sequences of functions $\{\tilde{E}^h(\cdot)\}_{h>0}$ 603 and $\{\tilde{F}^{h}(\cdot)\}_{h\geq 0}$ also converge uniformy in $\mathcal{C}^{0}([0,T];\mathbb{R}^{p})$ and $\mathcal{C}^{0}([0,T];\mathbb{R}^{m})$ to limits 604 $\tilde{E}(\cdot) = E_{II}(\cdot) + A_{II}x(\cdot)$ and $\tilde{F}(\cdot) = F(\cdot) + C_Ix(\cdot)$ as $h \to 0$, respectively. From (5.4) it 605 is deduced that $\dot{x}^h(\cdot)$ is uniformly bounded. Therefore from the Banach-Alaoglu The-606 orem, it follows that $\dot{x}^h(\cdot)$ converges weakly^{*} to a limit $v(\cdot)$. Since the domain [0,T] is 607 bounded, one has $\mathcal{L}^2([0,T];\mathbb{R}^{n_1}) \subseteq \mathcal{L}^1([0,T];\mathbb{R}^{n_1})$ [36, Corollary 1, p.80], and conse-608 quently $\dot{x}^h(\cdot)$ converges weakly in $\mathcal{L}^2([0,T];\mathbb{R}^{n_1})$. Since $x^h(t) = x^h(0) + \int_0^t \dot{x}^h(\tau) d\tau$, 609 we infer that $v(\cdot) = \dot{x}(\cdot)$ almost everywhere. 610

611 Now notice that using (3.2) the following holds:

612 (5.8)
$$\begin{cases} (a) P_I \dot{x}^h(t) = A_I x_h^*(t) + B_I z_h^*(t) + B_I \lambda_h^*(t) + E_I^h(t) \\ (b) \tilde{A} z_h^*(t) + \tilde{B} \lambda_h^*(t) + \tilde{E}_{II}^h(t) + A_{II} x^h(t) = 0 \\ (c) \ 0 \le \lambda_h^*(t) \perp \tilde{C} z_h^*(t) + \tilde{D} \lambda_h^*(t) + F^h(t) + C_I x^h(t) \ge 0, \end{cases}$$

for all
$$t \in [t_k, t_{k+1}), k \ge 0, x_h^{\star}(t) = x_{k+1}$$
 for all $t \in [t_k, t_{k+1})$ and similarly for the
other terms are step-functions. From the boundedness of $\|\dot{x}_k\| \triangleq \|\frac{x_{k+1}-x_x}{h}\| \le C$ for
some $C > 0$, we infer that $x_{\star}^h(\cdot)$ converges strongly in $\mathcal{L}^2([0,T];\mathbb{R}^{n_1})$ towards $x(\cdot)$.

(5.9)

617
$$\|x^{h} - x_{\star}^{h}\|_{\mathcal{L}_{2}([0,T];\mathbb{R}^{n_{1}})}^{2} \leq \sum_{k=0}^{n-1} \int_{t_{k}}^{t_{k+1}} (t_{k+1} - t_{k})^{2} \|\dot{x}_{k}\|^{2} dt$$
$$\leq C^{2} \sum_{k=0}^{n-1} (t_{k+1} - t_{k})^{3} = \frac{C^{2} nh^{3}}{3} = \frac{C^{2} Th^{2}}{3} \longrightarrow 0 \text{ as } h \to 0.$$

To conclude one notices that $||x^h - x^h_{\star}||^2_{\mathcal{L}_2([0,T];\mathbb{R}^{n_1})} = ||x^h \pm x - x^h_{\star}||^2_{\mathcal{L}_2([0,T];\mathbb{R}^{n_1})}$. Now we have $z^h_{\star}(t) \in \mathcal{G}^{-1}(E^h_{II,\star}(t) + A_{II}x^h_{\star}(t))$ for all $t \in [t_k, t_{k+1}), k \ge 0$. Since $E_I(\cdot)$, 618 619 $E_{II}(\cdot)$ are Lipschitz continuous, it follows that $\|\dot{E}_k\| \stackrel{\Delta}{=} \|\frac{E_{k+1}-E_k}{h}\| \leq k_E$ where k_E 620 is its Lipschitz constant. We infer that $E_{II,\star}^{h}(\cdot)$ converges strongly in $\mathcal{L}^{\tilde{2}}([0,T];\mathbb{R}^{p})$ 621 towards $E_{II}(\cdot)$, and using the Lipschitz continuity of $\mathcal{G}^{-1}(\cdot)$ it follows that $z^h_{\star}(\cdot)$ 622 converges strongly in $\mathcal{L}^2([0,T]; \mathbb{R}^p)$ towards a limit $z(\cdot) = \mathcal{G}^{-1}(E_{II}(\cdot) + A_{II}x(\cdot))$. Then from (3.3) $\lambda^h_{\star}(t) \in (\mathcal{N}_{\mathbb{R}^m_+} + \tilde{D})^{-1}(-\tilde{C}z^h_{\star}(t) - F^h_{\star}(t) + C_Ix^h_{\star}(t))$, so that $\lambda^h_{\star}(\cdot)$ 623 624 converges towards a limit $\lambda(\cdot) = (\mathcal{N}_{\mathbb{R}^m_+} + \tilde{D})^{-1}(-\tilde{C}z(\cdot) - F(\cdot) + C_I x(\cdot))$ strongly in 625 $\mathcal{L}^2([0,T]; \mathbb{R}^m)$, due to the single-valuedness, Lipschitzness and maximal monotonicity of $(\mathcal{N}_{\mathbb{R}^m_+} + \tilde{D})^{-1}(\cdot)$. It is concluded that $P_I \dot{x}^h(t) = A_I x_h^*(t) + B_I z_h^*(t) + \bar{B}_I \lambda_h^*(t) + C_I z_h^*(t) + C_I z_h^*(t$ 626 627 $E_I^h(t) \longrightarrow P_I \dot{x}(t) = A_I x(t) + B_I z(t) + \bar{B}_I \lambda(t) + E_I(t)$ weakly in $\mathcal{L}^2([0,T]; \mathbb{R}^{n_1})$ as $h \to 0$ 628 $0, \tilde{A}z_h^{\star}(t) + \tilde{B}\lambda_h^{\star}(t) + \tilde{E}_{II}^h(t) + A_{II}x^h(t) \rightarrow B_{II}z(t) + \bar{B}_{II}\lambda(t) + E_{II}(t) + A_{II}x(t) \ (=0),$ 629 and $0 \leq \lambda_h^{\star}(t) \perp \tilde{C} z_h^{\star}(t) + \tilde{D} \lambda_h^{\star}(t) + F^h(t) + C_I x^h(t) \geq 0 \rightarrow 0 \leq \lambda(t) \perp C_{II} z(t) + C_I z(t) +$ 630 $D\lambda(t) + F(t) + C_I x(t) \ge 0$ as $h \to 0$ strongly in $\mathcal{L}^2([0, T]; \mathbb{R}^l)$ where l depends on the 631 variables dimensions. Hence we have proved the existence of a continuous solution 632 633 (x, z, λ) to (2.1) on [0, T]. Consequently the following result is proved: **PROPOSITION 5.2.** Let the conditions of item 2 in Corollary 4.6 hold true uni-634

formly in $h \ (\Rightarrow D \succ 0)$. Assume that Assumptions 10 and 11 hold, $E(\cdot)$ and $F(\cdot)$ are bounded Lipschitz continuous functions, initial data are bounded. Then the DALCS in (2.1) possesses for each admissible bounded initial data, a continuous solution $(x, z, \lambda) : [0, T] \rightarrow \mathbb{R}^{n_1+n_2+m}, T > 0$. Moreover the Euler approximations converge to a solution as $h \rightarrow 0$.

5.1.3. The case blockdiagonal $\tilde{D} = \tilde{D}^{\top} \succeq 0$. This section is dedicated to 640 analyse an alternative to the conditions stated in sections 5.1.1 and 5.1.2. Let us 641 consider the case of Proposition 4.4 with Assumption 6. Then it follows from the 642 complementarity conditions that $\lambda_{1,k+1} = (\tilde{D}_{11} + \partial \psi_{K_1^*})^{-1} (-C_1 z_{k+1} - F_{1,k})$, where 643 the mapping $(\tilde{D}_{11} + \partial \psi_{K^*})^{-1}(\cdot)$ is maximal monotone (being the sum of two maximal 644 645 monotone mappings satisfying [9, Corollary 24.4 (ii)]), single-valued, well-defined and Lipschitz continuous, see Proposition B.1. Thus if $\bar{B}_I = (\bar{B}_I^1, 0)$ with $\bar{B}_I^1 \in \mathbb{R}^{n_1 \times m_1}$ one obtains $\bar{B}_I \lambda_{k+1} = \bar{B}_I^1 \lambda_{1,k+1} = \bar{B}_I^1 (\tilde{D}_{11} + \partial \psi_{K_1^*})^{-1} (-\tilde{C}_1 z_{k+1} - \tilde{F}_{1,k})$ in the right-646 647 hand side of (3.5) (c). 648

649 It may be less restrictive to compute upperbounds using the maximal monotonicity.

To start with let us consider the VI in (4.9) and let us state several features of its solutions. Here Assumption 7 has to be reinforced a little:

ASSUMPTION 12. The conditions of Assumption 7 hold where $X_2 \succ 0$ for any $h \ge 0$.

654 This secures that the variable change $\xi_{k+1} = -\tilde{R}_2^{-1} z_{k+1}$ remains well-defined as $h \to$

655 0. It is also important that item 1 in Proposition 4.4 holds, so that the VI in (4.9) is 656 equivalent to the MLCP in (3.2) (b) (c).

EEMMA 5.3. Assume that $\lambda_{\min}(-\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1})$ (= $\lambda_{\min}(-\tilde{A})$) is large enough, then any solution $\xi_{k+1} = -\tilde{R}_2^{-1} z_{k+1}$ of the VI in (4.9) is upperbounded by an affine function of $||x_k||$.

660 **Proof:** Let us manipulate the VI in (4.9) as follows: (5.10)

$$\begin{aligned} &-R_2 A R_2^{-1} \xi_{k+1} + R_2 E_k + G_2(t_k, \xi_{k+1}) \in -\partial f_2(\xi_{k+1}) \\ \Leftrightarrow & -\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1}(\xi_{k+1} - \xi_0) + \tilde{R}_2 \tilde{E}_k + \tilde{G}_2(t_k, \xi_{k+1}) + \zeta_0 - \tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} \xi_0 = -\zeta_{k+1} + \zeta_0 \\ & \text{with } \zeta_{k+1} \in \partial \tilde{f}_2(\xi_{k+1}), \zeta_0 \in \partial \tilde{f}_2(\xi_0) \\ \Rightarrow & \langle -\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1}(\xi_{k+1} - \xi_0), \xi_{k+1} - \xi_0 \rangle + \langle \tilde{R}_2 \tilde{E}_k, \xi_{k+1} - \xi_0 \rangle \\ & + \langle \tilde{G}_2(t_k, \xi_{k+1}) - \tilde{G}_2(t_k, \xi_0), \xi_{k+1} - \xi_0 \rangle + \langle \zeta_0 - \tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} \xi_0, \xi_{k+1} - \xi_0 \rangle \end{aligned}$$

$$+\langle \hat{G}_2(t_k,\xi_0),\xi_{k+1}-\xi_0\rangle = -\langle \zeta_{k+1}-\zeta_0,\xi_{k+1}-\xi_0\rangle \le 0,$$

where we have denoted $\tilde{G}_2(t_k, \xi_{k+1})$ to highlight the dependency of this mapping on t_k and x_k through \tilde{F}_k . We have $||\tilde{R}_2\tilde{B}_1(\tilde{D}_{11}+\partial\psi_{K_1^*})^{-1}(z_1)-\tilde{R}_2\tilde{B}_1(\tilde{D}_{11}+\partial\psi_{K_1^*})^{-1}(z_2)|| \leq \eta$ $\eta ||z_1-z_2||$ for any z_1 and z_2 , so that $||\tilde{G}_2(t_k,\xi_{k+1})-\tilde{G}_2(t_k,\xi_0)|| \leq \eta ||\tilde{C}_1\tilde{R}_2^{-1}(\xi_{k+1}-\xi_0)|| \leq \eta' ||\xi_{k+1}-\xi_0||$, where η depends on $\tilde{D}_{11}, \tilde{R}_2$ and \tilde{B}_1 (see the proof of Proposition B.1 in [16]), and η' depends on $\tilde{D}_{11}, \tilde{R}_2, \tilde{B}_1$, and \tilde{C}_1 . Thus we obtain: (5.11)

$$(-\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} (\xi_{k+1} - \xi_0), \xi_{k+1} - \xi_0) \leq \langle -\tilde{R}_2 \tilde{E}_k, \xi_{k+1} - \xi_0 \rangle - \langle \tilde{G}_2 (t_k, \xi_0), \xi_{k+1} - \xi_0 \rangle + \langle \zeta_0 + \tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} \xi_0, \xi_{k+1} - \xi_0 \rangle + \eta' ||\xi_{k+1} - \xi_0||^2$$

668 which implies

(5.12)

661

$$(\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1} - \eta' I_{n_{2}})(\xi_{k+1} - \xi_{0}), \xi_{k+1} - \xi_{0}) \leq ||\tilde{R}_{2}\tilde{E}_{k}|| ||\xi_{k+1} - \xi_{0}|| + ||\tilde{G}_{2}(t_{k},\xi_{0}) \pm \tilde{G}_{2}(t_{0},\xi_{0})|| ||\xi_{k+1} - \xi_{0}|| + ||\zeta_{0} + \tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1}\xi_{0}|| ||\xi_{k+1} - \xi_{0}||.$$

670 Now one has: $||\tilde{E}_{k}|| \leq ||E_{II,k} + hA_{II}E_{I,k}|| + ||A_{II}(P_{I}^{h})^{-1}P_{I}|| ||x_{k}||, \tilde{G}_{2}(t_{k},\xi_{0}) =$ $\tilde{R}_{2}\tilde{B}_{1}(\tilde{D}_{11} + \partial\psi_{K_{1}^{*}})^{-1}(\tilde{C}_{1}\tilde{R}_{2}^{-1}\xi_{0} - \tilde{F}_{1,k}), ||\tilde{G}_{2}(t_{k},\xi_{0}) - \tilde{G}_{2}(t_{0},\xi_{0})|| = ||\tilde{R}_{2}\tilde{B}_{1}(\tilde{D}_{11} + \partial\psi_{K_{1}^{*}})^{-1}(\tilde{C}_{1}\tilde{R}_{2}^{-1}\xi_{0} - \tilde{F}_{1,k}) - \tilde{R}_{2}\tilde{B}_{1}(\tilde{D}_{11} + \partial\psi_{K_{1}^{*}})^{-1}(\tilde{C}_{1}\tilde{R}_{2}^{-1}\xi_{0} - \tilde{F}_{1,0})|| \leq \eta ||-\tilde{F}_{1,k} +$ $\tilde{F}_{1,0}|| \leq \eta ||F_{k} + hC_{I}(P_{I}^{h})^{-1}E_{I,k}|| + \eta ||C_{I}(P_{I}^{h})^{-1}P_{I}|| ||x_{k}||.$ Let us denote $\alpha_{k} =$ $||\tilde{R}_{2}|| ||E_{II,k} + hA_{II}E_{I,k}|| + \eta ||F_{k} + hC_{I}(P_{I}^{h})^{-1}E_{I,k}||, \beta_{0} = ||\tilde{G}_{2}(t_{0},\xi_{0})|| + ||\zeta_{0} +$ $\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1}\xi_{0}||, \text{ and } \gamma = \eta ||C_{I}(P_{I}^{h})^{-1}P_{I}|| + ||\tilde{R}_{2}|| ||A_{II}(P_{I}^{h})^{-1}P_{I}||.$ Then: (5.13)

676
$$\langle (-\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} - \eta' I_{\mathbf{n}_2}) (\xi_{k+1} - \xi_0), \xi_{k+1} - \xi_0 \rangle \leq (\alpha_k + \beta_0 + \gamma ||x_k||) ||\xi_{k+1} - \xi_0|| \leq \frac{1}{2} \alpha_k^2 + \frac{1}{2} \beta_0^2 + \frac{1}{2} \gamma^2 ||x_k||^2 + \frac{1}{2} ||\xi_{k+1} - \xi_0||^2$$

677 which implies:

678
$$\langle (-\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} - \eta' I_{n_2} - \frac{1}{2} I_{n_2})(\xi_{k+1} - \xi_0), \xi_{k+1} - \xi_0 \rangle \leq \frac{1}{2} \alpha_k^2 + \frac{1}{2} \beta_0^2 + \frac{1}{2} \gamma^2 ||x_k||^2.$$

679 Therefore if $-\tilde{R}_2\tilde{A}\tilde{R}_2^{-1} - \eta' I_{n_2} - \frac{1}{2}I_{n_2} \succ 0$, one obtains

$$(5.15) \qquad \begin{aligned} \|\xi_{k+1} - \xi_0\| &\leq \frac{1}{\sqrt{\lambda_{\min}(-\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} - \eta' I_{n_2} - \frac{1}{2} I_{n_2})}} \sqrt{\frac{1}{2}} \alpha_k^2 + \frac{1}{2} \beta_0^2 + \frac{1}{2} \gamma^2 \|x_k\|^2} \\ &\leq \frac{1}{\sqrt{2} \sqrt{\lambda_{\min}(-\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} - \eta' I_{n_2} - \frac{1}{2} I_{n_2})}} \left(\sqrt{\alpha_k^2 + \beta_0^2} + \gamma \|x_k\|\right) \end{aligned}$$

681 which ends the proof. \boxtimes

Therefore the term $hB_I z_{k+1}$ in the right-hand side of (3.5) (c), admits also an affine upperbound in $||x_k||$. A further characterization of the iterates is given now.

684 LEMMA 5.4. Assume that $-\tilde{R}_2\tilde{A}\tilde{R}_2^{-1} \succ \frac{3}{2}I_{n_2}$. Assume also that the mapping 685 $\tilde{G}_2(\cdot)$ is monotone (see Lemma 4.5 for conditions), and that $E_I(\cdot)$, $E_{II}(\cdot)$ and $F(\cdot)$ 686 are uniformly bounded and Lipschitz continuous. Then for any k and $p \ge 0$ and any 687 solution of the VI in (4.9), one has:

688 (5.16)
$$\|\xi_k - \xi_p\| \le \alpha + \beta \|x_{k-1} - x_{p-1}\|$$

689 for some bounded constants $\alpha \geq 0$ and $\beta \geq 0$, and α is proportional to $(t_{k-1} - t_{p-1})$.

690 Proof: From

691 (5.17)
$$\begin{cases} -\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} \xi_k + \tilde{R}_2 \tilde{E}_{k-1} + \tilde{G}_2(t_{k-1}, \xi_k) = -\zeta_k \\ -\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} \xi_p + \tilde{R}_2 \tilde{E}_{p-1} + \tilde{G}_2(t_{p-1}, \xi_p) = -\zeta_p \end{cases}$$

692 with $\zeta_k \in \partial f_2(\xi_k), \, \zeta_p \in \partial f_2(\xi_p)$, it follows that

$$(5.18) \qquad \begin{array}{l} \langle -\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1}(\xi_{k}-\xi_{p}),\xi_{k}-\xi_{p}\rangle + \langle \tilde{R}_{2}(\tilde{E}_{k-1}-\tilde{E}_{p-1}),\xi_{k}-\xi_{p}\rangle \\ + \langle \tilde{G}_{2}(t_{k-1},\xi_{k}) - \tilde{G}_{2}(t_{k-1},\xi_{p}) + \tilde{G}_{2}(t_{k-1},\xi_{p}) - \tilde{G}_{2}(t_{p-1},\xi_{p}),\xi_{k}-\xi_{p}\rangle \\ = - \langle \zeta_{k}-\zeta_{p},\xi_{k}-\xi_{p}\rangle \end{array}$$

694 which implies

695 (5.19)
$$\langle -\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1}(\xi_k - \xi_p), \xi_k - \xi_p \rangle + \langle \tilde{R}_2 (\tilde{E}_{k-1} - \tilde{E}_{p-1}), \xi_k - \xi_p \rangle + \langle \tilde{G}_2 (t_{k-1}, \xi_p) - \tilde{G}_2 (t_{p-1}, \xi_p), \xi_k - \xi_p \rangle \le 0.$$

696 One has:

(5.20)

6

$$(5.26) \qquad \langle -\tilde{R}_{2}(\tilde{E}_{k-1} - \tilde{E}_{p-1}), \xi_{k} - \xi_{p} \rangle \\ \leq \langle -\tilde{R}_{2}(E_{II,k-1} - E_{II,p-1}), \xi_{k} - \xi_{p} \rangle - \langle \tilde{R}_{2}(hA_{II}(E_{I,k-1} - E_{I,p-1})), \xi_{k} - \xi_{p} \rangle \\ - \langle A_{II}(P_{I}^{h})^{-1}P_{I}(x_{k-1} - x_{p-1}), \xi_{k} - \xi_{p} \rangle \\ \leq \frac{1}{2} \|\tilde{R}_{2}((E_{II,k-1} - E_{II,p-1}) + hA_{II}(E_{I,k-1} - E_{I,p-1}))\|^{2} + \frac{1}{2} |\xi_{k} - \xi_{p}\|^{2} \\ + \frac{1}{2} \|A_{II}(P_{I}^{h})^{-1}P_{I}\|^{2} \|x_{k-1} - x_{p-1}\|^{2} + \frac{1}{2} \|\xi_{k} - \xi_{p}\|^{2} \\ \leq \|\tilde{R}_{2}\|^{2} \|E_{II,k-1} - E_{II,p-1}\|^{2} + \|\xi_{k} - \xi_{p}\|^{2} + h\|A_{II}\|^{2} \|E_{I,k-1} - E_{I,p-1}\|^{2} \\ + \frac{1}{2} \|A_{II}(P_{I}^{h})^{-1}P_{I}\|^{2} \|x_{k-1} - x_{p-1}\|^{2} \leq \alpha' + \beta' \|x_{k-1} - x_{p-1}\|^{2} + \|\xi_{k} - \xi_{p}\|^{2}, \end{cases}$$

for some $\alpha' \ge 0$ proportional to $(t_{k-1} - t_{p-1})^2$ since $E(\cdot)$ is Lipschitz continuous, $\beta' \ge 0$, and it was used that $x^\top y \le ||x|| ||y|| \le \frac{1}{2} (||x||^2 + ||y||^2)$. Also:

$$\langle -\tilde{G}_{2}(t_{k-1},\xi_{p}) + \tilde{G}_{2}(t_{p-1},\xi_{p}),\xi_{k} - \xi_{p} \rangle$$

$$\leq \|\tilde{G}_{2}(t_{k-1},\xi_{p}) - \tilde{G}_{2}(t_{p-1},\xi_{p})\| \|\xi_{k} - \xi_{p}\|$$

$$\leq \frac{1}{2}k_{\tilde{G}_{2}}\|\tilde{F}_{1,k-1} - \tilde{F}_{p-1}\|^{2} + \frac{1}{2}\|\xi_{k} - \xi_{p}\|^{2}$$

$$\leq k_{\tilde{G}_{2}}\|F_{k-1} - F_{p-1}\|^{2} + hk_{\tilde{G}_{2}}\|C_{I}(P_{I}^{h})^{-1}(E_{I,k-1} - E_{I,p-1})\|^{2}$$

$$+ k_{\tilde{G}_{2}}\|C_{I}(P_{I}^{h})^{-1}P_{I}(x_{k-1} - x_{p-1})\|^{2} + \frac{1}{2}\|\xi_{k} - \xi_{p}\|^{2}$$

$$\leq \alpha'' + \beta''\|x_{k-1} - x_{p-1}\|^{2} + \frac{1}{2}\|\xi_{k} - \xi_{p}\|^{2},$$

for some $\alpha'' \ge 0$ proportional to $(t_{k-1} - t_{p-1})^2$ (since $E(\cdot)$ and $F(\cdot)$ are Lipschitz continuous), and $\beta'' \ge 0$. Hence using (5.19):

$$\|\xi_{k} - \xi_{p}\| \leq \frac{1}{\sqrt{\lambda_{\min}(-\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1} - \frac{3}{2}I_{p})}} \sqrt{\alpha' + \alpha'' + (\beta' + \beta'')} \|x_{k-1} - x_{p-1}\|^{2} \\ \leq \frac{1}{\sqrt{\lambda_{\min}(-\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1} - \frac{3}{2}I_{p})}} (\sqrt{\alpha' + \alpha''} + \sqrt{\beta' + \beta''} \|x_{k-1} - x_{p-1}\|)$$

and the lemma is proved. \boxtimes

705

COROLLARY 5.5. Let h > 0. Assume that the conditions and assumptions of Lemma 5.4 are fulfilled. For all $k \ge 1$ and any solution of the VI in (4.9), one has:

708 (5.23)
$$\|\frac{\xi_{k+1} - \xi_k}{h}\| \le \alpha + \beta \|\frac{x_k - x_{k-1}}{h}\|$$

for some bounded constants $\alpha \geq 0$ and $\beta \geq 0$.

Proof: Let us consider the proof of Lemma 5.4, replacing k by k + 1 and p by k - 1. Then α in (5.16) is proportional to h. Redoing the calculations in (5.20) and (5.22), one obtains (5.23).

Clearly for all h > 0 (5.23) can be written equivalently as: $\|\xi_{k+1} - \xi_k\| \leq h\alpha + \beta \|x_k - x_{k-1}\|$. The bounds calculated for the variable ξ_k also hold for the variable $z_k = -\tilde{R}_2\xi_k$. Let us now analyse λ_{k+1} in (3.5) (a). As seen above, $\lambda_{1,k+1} = (\tilde{D}_{11} + \partial \psi_{K_1^*})^{-1}(-\tilde{C}_1 z_{k+1} - \tilde{F}_{1,k})$ and $\lambda_{2,k+1} \in -\mathcal{N}_{K_2}(\tilde{C}_2 z_{k+1} + \tilde{F}_{2,k})$. If the matrix \bar{B}_I has the above structure (*i.e.*, $\bar{B}_I = (\bar{B}_I^1, 0)$), then the term $\bar{B}_I \lambda_{k+1} = \bar{B}_I^1$ is Lipschitz continuous with respect to the variables z_{k+1} , x_k and F_k . In a more general setting one obtains:

720 (5.24)
$$P_{I}^{h}x_{k+1} \in P_{I}x_{k} + hB_{I}z_{k+1} + h\bar{B}_{I}^{1}(\tilde{D}_{11} + \partial\psi_{K_{1}^{*}})^{-1}(-\tilde{C}_{1}z_{k+1} - \tilde{F}_{1,k}) - h\bar{B}_{I}^{2}\mathcal{N}_{K_{2}}(\tilde{C}_{2}z_{k+1} + \tilde{F}_{2,k}).$$

The difficulty is that the last term in the right-hand side of (5.24) is set-valued in 721 general, and this is intrinsically due to the structure of the implicit scheme in (2.1). 722 723 A solution is to consider the VI in (4.9). Its right-hand side matches with the last term in the right-hand side of (5.24) and it is set-valued. However provided the VI 724 is well-posed, we know that one element of this set is equal to the left-hand side of 725 (4.9). We may assume that the scheme in (5.24) is implemented with this element. 726 This element appears to define a Lipschitz continuous function of ξ_{k+1} (let us recall 727 that the Lipschitz constant of the third term is equal to $\frac{1}{\lambda_{\min}(\tilde{D}_{11})}$, see Proposition 728 B.1). If this element is chosen then the numerical scheme (5.24) becomes: 729

730 (5.25)
$$P_{I}^{h}x_{k+1} = P_{I}x_{k} + hB_{I}z_{k+1} + h\bar{B}_{I}^{1}(\tilde{D}_{11} + \partial\psi_{K_{1}^{*}})^{-1}(\tilde{C}_{1}z_{k+1} - \tilde{F}_{1,k}) + h\bar{B}_{I}^{2}(-\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1}\xi_{k+1} - \tilde{R}_{2}\tilde{E}_{k} + \tilde{G}_{2}(t_{k},\xi_{k+1})).$$

Let us now study (5.25), using Lemmata 5.3, 5.4, Corollary 5.5. It is noteworthy that the constant α in both Lemmata and in Corollary can be made arbitrarily small by increasing $-\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1}$ arbitrarily (in the set of positive definite matrices).

PROPOSITION 5.6. Assume that the conditions of Lemmata 5.3, 5.4 and of Corollary 5.5 are verified, $E_I(\cdot)$, $E_{II}(\cdot)$ and $F(\cdot)$ are bounded Lipschitz continuous functions of time, and Assumption 13 below is satisfied. Then the sequences $\{x^h\}_{h\geq 0}$ and $\{\dot{x}^h\}_{h\geq 0}$ in (5.7) are uniformly bounded in h on [0,T]. Thus $\{x^h\}_{h\geq 0}$ is equi-Lipschitz continuous and there is a subsequence of it that converges uniformly in $\mathcal{C}^0([0,T]; \mathbb{R}^{n_1})$ towards a Lipschitz continuous limit denoted as $x(\cdot)$.

740 **Proof:** In the proof, use is made of positive bounded constants denoted generically as

741 k_i . Consider (5.25). First recall that the operator $\zeta \mapsto (\tilde{D}_{11} + \partial \psi_{K_1^{\star}})^{-1}(\zeta)$ is Lipschitz

742 continuous with Lipschitz constant $\frac{1}{\lambda_{\min}(\tilde{D}_{11})}$. Hence we obtain:

$$(5.26) \begin{array}{l} ||(P_{I}^{h})^{-1}\bar{B}_{I}^{1}(\tilde{D}_{11}+\partial\psi_{K_{1}^{*}})^{-1}(-\tilde{C}_{1}z_{k+1}-\tilde{F}_{1,k}) \\ -(P_{I}^{h})^{-1}\bar{B}_{I}^{1}(\tilde{D}_{11}+\partial\psi_{K_{1}^{*}})^{-1}(-\tilde{C}_{1}z_{p+1}-\tilde{F}_{1,p})|| \\ \leq k_{1}||\tilde{C}_{1}(z_{k+1}-z_{p+1}+\tilde{F}_{1,k}-\tilde{F}_{1,p})|| \leq k_{2}||z_{k+1}-z_{p+1}||+k_{3}||\tilde{F}_{1,k}-\tilde{F}_{1,p}|| \\ \leq k_{4}||\xi_{k+1}-\xi_{p+1}||+k_{4}||F_{1,k}-F_{1,p}||+k_{5}||E_{I,k}-E_{I,p}||+k_{6}||x_{k}-x_{p}|| \\ \leq k_{4}(\alpha+\beta||x_{k}-x_{p}||)+k_{7}|t_{k}-t_{p}|+k_{6}||x_{k}-x_{p}|| \leq k_{8}+k_{9}||x_{k}-x_{p}|| \\ +k_{7}|t_{k}-t_{p}|, \end{array}$$

for any k and $p \ge 0$, where we used Lemma 5.4 and the Lipschitzness property. In a quite similar way one has $||(P_I^h)^{-1}\bar{B}_I^2(\tilde{G}_2(\xi_{k+1}) - \tilde{G}_2(\xi_{p+1}))|| \le k_8 + k_9||x_k - x_p|| + k_7|t_k - t_p|$, where k_7 , k_8 and k_9 depend on the Lipschitz constants of $\tilde{G}_2(\cdot, \cdot)$, $F(\cdot)$, and on $||(P_I^h)^{-1}\bar{B}_I^2||$. One also has $||(P_I^h)^{-1}(B_I + \bar{B}_I^2(\tilde{R}_2\tilde{A}\tilde{R}_2^{-1})(z_{k+1} - z_{p+1})|| \le k_{10}(\alpha + \beta||x_k - x_p||) \le k_{11} + k_{12}||x_k - x_p||$, and $||\bar{B}_I^2\tilde{R}_2(\tilde{E}_k - \tilde{E}_p)|| \le k_{13}|t_k - t_p| + k_{14}||x_k - x_p||$. It is deduced that $||x_{k+1} - x_{p+1}|| \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k - x_p|| + hk_{16}|t_k - t_p| + hk_{17}$, where the term k_{17} stems from $k_8 + k_{11} = \alpha(k_4 + k_{10})$. Here we are led to make the following which is more stringent that Assumption 11:

ASSUMPTION 13. There exists $h_{\text{max}} > 0$ such that $||(P_I^h)^{-1}P_I|| + hk_{15} < 1$ for all 753 $0 < h \le h_{\text{max}}$.

Obviously for h = 0 and the Hölder-induced norm one has $||(P_I^h)^{-1}P_I||_{2,2} + hk_{15} = 1$. 754Let us choose p = 0, it follows that $||x_{k+1} - x_1|| \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k - h|^{-1}P_I|| + hk_{15}) ||x_k - h|^{-1}P_I|| + hk_{15} ||x_k - h|^{-1}P_I|| + hk_{15}$ 755 $|x_0|| + hk_{16}T + hk_{17}$, hence $||x_{k+1}|| \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k|| + hk_{16}T + hk_{17} + hk_{16}T + hk_{17} + hk_{16}T + hk_{16$ 756 $||x_1|| + (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_0||$. Now it has to be proved that $||x_1||$ is bounded. 757 From Lemma 5.3 and (5.15), $||\xi_1||$ and $||z_1||$ are upperbounded by an affine function of 758initial data $||x_0||$ and $||z_0||$. The operator $(D_{11} + \partial \psi_{K_1^{\star}})^{-1}(\cdot)$ is well-defined and single-759 valued. Taking k = 0 in (5.25) one infers that indeed $||x_1||$ is bounded for bounded 760 initial data. Consequently it follows that $||x_{k+1}|| \leq (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k|| + hk_{15}$ 761 $hk_{16}T + hk_{17} + hf(x_0, z_0) + ||(P_I^h)^{-1}P_I|| ||x_0||$ for some function $f(\cdot)$ that is bounded 762for bounded initial data. To simplify the notations let us set $h_1 \stackrel{\Delta}{=} ||(P_I^h)^{-1}P_I|| + hk_{15}$ 763 and $h_2 \stackrel{\Delta}{=} k_{16}T + k_{17} + f(x_0, z_0)$. We obtain for all $k \ge 1$: 764

765 (5.27)
$$\begin{aligned} ||x_k|| &\leq h_1^k ||x_0|| + \sum_{j=0}^{k-1} h_1^j (hh_2 + ||(P_I^h)^{-1} P_I|| ||x_0||) \\ &\leq h_1^k ||x_0|| + (hh_2 + ||(P_I^h)^{-1} P_I|| ||x_0||) \frac{1-h_1^{k+1}}{1-h_1}. \end{aligned}$$

Thus in view of Assumption 13 the sequence $\{x_k\}_{k\geq 0}$ generated by (5.25) is uniformly bounded for bounded initial data.

Let us now reconsider the inequality $||x_{k+1} - x_{p+1}|| \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k - x_p|| + hk_{16}|t_k - t_p| + hk_{17}$ obtained above. Setting p = k - 1 yields $||x_{k+1} - x_k|| \le ||x_k - x_p|| + hk_{17}$ obtained above. 768 769 $(||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k - x_{k-1}|| + h^2k_{16} + hk_{17} \Rightarrow \frac{||x_{k+1} - x_k||}{h} \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k - x_{k-1}|| + h^2k_{16} + hk_{17} \Rightarrow \frac{||x_{k+1} - x_k||}{h} \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k - x_{k-1}|| + h^2k_{16} + hk_{17} \Rightarrow \frac{||x_{k+1} - x_k||}{h} \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k - x_{k-1}|| + h^2k_{16} + hk_{17} \Rightarrow \frac{||x_{k+1} - x_k||}{h} \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k - x_{k-1}|| + h^2k_{16} + hk_{17} \Rightarrow \frac{||x_{k+1} - x_k||}{h} \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k - x_{k-1}|| + h^2k_{16} + hk_{17} \Rightarrow \frac{||x_{k+1} - x_k||}{h} \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k - x_{k-1}|| + h^2k_{16} + hk_{17} \Rightarrow \frac{||x_{k+1} - x_k||}{h} \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k - x_{k-1}|| + h^2k_{16} + hk_{17} \Rightarrow \frac{||x_{k+1} - x_k||}{h} \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k - x_{k-1}|| + h^2k_{16} + hk_{17} \Rightarrow \frac{||x_{k+1} - x_k||}{h} \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k - x_{k-1}|| + h^2k_{16} + hk_{17} \Rightarrow \frac{||x_{k+1} - x_k||}{h} \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k - x_{k-1}|| + h^2k_{16} + hk_{17} \Rightarrow \frac{||x_k - x_k||}{h} \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k - x_{k-1}|| + h^2k_{16} + hk_{17} \Rightarrow \frac{||x_k - x_k||}{h} \le (||x_k - x_k|^{-1}) ||x_k - x_k|^{-1} + hk_{17} \Rightarrow \frac{||x_k - x_k||}{h} \le (||x_k - x_k|^{-1}) ||x_k - x_k|^{-1} + hk_{17} \Rightarrow \frac{||x_k - x_k||}{h} \le (||x_k - x_k|^{-1}) ||x_k - x_k|^{-1} + hk_{17} \Rightarrow \frac{||x_k - x_k||}{h} \le (||x_k - x_k|^{-1}) ||x_k - x_k|^{-1} + hk_{17} \Rightarrow \frac{||x_k - x_k||}{h} \le (||x_k - x_k|^{-1}) ||x_k - x_k|^{-1} + hk_{17} \Rightarrow \frac{||x_k - x_k||}{h} \le (||x_k - x_k|^{-1}) ||x_k - x_k|^{-1} + hk_{17} \Rightarrow \frac{||x_k - x_k||}{h} \le (||x_k - x_k|^{-1}) ||x_k - x_k|^{-1} + hk_{17} \Rightarrow \frac{||x_k - x_k||}{h} \le (||x_k - x_k|^{-1}) ||x_k - x_k|^{-1} + hk_{17} \Rightarrow \frac{||x_k - x_k||}{h} \le (||x_k - x_k|^{-1}) ||x_k - x_k|^{-1} + hk_{17} = hk_{17} + hk_{17} = hk_{17$ 770 $(||(I_I) - I_I|| + hk_{15}) ||x_k - x_{k-1}|| + hk_{16} + hk_{18} + hk_{16} + hk_{18} + hk_{18$ 771 772 773 774 bounded for all $k \ge 1$. Evaluating (5.25) at k = 0 and using the boundedness of x_1 775 as shown above, proves that this is indeed the case. The proposition is proved using 776 the Arzela-Ascoli Theorem. \boxtimes 777

778 Clearly in practice one needs to compute the constants k_i , for instance k_{15} is needed

to determine h_{max} in Assumption 13.

COROLLARY 5.7. Under the conditions of Proposition 5.6, the sequences $\{z^h\}_{h\geq 0}$ and $\{\dot{z}^h\}_{h\geq 0}$ (respectively $\{\xi^h\}_{h\geq 0}$ and $\{\dot{\xi}^h\}_{h\geq 0}$) are uniformly bounded in h on [0,T], $0 < T < +\infty$. Thus $\{z^h\}_{h\geq 0}$ (respectively $\{\xi^h\}_{h\geq 0}$) is equi-Lipschitz continuous and there is a subsequence of it that converges uniformly in $\mathcal{C}^0([0,T];\mathbb{R}^{n_2})$ towards a Lipschitz continuous limit denoted as $z(\cdot)$ (respectively $\xi(\cdot)$).

Proof: This follows from Proposition 5.6 and using Lemmata 5.3, 5.4, Corollary 5.5.
⊠

187 Let us now pass to the analysis of the multiplier convergence. Recall that $\lambda_{1,k+1} = (\tilde{D}_{11} + \partial \psi_{K_1^*})^{-1}(-\tilde{C}_1 z_{k+1} - \tilde{F}_{1,k})$ and $\lambda_{2,k+1} = -\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} \xi_{k+1} + \tilde{R}_2 \tilde{E}_k + \tilde{G}_2(\xi_{k+1}) \in$ $(\tilde{D}_{11} + \partial \psi_{K_1^*})^{-1}(-\tilde{C}_1 z_{k+1} - \tilde{F}_{1,k})$ and $\lambda_{2,k+1} = -\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} \xi_{k+1} + \tilde{R}_2 \tilde{E}_k + \tilde{G}_2(\xi_{k+1}) \in$ $-\mathcal{N}_{K_2}(\tilde{C}_2 z_{k+1} + \tilde{F}_{2,k})$, therefore for all $t \ge 0$, $\lambda_{1,\star}^h(t) = (\tilde{D}_{11} + \partial \psi_{K_1^*})^{-1}(-\tilde{C}_1 z_{\star}^h(t) \tilde{F}_{1,\star}^h(t)$ and $\lambda_{2,\star}^h(t) = -\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} \xi_{\star}^h(t) + \tilde{R}_2 \tilde{E}_{\star}^h(t) + \tilde{G}_2(\xi_{\star}^h(t)) \in -\mathcal{N}_{K_2}(\tilde{C}_2 z_{\star}^h(t) +$ $\tilde{F}_{2,\star}^h(t)$), where we recall that the step-functions are defined as $f_{\star}^h(t) = f_{k+1}$ for $t \in [t_k, t_{k+1})$. We know that $\dot{\xi}_k \triangleq \|\frac{\xi_{k+1}-\xi_k}{h}\| \le C$ for some C > 0, hence: $\|\xi^h - \xi^h\|^2 = \sum_{k=1}^{k+1} |f_{k+1}(t_k) - |f_k|^2 |\dot{\xi}_k|\|^2 dt$

$$(5.28) \qquad \qquad \|\xi^{n} - \xi^{n}_{\star}\|_{\mathcal{L}^{2}([0,T];\mathbb{R}^{n_{2}})}^{2} \qquad \leq \sum_{k=0}^{n-1} \int_{t_{k}}^{t_{k+1}} (t_{k+1} - t_{k})^{2} \|\xi_{k}\|^{2} dt \\ \leq C^{2} \sum_{k=0}^{n-1} (t_{k+1} - t_{k})^{3} = \frac{C^{2} n h^{3}}{3} = \frac{C^{2} T h^{2}}{3}.$$

Consequently $\|\xi^h - \xi^h_\star\|_{\mathcal{L}_2([0,T];\mathbb{R}^{n_2})} \to 0$ as $h \to 0$. Since $\|\xi^h - \xi^h_\star\|_{\mathcal{L}_2([0,T];\mathbb{R}^{n_2})} = \|\xi^h \pm \xi - \xi^h_\star\|_{\mathcal{L}_2([0,T];\mathbb{R}^{n_2})} \to \|\xi - \xi^h_\star\|_{\mathcal{L}_2([0,T];\mathbb{R}^{n_2})} \to 0$, we conclude that the sequence $\{\xi^h_\star\}_{h\geq 0}$ converges strongly in $\mathcal{L}_2([0,T];\mathbb{R}^{n_2})$ towards the limit $\xi(\cdot)$. The same result holds for the sequence $\{x^h_\star\}_{h\geq 0}$ with the limit $x(\cdot)$. Therefore the sequence $\{-\tilde{C}_1 z^h_\star - \tilde{F}^h_{1,\star}\}_{h\geq 0}$ also converges strongly in \mathcal{L}_2 towards a limit, and we infer that the same conclusion holds for both $\{\lambda^h_{1,\star}(\cdot)\}_{h\geq 0}$ and $\{\lambda^h_{2,\star}(\cdot)\}_{h\geq 0}$. Since $\mathcal{N}_{K_2}(\cdot)$ is maximal monotone, it follows from [9, Proposition 20.33] that the limit $\lambda_2(\cdot)$ satisfies $\lambda_2(t) \in -\mathcal{N}_{K_2}(C_{II,2}z(t) + F_2(t) + C_{I,2}x(t))$. Thus we proved the following result.

802 LEMMA 5.8. Let the conditions of Lemmata 5.3, 5.4, Corollary 5.5 be satisfied for 803 all h > 0. The sequences $\{\xi_{\star}^{h}\}_{h \ge 0}$, $\{x_{\star}^{h}\}_{h \ge 0}$, $\{\lambda_{1,\star}^{h}(\cdot)\}_{h \ge 0}$ and $\{\lambda_{2,\star}^{h}(\cdot)\}_{h \ge 0}$ converge 804 strongly in $\mathcal{L}_{2}([0,T]; \mathbb{R}^{n_{2}})$ towards the limits $\xi(\cdot)$, $x(\cdot)$, $\lambda_{1}(\cdot)$ and $\lambda_{2}(\cdot)$, respectively.

To finish this section, one can consider again (5.8) and the analysis which follows it, to obtain a proposition similar to Proposition 5.2, relying on Lemma 5.8, Proposition 5.6, and on Corollary 5.5. It is noteworthy that the three fundamental assumptions in this section, are that $-\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1}$ be sufficiently large as a positive definite matrix, that a suitable selection is chosen for the multiplier $\lambda_{2,k}$, and that a coupling as stated in Assumptions 2 or 7 exists.

811 PROPOSITION 5.9. Assume that the conditions of Proposition 5.6, Corollary 5.7 812 and Lemma 5.8 are verified. Then the limit functions $x(\cdot)$, $z(\cdot)$, $\lambda(\cdot)$, are solutions of 813 the DALCS (2.1).

814 Proof. Let us consider (5.8). From Lemma 5.8, $\{x_{\star}^{h}(\cdot)\}, \{z_{\star}^{h}(\cdot)\}, \{\lambda_{\star}^{h}(\cdot)\}$ converge 815 strongly in $\mathcal{L}^{2}([0,T];\mathbb{R}^{n_{1}})$ towards limits $x(\cdot), z(\cdot)$ and $\lambda(\cdot)$. Using the same argu-816 ments as those developed before Proposition 5.2, it is inferred from (5.8) (b) (c) that 817 the limits verify (2.1) (b) (c). Using the same arguments as those developed just above 818 (5.8), it is deduced that $\{\dot{x}^{h}(\cdot)\}$ converges weakly^{*} to a limit $v(\cdot)$ where $v(\cdot) = \dot{x}(\cdot)$ 819 almost everywhere. Consequently $P_{I}\dot{x}^{h}(t) = A_{I}x_{h}^{*}(t) + B_{I}z_{h}^{*}(t) + \overline{B}_{I}\lambda_{h}^{*}(t) + E_{I}^{h}(t) \longrightarrow$ 820 $P_{I}\dot{x}(t) = A_{I}x(t) + B_{I}z(t) + \overline{B}_{I}\lambda(t) + E_{I}(t)$ weakly in $\mathcal{L}^{2}([0,T];\mathbb{R}^{n_{1}})$ as $h \to 0$. \Box

5.2. Convergence results from (3.1). In this section the properties of the plant's model (2.1) are used, thus the numerical scheme in (3.1) and the results in

section 4.2 are analysed. This section is devoted to the case $D \succeq 0$ as in Proposition 823 4.8. Indeed, since the strict passivity of $(B_{II}, \overline{B}_{II}, C_{II}, D)$ implies Assumption 10 for 824 small enough h > 0, it is unnecessary to repeat the developments in sections 5.1.1 and 825 5.1.2. However, it is not possible in general to determine the implications between 826 Assumption 7 and Assumption 9, nor the relations between $\lambda_{\min}(A)$ and $\lambda_{\min}(B_{II})$ 827 for h > 0. Therefore let us deal with the VI in (4.17) where x_{k+1} is replaced by 828 $x_{k+\theta} = \theta x_k + (1-\theta) x_{k+1}, \ \theta \in [0,1]$. The following holds, which is the counterpart of 829 Lemma 5.3. 830

831 LEMMA 5.10. Assume that $\lambda_{\min}(-R_2B_{II}R_2^{-1}) (= \lambda_{\min}(-B_{II}))$ is large enough, 832 $E_I(\cdot), E_{II}(\cdot)$ and $F(\cdot)$ are uniformly bounded and Lipschitz continuous functions of 833 time. Then any solution $\xi_{k+1} = R_2^{-1} z_{k+1}$ of the VI in (4.17) is upperbounded by an 834 affine function of $||x_{k+\theta}||$.

The proof follows the same steps as the proof of Lemma 5.3 and it is therefore omitted. Lemma 5.10 implies that the term $B_I z_{k+1}$ in the right-hand side of (3.1) (a), is upperbounded by an affine function of $||x_{k+\theta}||$. The coefficients of the affine function depend only on the system's data (constant matrices, Lipschitz constants) and on the initial conditions. Let us now propose the counterpart of Lemma 5.4.

LEMMA 5.11. Assume that $-R_2B_{II}R_2^{-1} \succ I_{n_2}$. Assume also that the well-defined, single-valued and Lipschitz continuous mapping $G_2(\cdot)$ is monotone, and that $E_I(\cdot)$, $E_{II}(\cdot)$ and $F(\cdot)$ are Lipschitz continuous functions of time. Then for any k and $p \ge 0$ and any solution of the VI in (4.17), one has:

844 (5.29)
$$\|\xi_k - \xi_p\| \le \alpha + \beta \|x_{k-1+\theta} - x_{p-1+\theta}\|$$

for all $k \ge 1$ and some bounded constants $\alpha \ge 0$ (proportional to $|t_{k-1} - t_{p-1}|$) and $\beta \ge 0$.

Proof: The proof is similar to the proof of Lemma 5.4, since both VIs in (4.17) and (4.9) possess similar structures. Using (4.17) one finds:

849 (5.30)
$$\begin{cases} -R_2 B_{II} R_2^{-1} \xi_{k+1} - R_2 A_{II} x_{k+\theta} - R_2 E_{II,k} + G_2(t_k, \xi_{k+1}) = -\zeta_{k+1} \\ -R_2 B_{II} R_2^{-1} \xi_{p+1} - R_2 A_{II} x_{p+\theta} - R_2 E_{II,p} + G_2(t_p, \xi_{p+1}) = -\zeta_{p+1} \end{cases}$$

with $\zeta_k \in \partial f_2(\xi_k), k \ge 0$. Using the monotonicity of $G_2(t_k, \cdot)$ and of $\partial f_2(\cdot)$ it follows that:

$$\begin{cases} \langle -R_2 B_{II} R_2^{-1}(\xi_{k+1} - \xi_{p+1}), \xi_{k+1} - \xi_{p+1} \rangle \\ \leq \langle R_2 A_{II}(x_{k+\theta} - x_{p+\theta}) + R_2(E_{II,k} - E_{II,p}), \xi_{k+1} - \xi_{p+1} \rangle \\ - \langle G_2(t_k, \xi_{k+1}) - G_2(t_k, \xi_{p+1}), \xi_{k+1} - \xi_{p+1} \rangle \\ + \langle G_2(t_k, \xi_{p+1}) - G_2(t_p, \xi_{p+1}), \xi_{k+1} - \xi_{p+1} \rangle \\ \leq \|R_2 A_{II}\|^2 \|x_{k+\theta} - x_{p+\theta}\|^2 + \|R_2\|^2 k_E^2 (t_k - t_p)^2 \\ + \frac{1}{2} k_{G_2}^2 \|F_k - F_p\|^2 + \|\xi_{k+1} - \xi_{p+1}\|^2 \\ \leq \|R_2 A_{II}\|^2 \|x_{k+\theta} - x_{p+\theta}\|^2 + \|R_2\|^2 k_E^2 (t_k - t_p)^2 + \frac{1}{2} k_{G_2}^2 k_F^2 (t_k - t_p)^2 \\ + \|\xi_{k+1} - \xi_{p+1}\|^2. \end{cases}$$

853 The lemma is proved. \boxtimes

Thus Corollary 5.5 extends as well, and in particular one has $\|\xi_{k+1} - \xi_k\| \le h\alpha + \beta \|x_{k+\theta} - x_{k-1+\theta}\|$ for all $k \ge 1$, for bounded constants $\alpha \ge 0$ and $\beta \ge 0$. The scheme in (5.25) is now replaced by:

857
$$P_{I}^{h}x_{k+1} = P_{I}x_{k} + hB_{I}z_{k+1} + h\bar{B}_{I}^{1}(D_{11} + \partial\psi_{K_{1}^{\star}})^{-1}(C_{II}^{1}R_{2}^{-1}\xi_{k+1} + F_{1,k}) + h\bar{B}_{I}^{2}(-R_{2}B_{II}R_{2}^{-1}\xi_{k+1} - R_{2}A_{II}x_{k+\theta} - R_{2}E_{II,k} + G_{2}(t_{k},\xi_{k+1})),$$

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858 equivalently:
(5.33)

$$(P_{I}^{h} + h(1-\theta)\bar{B}_{I}^{2}R_{2}A_{II})x_{k+1} = P_{I}x_{k} + hB_{I}z_{k+1} + h\bar{B}_{I}^{1}(D_{11} + \partial\psi_{K_{1}^{\star}})^{-1}(C_{II}^{1}R_{2}^{-1}\xi_{k+1} + F_{1,k}) + h\bar{B}_{I}^{2}(-R_{2}B_{II}R_{2}^{-1}\xi_{k+1} - \theta R_{2}A_{II}x_{k} - R_{2}E_{II,k} + G_{2}(t_{k},\xi_{k+1})),$$

and it appears that both schemes in (5.25) and in (5.33) possess the same structure. Consequently the results of convergence obtained for the foregoing implicit scheme (Proposition 5.6, Corollary 5.7, Lemma 5.8) extend to the scheme in this section, for suitably chosen numerical and system parameters. For the sake of briefness we do not repeat the developments.

Remark 5.12. Using Lemma 4.1 item 1, Assumption 10 could be replaced by assuming the strict passivity of the quadruple $(B_{II}, \overline{B}_{II}, C_{II}, D)$ (a system property) and h > 0 small enough.

Remark 5.13. Obviously when B_{II} is invertible it is possible to recast (2.1) into 868 an LCS by solving the equality constraint for z. This yields the LCS quadruple $(P_I^{-1}(A_I - B_I B_{II}^{-1} A_{II}), P_I^{-1}(\bar{B}_I - B_I B_{II}^{-1} \bar{B}_{II}), C_I - C_{II} B_{II}^{-1} A_{II}, D - C_{II} B_{II}^{-1} \bar{B}_{II}).$ It is unclear how to relate in general the passivity of this quadruple, to the above 869 870 871 assumptions which focus on the MLCP part (2.1) (b) (c) of the system. Especially 872 no specific assumptions are made on the matrices A_I , B_I , \bar{B}_I , C_I in the above analy-873 ses. For instance Assumption 3 implies that the feedthrough matrix of the LCS is 874 $D - C_{II}B_{II}^{-1}X^{-1}C_{II}^{\top}$, which may be indefinite in general. It is also noteworthy that 875 876 transforming the MLCP into an LCP may be undesirable [59].

5.3. Recapitulation. Let us recapitulate the various assumptions upon which 877 the results in this article rely. Assumptions 2, 4, 7, 10, 12 (respectively Assumptions 878 (3, 5, 9), characterize the passivity properties of an auxiliary operator defined by the 879 880 DALCS constraints in (3.2) (b) (c) (respectively in (2.1) (b) (c)). These passivity properties constrain the relative degree between the "input" λ_{k+1} and the "output" 881 $w_{k+1} = \tilde{C}z_{k+1} + \tilde{D}\lambda_{k+1}$ (respectively λ and $w = C_{II}z + D\lambda$). Assumptions 6 and 8 882 have practical importance (because some systems possess such D matrix), they are 883 also related to passivity in the sense that passive systems possess positive semidefinite 884 throughout matrices. Thus all these assumptions are fundamental ones for the devel-885 886 opments in this article, which basically deals with the analysis of two types of generalised equations (GE): $0 \in A(u) + B(u)$ in (4.7) and (4.15), and $0 \in A(u) + B(u) + C(u)$ 887 in (4.9) and (4.17), where the mappings $A(\cdot)$, $B(\cdot)$ and $C(\cdot)$ are maximal monotone 888 (notice in passing that due to time-dependency, these GEs do not fit with those in [28, 889 Chapter 2], so that in particular [28, Theorem 2F.6] cannot be applied in general). 890 Most importantly, the convergence results (Proposition 5.2, Lemmata 5.3 and 5.4, 891 Corollary 5.5, Proposition 5.6, Corollary 5.7, Lemmata 5.8 and 5.10) rely upon the 892 positive definiteness of $-R_2B_{II}R_2^{-1}$ and $-\tilde{R}_2\tilde{A}\tilde{R}_2^{-1}$, which is an index one condition 893 and guarantees strong monotonicity properties, allowing us to characterize the prop-894 erties of the solutions to the GEs. It would be worth relaxing it in future studies. 895 One path could be to use least norm elements [34]. Finally, it is worth noting that 896 many practical systems like circuits with nonlinear resistors and inductances contain 897 nonlinearities [19, section 4.1]. An extension of the above developments, consists of 898 replacing the linear term $A_I x$ in (2.1), by a nonlinear term $A_I(x)$, where $A_I(\cdot)$ is 899 globally Lipschitz continuous, and setting $A_I(x_k)$ in (3.1). The matrices $\tilde{A}, \tilde{B}, \tilde{C}, \tilde{D}$ 900 keep the same structure with P_I^{-1} instead of $(P_I^h)^{-1}$, while both \tilde{E}_k and \tilde{F}_k contain 901

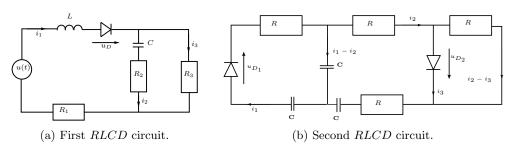


Fig. 1: Two circuits with ideal diodes.

902 $A_I(x_k)$. Nonlinearities in the equality and in the complementarity constraints should 903 also be considered, where the results in [27] could be useful.

904 **6. Examples and Numerical Simulations.** The next examples are taken 905 from circuits, mechanics and state-dependent-switch DAEs. They aim at illustrat-906 ing the above developments, and the limitations imposed by some assumptions.

6.1. Example 1 (circuit). Let us consider the circuit in Figure 1b. Using
Kirchhoff laws its dynamics are given by:

909 (6.1)
$$\begin{cases} \dot{x}_1(t) = \frac{-2}{RC} x_1(t) + \frac{1}{RC} x_2(t) + \frac{1}{R} u_{D_1}(t) \\ \dot{x}_2(t) = \frac{-1}{RC} x_2(t) + \frac{1}{2RC} x_1(t) + \frac{1}{2R} u_{D_2}(t) \\ 0 = \frac{3}{2} u_{D_2}(t) - Rx_3(t) - \frac{1}{C} x_2(t) + \frac{1}{2C} x_1(t) \\ 0 \le u_{D_1}(t) \perp \frac{-2}{RC} x_1(t) + \frac{1}{RC} x_2(t) + \frac{1}{R} u_{D_1}(t) \ge 0 \\ 0 \le u_{D_2}(t) \perp x_3(t) \ge 0, \end{cases}$$

910 with $x_1(t) = \int_0^t i_1(s)ds$, $x_2(t) = \int_0^t i_2(s)ds$, $x_3 = i_3 = z$, $\lambda = (u_{D_1} \quad u_{D_2})^{\top}$ (the 911 voltages across the diodes), $K = K^* = I\!\!R_+^2$, $n_1 = 2$, $p = n_2 = 1$, m = 2. One 912 has $D = \begin{pmatrix} \frac{1}{R} & 0\\ 0 & 0 \end{pmatrix}$, $P_I = I_2$, $A_I = \begin{pmatrix} \frac{-2}{RC} & \frac{1}{RC}\\ \frac{1}{2RC} & \frac{-1}{RC} \end{pmatrix}$, $B_I = \begin{pmatrix} 0\\ 0 \end{pmatrix}$, $\bar{B}_I = \begin{pmatrix} \frac{1}{R} & 0\\ 0 & \frac{1}{2R} \end{pmatrix}$, 913 $A_{II} = (\frac{1}{2C} \quad -\frac{1}{C})$, $B_{II} = -R$, $\bar{B}_{II} = (0 \quad \frac{3}{2})$, $C_I = \begin{pmatrix} \frac{-2}{RC} & \frac{1}{RC}\\ 0 & 0 \end{pmatrix}$, $C_{II} = \begin{pmatrix} 0\\ 1 \end{pmatrix}$. Assumption 8 holds with $D_{11} = \frac{1}{R}$. Thus $\bar{B}_{II}^2 = \frac{3}{2}$, $C_{II}^2 = 1$, $X_2 = \frac{2}{3}$ and Assumption 9 is verified. Item 1 in Proposition 4.8 applies since $\text{Im}(C_{II}^2) = IR$. Since $\bar{B}_{II}^1 = 0$ 914 915the mapping $G_2(\cdot)$ is zero, and item 3 in Proposition 4.8 applies. The quadruple 916 ($B_{II}, \bar{B}_{II}, C_{II}, D$) is passive because the LMI (C.1) is verified with $X = \frac{2}{3}$, hence Assumption 5 is satisfied, but it is not strictly passive since $D \geq 0$. It can be checked that the transfer matrix $H(s) = D + C_{II}(s - B_{II})^{-1}\bar{B}_{II} = \begin{pmatrix} \frac{1}{R} & 0\\ 0 & \frac{3}{2(s+R)} \end{pmatrix}$ is positive 917 918 919 real [18, Theorem 2.48]. The conditions of Lemmata 5.10 and 5.11 are verified by 920 taking R large enough (notice that all other parameters depend on $\frac{1}{R}$). Here A = 921 $B_{II} = -R$, but Assumptions 2 and 6 do not hold, hence the material in Proposition 922

923 4.4 cannot be used.

6.2. Example 2 (circuit). Let us consider now the circuit in Figure 1a. Using Kirchhoff laws its dynamics may be written as:

926 (6.2)
$$\begin{cases} \dot{x}_1(t) = -\frac{R_1}{L}x_1(t) - \frac{R_3}{L}z(t) + \frac{1}{L}u(t) + \frac{1}{L}u_D(t) \\ \dot{x}_2(t) = \frac{R_3}{R_2}z(t) - \frac{1}{R_2\mathbf{C}}x_2(t) \\ 0 = x_1(t) + \frac{1}{R_2\mathbf{C}}x_2(t) - (1 + \frac{R_3}{R_2})z(t) \\ 0 \le u_D(t) \perp x_1(t) \ge 0, \end{cases}$$

927 where $x_1 = i_1, x_2(t) = \int_0^t i_2(s) ds, z = i_3, \lambda = u_D$ (the voltage across the diode), 928 $K = K^* = I\!\!R_+, n_1 = 2, p = n_2 = 1, m = 1$. Let apply the feedback control $u = 2R_3 z$, 929 which modifies B_I . We get: $A_I = \begin{pmatrix} -\frac{R_I}{L} & 0\\ 0 & -\frac{1}{R_2 C} \end{pmatrix}, B_I = \begin{pmatrix} \frac{R_3}{L}\\ \frac{R_3}{R_2} \end{pmatrix}, \bar{B}_I = \begin{pmatrix} \frac{1}{L}\\ 0 \end{pmatrix},$ 930 $A_{II} = (1 \ \frac{1}{R_2 C}), B_{II} = -(1 + \frac{R_3}{R_2}), \bar{B}_{II} = 0, C_I = (1 \ 0), C_{II} = 0, D = 0, E_I = \begin{pmatrix} \frac{1}{L}u\\ 0 \end{pmatrix},$ 921 $E_{II} = 0, E_I = L_2, P^h = \begin{pmatrix} 1 + h\frac{R_1}{L} & 0\\ 0 & - \frac{1}{R_2 C} \end{pmatrix} > 0$ for all h > 0. Thus

931
$$E_{II} = 0, F(t) = 0, P_I = I_2, P_I^h = \begin{pmatrix} 1 + n \frac{1}{L} & 0 \\ 0 & 1 + h \frac{1}{R_2 C} \end{pmatrix} \succ 0 \text{ for all } h \ge 0.$$
 Thus

 $\tilde{B} = \frac{h}{hR_1 + L}, \ \tilde{C} = \frac{hR_3}{L + hR_1}, \ \tilde{D} = \frac{h}{L + hR_1}, \ \tilde{A} = -\frac{R_2 + R_3}{R_2} + hR_3 \left(\frac{1}{L + hR_1} + \frac{1}{R_2(h + R_2\mathbf{C})}\right).$ 932 Assumption 2 holds with $\tilde{X} = R_3$, $\tilde{D} > 0$ for all h > 0, and Assumption 4 holds 933 for h > 0 small enough such that $\tilde{A} < 0$. For h > 0 small enough the set $\tilde{S} = \{0\}$ 934 in (4.6) since $RAR^{-1} = A < 0$, hence the VI in (4.7) always has a solution z_{k+1} by 935 Proposition 4.3 item 3 (a). By item (c) (or (f)) of the same proposition, the solution is 936 unique. When $\tilde{A} < 0$, the quadruple $(\tilde{A}, \tilde{B}, \tilde{C}, \tilde{D})$ is strictly passive for all h > 0, hence 937 items 1 and 3 of Corollary 4.6 apply. However Assumption 10 does not hold since 938 $\dot{D} \to 0$ as $h \to 0$. Consider now the quadruple $(B_{II}, \bar{B}_{II}, C_{II}, D) = (B_{II}, 0, 0, 0)$. 939 It is passive since (C.1) is verified, it is even strictly state passive but not strictly 940 passive since D = 0. Assumptions 3 and 9 are trivially satisfied since both sides of 941 the passivity equality are null. Provided that $\{0\} \in \mathbb{R}_+ - x_{1,k+1} \Leftrightarrow x_{1,k+1} \ge 0$, item 1 942 in Proposition 4.7 applies and the OSNSP for (3.1) has a unique solution (but clearly 943 there is no guarantee that this inequality is satisfied along the integration process, 944 945rendering Proposition 4.7 unapplicable). Proposition 4.8 item 3 also applies here for 946 the VI in (4.17), however one has to check the equivalence stated in item 1 of the same proposition. 947

As noted above, the matrix A_I intervenes only in Assumption 1, hence its modification does not alter the results. Setting $u = 2R_3z + R_1x_1 + aLx_1 + v(t)$ yields the modified

950
$$A_I = \begin{pmatrix} a & 0 \\ 0 & -\frac{1}{R_2 \mathbf{C}} \end{pmatrix}$$
, and B_I as above. This yields (see Remark 5.13): $P_I^{-1}(A_I - \frac{1}{R_2 \mathbf{C}})$

951
$$B_I B_{II}^{-1} A_{II}$$
 = $\begin{pmatrix} a + \frac{R_2 R_3}{L(R_2 + R_3)} & \frac{R_3}{L(R_2 + R_3)} \\ \frac{R_3}{R_2 + R_3} & -\frac{1}{(R_2 + R_3)\mathbf{C}} \end{pmatrix}$, which is unstable if $a + \frac{R_2 R_3}{(R_2 + R_3)L} > 0$

952
$$\overline{(R_2+R_3)\mathbf{C}}$$
. The other elements of the reduced LCS quadruple are given by $P_I^{-1}(B_I - B_I)$
953 $B_I B_{II}^{-1} \bar{B}_{II} = \begin{pmatrix} \frac{1}{L} \\ 0 \end{pmatrix}$, $C_I - C_{II} B_{II}^{-1} A_{II} = (1 \ 0)$, $D - C_{II} B_{II}^{-1} \bar{B}_{II} = 0$. Consequently
954 the reduced LCS is not passive with the above feedback controller. Also $P_I^h = \text{diag}(1 - B_I)$

955
$$ha, 1 + \frac{h}{R_2 C}) > 0$$
 for all $h < \frac{1}{a}$.

To simplify the calculations let us set $R_2 = R_3 = 1\Omega$, L = 1H, $\mathbf{C} = 1$ F, a = 1. Then

957 the discretization (3.1) of (6.2) writes as:

958 (6.3)
$$\begin{cases} x_{1,k+1} - x_{1,k} = hx_{1,k+1} + hz_{k+1} + h\lambda_{k+1} + hv_k \\ x_{2,k+1} - x_{2,k} = -hx_{2,k+1} + hz_{k+1} \\ 0 = x_{1,k+1} + x_{2,k+1} - 2z_{k+1} \\ 0 \le \lambda_{k+1} \perp x_{1,k+1} \ge 0. \end{cases}$$

Equivalently the discretization in (3.2) can be used with the above matrices $\tilde{A}, \tilde{B}, \tilde{C}, \tilde{D}$.

960 **6.3. Example 3 (switching DAEs).** Let us consider the following system, 961 which is a nonsmooth DAE (similar to those studied in [53, 54, 52, 51], and which 962 can be interpreted as a DAE with state-dependent switches):

963 (6.4)
$$\begin{cases} (a) \dot{x}(t) = A_I x(t) + B_I z(t) + E_I(t) \\ (b) 0 = a |x(t)| + b |z(t)| - c z(t) + E_{II}(t), \end{cases}$$

with $x(t) \in \mathbb{R}$, $z(t) \in \mathbb{R}$, $a, b, c \in \mathbb{R}$, c > 0. The constraint in (6.4) (b) is equivalently rewritten as:

966 (6.5)
$$0 = a\lambda_1(t) + a\lambda_2(t) + b\lambda_3(t) + b\lambda_4(t) - cz(t) + E_{II}(t)$$
$$0 \le \begin{pmatrix} \lambda_1(t) \\ \lambda_2(t) \\ \lambda_3(t) \\ \lambda_4(t) \end{pmatrix} \perp \begin{pmatrix} \lambda_1(t) \\ \lambda_2(t) \\ \lambda_3(t) \\ \lambda_4(t) \end{pmatrix} + \begin{pmatrix} -x(t) \\ x(t) \\ 0 \\ 0 \end{pmatrix} + \begin{pmatrix} 0 \\ 0 \\ -z(t) \\ z(t) \end{pmatrix} \ge 0.$$

967 Therefore (6.4) (a) (6.5) fits with (2.1) with $A_{II} = 0$, $B_{II} = -c$, $\bar{B}_{II} = (a \ a \ b \ b)$, 968 $D = I_4$, $C_I = (-1 \ 1 \ 0 \ 0)^{\top}$, $C_{II} = (0 \ 0 \ -1 \ 1)^{\top}$, F(t) = 0, $\bar{B}_I = 0$. Let us check the 969 strict passivity of the quadruple $(B_{II}, \bar{B}_{II}, C_{II}, D)$. This boils down to checking that 970 (see (C.1)):

971 (6.6)
$$M \stackrel{\Delta}{=} \begin{pmatrix} 2cX & -aX & -aX & -bX-1 & -bX+1 \\ -aX & 2 & 0 & 0 & \\ -aX & 0 & 2 & 0 & 0 \\ -bX-1 & 0 & 0 & 2 & 0 \\ -bX+1 & 0 & 0 & 0 & 2 \end{pmatrix} \succ 0,$$

for some X > 0. Let us denote M = P + N with $P = \text{diag}(2cX, 2I_4) \succ 0$, and 972 N_{00} the matrix obtained for a = b = 0. Applying [25, Theorem 2.11] or [10, Fact 973 8.19.2], it follows that $P + N_{00} \succ 0$ if $\sigma_{\max}^2(N_{00}) = 2 < \lambda_{\min}(2cX)\lambda_{\min}(2I_4) = 4cX$. 974Using continuity of the eigenvalues [10, Fact 10.11.8], it is deduced that given c > 0, 975 X > 0, there exists a and b small enough such that (6.6) holds (or, given a, b and 976 X > 0, there exists c > 0 large enough). Thus, item 1 in Lemma 4.1 holds, item 2 977 in Proposition 4.3 holds for h > 0 small enough (guaranteeing $D \succ 0$), Assumption 978 979 10 holds true for small enough h > 0 (see the proof of item 1 in Lemma 4.1), and the arguments stated at the beginning of section 5.2 about the convergence analysis 980led in sections 5.1.1 and 5.1.2, hold. We have: $\tilde{A} = -c$, $\tilde{B} = (a \ a \ b \ b)$, $\tilde{C} = (0 \ 0 \ -1 \ 1)^{\top} + h(-1 \ 1 \ 0 \ 0)^{\top}(1 - hA_I)^{-1}B_I = (\frac{-hB_I}{1 - hA_I} \ \frac{hB_I}{1 - hA_I} \ 0 \ 0)^{\top}$, $\tilde{D} = D$. 981 982Integration can proceed using (3.2). 983

984 Let us now pass to a second example of state-dependent-switch DAE:

985 (6.7)
$$\begin{cases} (a) \dot{x}_1(t) = a_1 x_1(t) + b_1 x_2(t) + c_1 z(t) + e_1(t) \\ (b) \dot{x}_2(t) = a_2 x_1(t) + b_2 x_2(t) + c_2 z(t) + e_2(t) \\ (c) \ 0 = a_3 x_1(t) + b_3 x_2(t) + c_3 z(t) + e_3(t) \text{ if } x_1(t) < 0 \\ (d) \ 0 = a_4 x_1(t) + b_3 x_2(t) + c_3 z(t) + e_3(t) \text{ if } x_1(t) > 0 \end{cases}$$

986 where the switching conditions in (c) and (d) are equivalently rewritten as:

987 (6.8)
$$\begin{cases} 0 = -a_3\lambda_1(t) + a_4\lambda_2(t) + b_3x_2(t) + c_3z(t) + e_3(t) \\ 0 \le \lambda_1(t) \perp \lambda_1(t) + x_1(t) \ge 0 \\ 0 \le \lambda_2(t) \perp \lambda_2(t) - x_1(t) \ge 0. \end{cases}$$

Thus we have: $A_I = \begin{pmatrix} a_1 & b_1 \\ a_2 & b_2 \end{pmatrix}, B_I = (c_1 \ c_2)^{\top}, \bar{B}_I = (0 \ 0)^{\top}, E_I = (e_1 \ e_2)^{\top}, \bar{B}_{II} =$ 988 $(-a_3 \ a_4)^{\top}, A_{II} = (0 \ b_3), B_{II} = c_3, E_{II} = e_3, C_I = \begin{pmatrix} 1 & 0 \\ -1 & 0 \end{pmatrix}, C_{II} = (0 \ 0)^{\top}, D = I_2,$ 989 $F(t) = (0 \ 0)^{\top}$. Also $\tilde{D} = D = I_2$, $\tilde{A} = c_3 + h(0 \ b_3)(I_2 - hA_I)^{-1} \begin{pmatrix} c_1 \\ c_2 \end{pmatrix} = c_3 + \mathcal{O}(h)$, 990 $\tilde{B} = \bar{B}_{II} = (-a_3 \ a_4)^{\top}, \tilde{C} = h \begin{pmatrix} 1 & 0 \\ -1 & 0 \end{pmatrix} (I_2 - hA_I)^{-1} \begin{pmatrix} c_1 \\ c_2 \end{pmatrix} = \mathcal{O}(h). \text{ It can be checked}$ that Proposition 4.3 items 1 and 2, Corollary 4.6 items 1, 2 and 3 (if the passivity LMI holds strictly) apply. The passivity LMI $\begin{pmatrix} -2X\tilde{A} & -X(a_3 \ a_4) + \tilde{C} \\ (-X(a_3 \ a_4) + \tilde{C})^{\top} & 2I_2 \end{pmatrix} \succ 0$ is verified for a_3, a_4 and h > 0 small enough, given $c_3 < 0, X > 0$, implying the 991 992 993 994 strict passivity of the quadruple $(\tilde{A}, \tilde{B}, \tilde{C}, \tilde{D})$, uniformly as $h \to 0$ (Assumption 10) 995 996 holds true). It is of interest to investigate what happens if $c_3 = 0$, which makes the "index" increase. Let us assume for simplicity that $A_I = 0$, the matrix of the passivity LMI becomes $\begin{pmatrix} -2hXb_3c_2 & -X(-a_3 a_4) + h(c_1 - c_1) \\ (-X(-a_3 a_4) + h(c_1 - c_1))^\top & 2I_2 \end{pmatrix}$. For given X > 0 and h > 0, there always exists $b_3c_2 < 0$ such that this matrix is $\succ 0$, 997 998 999 showing that there exists conditions under which the OSNP is still well-posed when 1000 $c_3 = 0$. It is interesting to note that the term $b_3 c_2 z(t)$ appears when an index reduction 1001 is applied to the equality in (6.8) by differentiating it once. Simulations are shown 1002 in Figures 2 and 3. They are performed with the INRIA software package $SICONOS^1$ 1003 [1]. The MLCP is solved with an enumeration method and the differential variables 1004 $x_1(\cdot)$ and $x_2(\cdot)$ are calculated in a second step (hence we follow the philosophy of the 1005 foregoing sections). The parameters are chosen as follows: $a_1 = a_2 = b_1 = b_2 = a_3 =$ 1006 $a_4 = b_3 = c_3 = 1, c_1 = 10, c_2 = 20, e_1(t) = 40\sin(25t), e_2(t) = \sin(t), e_3(t) = \sin(10t),$ 1007 $h = 2.10^{-3}$ s, $x_1(0) = 0$, $x_2(0) = 0.1$. Simulations are done on [0,2]s. The errors are 1008 calculated as $\max_{t\in[0,2]} ||x_i(t) - x_i^{ref}(t)||$, i = 1, 2, with 14 different values of step size h between 10^{-1} s and 10^{-4} s, where $x_i^{ref}(\cdot)$ is calculated with $h = 10^{-4}$ s. It is 1009 1010 observed in Figure 2 that the order is approximately equal to one for small enough 1011 step size. Figure 3 depicts the time evolution of various variables, demonstrating that 1012 the complementarity constraints are respected. Clearly the multipliers are continuous but non differentiable when the complementarity switches from one mode to another. 1014 1015

1016 **6.4. Example 4 (mechanical system).** As an extension of the systems studied 1017 in [37], let us consider a linear mechanical system with both unilateral springs and 1018 bilateral constraints, whose dynamics are:

1019 (6.9)
$$\begin{cases} (a) M\ddot{q}(t) + K_v \dot{q}(t) + K_p q(t) = \tau(t) + H^\top \lambda_u(t) + G^\top \lambda_b(t) \\ (b) 0 = Gq(t) - f(t) \\ (c) 0 \le \lambda_u(t) \perp D\lambda_u(t) + C_I x(t) + F(t) \ge 0, \end{cases}$$

¹https://nonsmooth.gricad-pages.univ-grenoble-alpes.fr/siconos/index.html

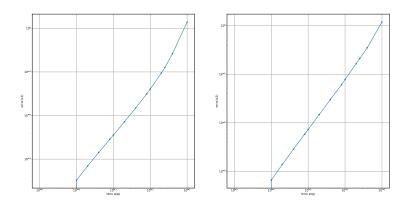


Fig. 2: Log-log error diagram: x_1 (left) and x_2 (right).

1020 where
$$M = M^{\top} \succ 0$$
, $D \succ 0$, $x = (q^{\top} \dot{q}^{\top})^{\top}$, $q(t) \in \mathbb{R}^{\frac{n}{2}}$, $z = \lambda_b \in \mathbb{R}^{n_2}$, $\lambda = \lambda_u \in \mathbb{R}^m$, $G \in \mathbb{R}^{n_2 \times \frac{n_1}{2}}$, $f : \mathbb{R}^+ \to \mathbb{R}^{n_2}$ is twice differentiable, $\tau(\cdot)$ is an exogenous force.
1021 Here an index reduction is possible, in order to replace (6.9) (b) by:

1023 (6.10)
$$0 = G\ddot{q}(t) - \ddot{f}(t) = GM^{-1}H^{\top}\lambda(t) + GM^{-1}G^{\top}z(t) - GM^{-1}(K_p K_v)x(t) - \ddot{f}(t),$$

10241025

and Gq(0) - f(0) = 0, $G\dot{q}(0) - \dot{f}(0) = 0$. Thus $B_{II} = -GM^{-1}G^{\top}$ ($\prec 0$ if bilateral constraints are independent, *i.e.*, rank(G) = $p = n_2 \Rightarrow n_1 \ge 2n_2$), $A_{II} = GM^{-1}(K_p K_v)$, $\bar{B}_{II} = -GM^{-1}H^{\top}$, $B_I = \begin{pmatrix} 0 \\ -M^{-1}G^{\top} \end{pmatrix}$, and $C_{II} = 0$ (similarly to the example in 1026 1027 section 6.2, one major difference being that here $D \succ 0$). Let us examine the OSNSP

for (6.9) on a particular example as depicted in Figure 4, whose dynamics are given 10281029 by:

1030 (6.11)
$$\begin{cases} m\ddot{q}_{1}(t) = \tau(t) + \lambda_{1}(t) - \lambda_{2}(t) \\ m\ddot{q}_{2}(t) = z(t) \\ 0 = q_{2}(t) - f(t) \Leftrightarrow \frac{-z(t)}{m} + \ddot{f}(t) = 0, \ q_{2}(0) - f(0) = 0, \ \dot{q}_{2}(0) - \dot{f}(0) = 0 \\ 0 \le \lambda_{1}(t) \perp \lambda_{1}(t) + k(q_{1}(t) - L - \alpha(t) - l) \ge 0 \\ 0 \le \lambda_{2}(t) \perp \lambda_{2}(t) - k(q_{1}(t) + L - \beta(t) + l) \ge 0. \end{cases}$$

1031 where $q = (q_1 \ q_2)^{\top}$ are the mass coordinates, k > 0 is a stiffness constant, $\alpha(\cdot)$ and $\beta(\cdot)$ are the massless springs' varying positions, l > 0 is the spring's free length, $\beta(t) > \alpha(t) + 2L + 2l$ for all $t \ge 0$. We have: $x = (q_1 \ \dot{q}_1 \ q_2 \ \dot{q}_2)^{\top}$, $G = (0 \ 1)$, H =1032 1033 $\begin{array}{l} 1035 \quad \tilde{p}(t) > \alpha(t) + 2D + 2t \text{ for all } t \geq 0, \text{ we have, } t = (q_1 q_1 q_2 q_2)^{-}, \text{ } t = (0 1), \text{ } t = \\ 1034 \quad \begin{pmatrix} 1 & 0 \\ -1 & 0 \end{pmatrix}, M = \text{diag}(m, m), D = I_2, A_{II} = 0, B_{II} = \frac{-1}{m}, \bar{B}_{II} = 0, E_{II}(t) = \ddot{f}(t), \\ 1035 \quad C_I = \begin{pmatrix} k & 0 & 0 \\ -k & 0 & 0 \end{pmatrix}, F(t) = \begin{pmatrix} k(-L - \alpha(t) - l) \\ -k(L - \beta(t) + l) \end{pmatrix}, C_{II} = 0, B_I = (0 \ 0 \ 0 \ \frac{1}{m})^{\top}, \\ 1036 \quad P_I = I_4 \text{ (and Assumption 1 holds for } h > 0 \text{ small enough}), \bar{B}_I = \begin{pmatrix} 0 & \frac{1}{m} & 0 & 0 \\ 0 & \frac{-1}{m} & 0 & 0 \end{pmatrix}^{\top}, \\ 1037 \quad \tilde{A} = \frac{-1}{m}, \ \tilde{B} = 0, \ \tilde{C} = hC_I(P_I^h)^{-1}B_I, \ \tilde{D} = I_2 + hC_I(P_I^h)^{-1}\bar{B}_I \succ 0 \text{ if } h > 0 \text{ is small enough}. \end{array}$ enough. Let us check Assumption 4. The passivity LMI is $\begin{pmatrix} \frac{2X}{m} & \tilde{C} \\ C^{\top} & \tilde{D} + \tilde{D}^{\top} \end{pmatrix} \succ 0$, 1038

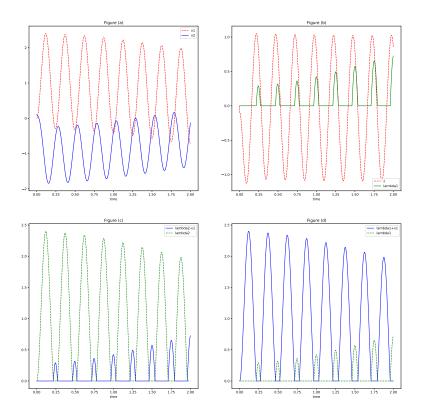


Fig. 3: (a) $x_1(t)$ (dashed) and $x_2(t)$, (b) z(t) (dashed) and $\lambda_1(t)$, (c) $\lambda_2(t) - x_1(t)$ and $\lambda_2(t)$ (dashed), (d) $\lambda_1(t) + x_1(t)$ and $\lambda_1(t)$ (dashed).

X > 0, which is verified for h > 0 small enough and X > 0 large enough. Hence under 1039 these conditions the quadruple $(\hat{A}, \hat{B}, \hat{C}, \hat{D})$ is strictly passive. It can be checked that 1040 items 1 and 2 in Proposition 4.3 hold, as well as items 1, 2 and 3 in Corollary 4.6. 1041 Therefore the dynamics can be integrated using (3.2) with $q_{2,0} + f_0 = 0$, $\dot{q}_{2,0} + \dot{f}_0 = 0$. 1042We note also that the quadruple $(B_{II}, \bar{B}_{II}, C_{II}, D) = (\frac{-1}{m}, 0, 0, I_2)$ is strictly passive 1043 (Assumption 5 is verified), hence item 1 in Lemma 4.1 holds. Assumption 10 is 1044 also verified (for h > 0 small enough but uniformly in $h \ge 0$). Let us analyse the 1045convergence as in section 5.1. Assumption 11 may not be verified in this case, however 1046as noted after this assumption it is easy to modify slightly the numerical scheme so 1047 that it is satisfied with $P_I = I_4$. Hence it is inferred that Proposition 5.2 applies. 1048

1049 **6.5. Example 5 (circuit).** Let us consider the circuit in Figure 5, with $R_1 =$ 1050 $R_3 = 0 \ \Omega$. Let $x_1 = i_2, x_2 = \int_0^t i_3(t) dt, z = i_1, \lambda = u_D$ is the voltage across the

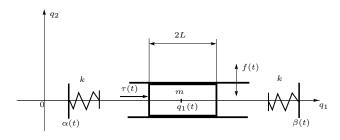


Fig. 4: A mechanical DALCS.

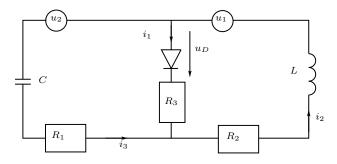


Fig. 5: An *RLCD* passive circuit.

1051 diode, $u_1(\cdot)$ and $u_2(\cdot)$ are two voltage sources. One way to write its dynamics is:

1052 (6.12)
$$\begin{cases} \dot{x}_1(t) = -\frac{R_2}{L}x_1(t) + \frac{1}{L}\lambda(t) + \frac{1}{L}u_1(t) \\ \dot{x}_2(t) = x_1(t) - z(t) \\ 0 = \frac{1}{C}x_2(t) + \lambda(t) + u_2(t) \\ 0 < \lambda(t) \perp z(t) > 0. \end{cases}$$

Assume that the feedback $u_2 = -az + u_3$ is applied, so that the equality constraint 1053 becomes $0 = \frac{1}{C}x_2(t) + u_D(t) - az(t) + u_3(t)$. We have: $A_{II} = (0 \quad \frac{1}{C}), C_I = 0,$ 1054 $C_{II} = 1, B_{II} = -a, \bar{B}_{II} = 1, D = 0, \bar{B}_{I} = (\frac{1}{L} \ 0)^{\top}, B_{I} = (0 \ -1)^{\top}, \text{ and } \tilde{D} = D = 0, \\ \tilde{C} = C_{II} = 1, \tilde{A} = -a + hA_{II}(P_{I}^{h})^{-1}B_{I}, \tilde{B} = 1 + hA_{II}(P_{I}^{h})^{-1}\bar{B}_{I}. \text{ The condition} \\ \operatorname{Im}(C_{II}) \cap (K - C_{I}x_{k+1} + F_{k}) = \mathbb{R} \cap \mathbb{R}_{+} - \frac{1}{\mathbf{C}}x_{2,k+1} \neq \emptyset \text{ is verified always. The set } \mathcal{S} \text{ in} \\ (4.16) \text{ is equal to } \{0\} \text{ hence Proposition } 4.7 \text{ item 1 applies and the VI in } (4.15) \text{ always} \end{cases}$ 10551056 10571058has a solution. Proposition 4.8 items 1, 2, 3 apply, with $D_{11} = 0$ in Assumption 8. 1059The system $(B_{II}, B_{II}, C_{II}) = (-a, 1, 1)$ is passive and Assumptions 3 and 5 hold true. 1060Also the system $(\tilde{A}, \tilde{B}, \tilde{C})$ is passive for h > 0 small enough and Assumptions 2 and 1061 4 hold true. The set S in (4.6) is equal to $\{0\}$ so Proposition 4.3 item 3 (a) applies. 1062For h > 0 small enough, $\tilde{A} < 0$ so Corollary 4.6 items 1 and 2 apply. Assumption 12 1063 is verified for $h \ge 0$ small enough with $\tilde{X} = \tilde{B}^{-1}$. It is inferred that if a > 0 is large 1064enough, the convergence results in sections 5.1.3 and 5.2 apply. 1065

7. Conclusion. This article is devoted to analysing the implicit Euler discretisation of a class of nonsmooth nonlinear systems with equality constraints, named differential-algebraic linear complementarity systems. The well-posedness of the onestep nonsmooth problem is studied first, then the convergence of the discrete-time solutions is analysed. This article only scratches the surface of singular nonsmooth 1071 systems analysis. The above results could be extended in various directions. Most

1072 importantly, it would be worth extending the results to more general indices and rela-

1073 tive degrees, to deeper analyse how the algebraic and the complementarity constraints

are coupled, and to consider state discontinuities and solutions of bounded variation.

Appendix A. Well-posedness of Variational Inequalities. 1075 Some basic notions of Convex Analysis and Complementarity Theory are briefly introduce now 1076[11, 29, 49, 31]. Let $f: \mathbb{R}^n \to \mathbb{R} \cup \{+\infty\}$ be a proper convex and lower semicontin-1077 uous (lsc) function, we denote by dom $(f) \stackrel{\Delta}{=} \{x \in \mathbb{R}^n | f(x) < +\infty\}$ the domain of 1078 the function $f(\cdot)$. The Fenchel transform $f^*(\cdot)$ of $f(\cdot)$ is the proper, convex and lsc 1079 function defined by: for all $z \in \mathbb{R}^n$, $f^*(z) = \sup_{x \in \operatorname{dom}(f)} \{ \langle x, z \rangle - f(x) \}$. The sub-1080 differential $\partial f(x)$ of $f(\cdot)$ at $x \in \mathbb{R}^n$ is defined by: $\partial f(x) = \{\omega \in \mathbb{R}^n | f(v) - f(x) \geq 0\}$ 1081 $\langle \omega, v - x \rangle, \forall v \in \mathbb{R}^n$. We denote by $\text{Dom}(\partial f) \stackrel{\Delta}{=} \{x \in \mathbb{R}^n | \partial f(x) \neq \emptyset\}$ the domain of 1082 the subdifferential operator $\partial f : \mathbb{R}^n \to \mathbb{R}^n$. 1083

Let x_0 be any element in the domain dom(f), the recession function $f_{\infty}(\cdot)$ of $f(\cdot)$ is defined by: for all $x \in \mathbb{R}^n : f_{\infty}(x) = \lim_{\lambda \to +\infty} \frac{1}{\lambda} f(x_0 + \lambda x)$. The function $f_{\infty} : \mathbb{R}^n \to \mathbb{R} \cup \{+\infty\}$ is a proper convex and lsc function which 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 [49]:

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

1084 The set K_{∞} is a nonempty closed convex cone that is described in terms of the 1085 directions which recede from K. The indicator function of a set $K \subseteq \mathbb{R}^n$ is $\psi_K(x) = 0$ 1086 if $x \in K$, $\psi_K(x) = +\infty$ if $x \notin K$. If K is closed non empty convex, we have 1087 $\partial \psi_K(x) = \mathcal{N}_K(x)$, the so-called normal cone to K at x, defined as $\mathcal{N}_K(x) = \{v \in$ 1088 $\mathbb{R}^n \mid v^{\top}(s-x) \leq 0$ for all $s \in K\}$.

- 1089 Important properties of the recession function and recession cone are recalled now:
- 1090 PROPOSITION A.1. [11, Proposition 1.4.8] The following statements hold:

1091 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 1092 lsc functions. Suppose that $f_1 + f_2$ is proper. Then for all $x \in \mathbb{R}^n : (f_1 + f_2)_{\infty}(x) =$ 1093 $(f_1)_{\infty}(x) + (f_2)_{\infty}(x)$.

- 1094 b) Let $f : \mathbb{R}^n \to \mathbb{R} \cup \{+\infty\}$ be a proper, convex and lsc function, and let K be a 1095 nonempty closed convex set, such that $f + \Psi_K$ is proper (equivalently dom $(f) \cap K$ is 1096 non empty). Then for all $x \in \mathbb{R}^n$: $(f + \Psi_K)_{\infty}(x) = f_{\infty}(x) + (\Psi_K)_{\infty}(x)$.
- 1097 c) Let $K \subset \mathbb{R}^n$ be a nonempty, closed and convex set. Then for all $x \in \mathbb{R}^n$: 1098 $(\Psi_K)_{\infty}(x) = \Psi_{K_{\infty}}(x)$. Moreover for all $x \in K$ and $e \in K_{\infty}$: $x + e \in K$.
- 1099 d) If $K \subseteq \mathbb{R}^n$ is a nonempty closed and convex cone, then $K_{\infty} = K$.
- 1100 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 1101 $K_{\infty} = P(A, 0) = \{x \in \mathbb{R}^n | Ax \ge 0\}.$
- 1102 Let us now concatenate [4, Theorem 3, Corollaries 3 and 4]. They concern variational 1103 inequalities (VIs) of the form: Find $u \in \mathbb{R}^n$ such that
- 1104 (A.1) $\langle \mathbf{M}u + \mathbf{q}, v u \rangle + \varphi(v) \varphi(u) \ge 0$, 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 lsc function. The VI in (A.1) is equivalent to the inclusion

1107 $\mathbf{M}u + \mathbf{q} \in -\partial \varphi(u)$, which is a generalized equation (GE) of the form: $0 \in A(u) + B(u)$ 1108 where $A(\cdot)$ and $B(\cdot)$ are two maximal monotone operators (since it will be assumed 1109 $\mathbf{M} \succeq 0$). In fact the problems encountered in this paper (see (4.7), (4.9), (4.14), (4.15), 1110 (4.17)) all belong to this class of GEs. Numerical schemes that compute solutions of 1111 with CEs have have tailed for a laws time in the mathematical literature are [77]

1111 such GEs have been studied for a long time in the mathematical literature, see [57] 1112 and references therein. The next proposition states conditions about existence and

1113 uniqueness of solutions.

1114 The problem in (A.1) is denoted as $VI(\mathbf{M}, \mathbf{q}, \varphi)$, and we set:

1115 (A.2)
$$\mathcal{K}(\mathbf{M},\varphi) = \{ x \in \mathbb{R}^n | \mathbf{M}x \in (\operatorname{dom}(\varphi_\infty))^* \}.$$

1116 Note that $(\operatorname{dom}(\varphi_{\infty}))^*$ is the dual cone of the domain of the recession function φ_{∞} 1117 while $(\operatorname{dom}(\varphi))_{\infty}$ (that we may denote also as $\operatorname{dom}(\varphi)_{\infty}$) is the recession cone of 1118 $\operatorname{dom}(\varphi)$.

1119 PROPOSITION A.2. [4] Let $\varphi : \mathbb{R}^n \to \mathbb{R} \cup \{+\infty\}$ be a proper, convex and lsc 1120 function with closed domain, $\mathbf{M} \in \mathbb{R}^{n \times n}$, and suppose that $\mathbf{M} \succeq 0$.

1121 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 1122 $VI(\mathbf{M}, \mathbf{q}, \varphi)$ has at least one solution.

1123 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 $x_0 \in$ 1124 $\operatorname{dom}(\varphi)$ such that

(A.3)
1125
$$\langle \mathbf{q} - \mathbf{M}^{\top} x_0, v \rangle + \varphi_{\infty}(v) > 0, \ \forall v \in (dom(\varphi))_{\infty} \cap \operatorname{Ker} \{ \mathbf{M} + \mathbf{M}^{\top} \} \cap \mathcal{K}(\mathbf{M}, \varphi), \ v \neq 0,$$

1126 then problem $VI(\mathbf{M}, \mathbf{q}, \varphi)$ has at least one solution.

1127 c) If u_1 and u_2 are two solutions of problem $VI(\mathbf{M}, \mathbf{q}, \varphi)$ then $u_1 - u_2 \in \text{Ker}\{\mathbf{M} + 1128 \ \mathbf{M}^{\mathsf{T}}\}$.

1129 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 1130 $\langle \mathbf{q}, u_1 - u_2 \rangle = \varphi(u_2) - \varphi(u_1)$.

1131 e) 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 1132 the optimization problem $\min_{x \in \mathbb{R}^n} \frac{1}{2}x^{\top}\mathbf{M}x + \langle \mathbf{q}, x \rangle + \varphi(x)$.

1133 Appendix B. Some Convex Analysis and Complementarity Theory 1134 tools. If $K \subseteq \mathbb{R}^n$ is a set, then $K^* = \{z \in \mathbb{R}^n | \langle z, x \rangle \ge 0 \text{ for all } x \in K\}$ is 1135 its dual cone. Let K be a nonempty closed convex cone, then:

1136 (B.1)
$$K^* \ni x \perp y \in K \iff x \in -\mathcal{N}_K(y) \iff y \in -\mathcal{N}_{K^*}(x).$$

1137 Let $M = M^{\top} \succ 0$, x and y two vectors, then [29]

1138 (B.2)
$$M(x-y) \in -\mathcal{N}_K(x) \Leftrightarrow x = \operatorname{proj}_M[K;y] \Leftrightarrow x = \min_{z \in K} \frac{1}{2} (z-y)^\top M(z-y).$$

1139 We note that this is a particular case of (A.1), so that Proposition A.2 can be con-1140 sidered as the characterization of a generalized projection operator $VI(\mathbf{M}, \mathbf{q}, \varphi)$. The 1141 next proposition is used several times in the article:

1142 PROPOSITION B.1. Suppose that the $m \times m$ matrix $D \succ 0$, and $\mathcal{M} : \mathbb{R}^m \rightrightarrows \mathbb{R}^m$ 1143 is a maximal monotone operator. Then the operator $(D + \mathcal{M})^{-1}(\cdot)$ is well-defined 1144 (i.e., for each $y \in \mathbb{R}^m$ there is an $x \in \mathbb{R}^m$ such that $y \in Dx + \mathcal{M}(x)$), single-valued 1145 and Lipschitz continuous with modulus $L = \frac{2}{\lambda_{\min}(D + D^+)}$.

It was proposed in [16, Proposition 1] when $\mathcal{M} = \partial \phi$ for some convex proper lsc $\phi(\cdot)$, 1146 extended to maximal monotone $\mathcal{M}(\cdot)$ in [5]. 1147

1148 **Appendix C. Passive Systems.** A quadruple (A, B, C, D) is said passive if:

1149 (C.1)
$$\begin{pmatrix} -XA - A^{\top}X & -XB + C^{\top} \\ -B^{\top}X + C & D + D^{\top} \end{pmatrix} \succeq 0, \quad X = X^{\top} \succ 0.$$

for some X. It is said *strictly* passive if the first matrix inequality is satisfied with $\succ 0$ 1150 (this implies that $D \succ 0$). This is not to be confused with the *strict state* passivity [18, 11511152Definition 4.54 (strict passivity implies strict state passivity). In fact strict passivity is directly related to strongly SPR transfer matricess [18, p.62 and section 3.12.2], 11531154while SPR transfer matrices with minimal realization are strictly state passive [18, Theorem 4.73 [43]. Strict state passivity holds when $-XA - A^{\top}X \succ 0$ in (C.1), with 1155 $D + D^{\top} \geq 0$. The two inequalities in (C.1) make a linear matrix inequality (LMI) with 1156 unknown X, equivalently a Riccati inequality using the Schur complement Theorem 1157[18, Theorem A.65]. Numerical algorithms exist that efficiently solve LMIs, see [13], 1158

and can be used to test whether or not a quadruple is passive. 1159

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