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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 well-posed 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.

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

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

1. Introduction. The analysis of non-smooth dynamical systems with multivalued right-hand sides satisfying maximal monotone properties has witnessed a large number of contributions [19]. Within this class one finds linear and nonlinear complementarity dynamical systems, projected systems, differential variational inequalities, differential inclusions with maximal monotone right-hand sides, Moreau's sweeping process, some switching dynamics, etc. In this article we focus on singular, or differential-algebraic linear complementarity systems (DALCS), which may be viewed either as an extension of classical differential-algebraic equations (DAE), or of "classical" linear complementarity systems. As will become clear later, DALCS represent an extension of the set-valued Lur'e systems studied, e.g., in [39, 40]. The associated canonical form is motivated by applications in chemistry [53, 52, 62], switching DAE analysis [48], circuits with nonsmooth electronic components [2, section 3.5, chapter 4] [22], optimization-constrained differential equations [19, Example 2] with applications in 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 do not apply to all of these systems).

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]. 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 DAE 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 relying on a special Weierstrass' representation. However the time-discretisation is not tackled

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tools.

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42 in the above references.

43 Our goal is to analyse the well-posedness of the one-step nonsmooth problem (OS-NSP) obtained after an implict (backward) Euler discretization, and the convergence 44 of the discrete-time solutions towards a continuous-time limit. The second step is useful not only for proving the existence of solutions to the continuous-time system, but 46 also, perhaps most importantly, to show that the proposed time-discretization makes 47 sense for numerical simulations. Passivity properties of suitable operators are pivotal 48 throughout the article. It is noteworthy that implicit Euler methods are well-known 49 in DAE [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 53 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 their application to DALCS has not been studied yet. The article is organised as follows: Section 2 introduces the general dynamics of 56 the nonsmooth singular systems, section 3 presents their time-discretization, section 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, while section 5.2 uses properties of the system. Section 6 is dedicated to illustrating 60 examples and simulations, from circuits, mechanics and switching DAE. Conclusions 61 are drawn in section 7. The appendix is dedicated to recall various mathematical 62

Notation and definitions: for any vector $x \in \mathbb{R}^n$ and any matrix $M \in \mathbb{R}^{m \times n}$. 64 ||x|| is the Euclidean norm and ||M|| is the Frobenius norm, which are compatible 65 norms [10, Proposition 9.3.5], i.e., $||Mx|| \le ||M|| \, ||x||$. Let $M \in \mathbb{R}^{n \times m}$, then Im(M)66 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 67 definite matrix: $M \succ 0$ if $x^{\top}Mx > 0$ for all $x \neq 0$, positive semidefinite matrix: 68 $M \geq 0$ if $x^{\top} M x \geq 0$ for all x (such M is not necessarily symmetric). The maximum 69 singular value is denoted as $\sigma_{\max}(M)$, and the minimum eigenvalue as $\lambda_{\min}(M)$. The 70 ith row of M is denoted as $M_{i\bullet}$. The $n \times n$ identity matrix is denoted I_n . A set-valued 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 72 has $\langle x_1 - x_2, y_1 - y_2 \rangle \geq 0$. It is maximal monotone if its graph cannot be enlarged 73 without destroying monotonicity. It is η -strongly monotone if there exists $\eta > 0$ such 74 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. 75

2. The class of nonsmooth singular dynamical systems. Let us consider the following differential algebraic linear complementarity system (DALCS):

(2.1)
$$\begin{cases} (a) \ P_{I}\dot{x}(t) = A_{I}x(t) + B_{I}z(t) + \bar{B}_{I}\lambda(t) + E_{I}(t) \\ (b) \ 0 = A_{II}x(t) + B_{II}z(t) + \bar{B}_{II}\lambda(t) + E_{II}(t) \\ (c) \ K^{\star} \ni \lambda(t) \perp w(t) = C_{I}x(t) + C_{II}z(t) + D\lambda(t) + F(t) \in K, \end{cases}$$

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) 86 is equivalently rewritten as $\phi(\lambda, w) = 0$ where $\phi(\cdot)$ is a C-function [3, 29]. Therefore 87 (2.1) (b) (c) can be rewritten equivalently as $\Phi(x, z, \lambda, t) = 0$ for some function $\Phi(\cdot)$, 88 usually nondifferentiable. The system (2.1) thus possesses the form of a semi-explicit DAE as studied in [53, 52], assuming piecewise differentiable functions $\Phi(\cdot)$, with x 90 the differential state, z and λ the algebraic state variables. The so-called regularity 91 conditions (which are an extension of index 1 conditions for DAE) used in [53, 54, 52, 51] may hold in particular cases for (2.1). Examples in section 6.3 illustrate the links 93 between both approaches in [53, 54, 52, 51] and in this article.

- Consider (2.1), then the initial data $x(0) = x_0$ and $z(0) = z_0$ are compatible with the 95 system's constraints if and only if: 96
- $0 \in A_{II}x_0 + B_{II}z_0 + E_{II}(0) \bar{B}_{II}(D + \mathcal{N}_{K^*})^{-1}(C_Ix_0 + C_{II}z_0 + F(0))$ $\lambda(0) \in -(D + \mathcal{N}_{K^*})^{-1}(C_Ix_0 + C_{II}z_0 + F(0)).$ (2.2)97
- This implies that $C_I x_0 + C_{II} z_0 + F(0) \in \operatorname{Im}(D + \mathcal{N}_{K^*})$ and $A_{II} x_0 + B_{II} z_0 + E_{II}(0) \in$ 98 $\bar{B}_{II}\mathrm{Dom}(D+\mathcal{N}_{K^{\star}}).$ 99
- 3. Time-discretization of DALCS with implicit Euler methods. Let us 100 consider the time interval $[0,T], T>0, h=\frac{T}{N}, N\in \mathbb{N}_+, t_k=kh, k\in\{0,1,\ldots,N\},$ 101 $t_0 = 0, t_N = T$. For a continuous function $f(\cdot), f_k \stackrel{\Delta}{=} f(t_k)$. The implicit Euler 102 discretization of (2.1) reads as: 103

$$\begin{cases}
(a) \ P_{I}(x_{k+1} - x_{k}) = hA_{I}x_{k+1} + hB_{I}z_{k+1} + h\bar{B}_{I}\lambda_{k+1} + hE_{I,k} \\
(b) \ A_{II}x_{k+1} + B_{II}z_{k+1} + \bar{B}_{II}\lambda_{k+1} + E_{II,k} = 0 \\
(c) \ K^{\star} \ni \lambda_{k+1} \perp C_{I}x_{k+1} + C_{II}z_{k+1} + D\lambda_{k+1} + F_{k} \in K.
\end{cases}$$

- The problem in (3.1) (b) (c) is a mixed LCP (MLCP). Let us now state a basic 105 assumption.
- Assumption 1. The matrix $P_I^h \stackrel{\triangle}{=} P_I hA_I$ is full-rank. 107
- In case $P_I > 0$, [25, Theorem 2.11] can be used to set conditions on hA_I so that 108 $P_I^h > 0$: there always exists h_{max} such that the assumption holds for all $h \in [0, h_{\text{max}})$ 109 (with $h_{\text{max}} = +\infty$ if $A_I \leq 0$). Let us define the following matrices and vectors: 110
- 1. $\tilde{A} = hA_{II}(P_I^h)^{-1}B_I + B_{II},$ 111
- 112
- 1. $A = hH_{II}(I_I)$ $E_I + E_{II}$, 2. $\tilde{B} = hA_{II}(P_I^h)^{-1}\bar{B}_I + \bar{B}_{II}$, 3. $\tilde{E}_k = E_{II,k} + hA_{II}(P_I^h)^{-1}E_{I,k} + A_{II}(P_I^h)^{-1}P_Ix_k$, 113
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- 4. $\tilde{C} = hC_I(P_I^h)^{-1}B_I + C_{II},$ 5. $\tilde{D} = hC_I(P_I^h)^{-1}\bar{B}_I + D,$ 6. $\tilde{F}_k = F_k + hC_I(P_I^h)^{-1}E_{I,k} + C_I(P_I^h)^{-1}P_Ix_k.$
- Then we deduce from (3.1) the difference equations: 117

118 (3.2)
$$\begin{cases} (a) \ x_{k+1} = (P_I^h)^{-1} (P_I x_k + h B_I z_{k+1} + h \bar{B}_I \lambda_{k+1} + h E_{I,k}) \\ (b) \ \tilde{A} z_{k+1} + \tilde{B} \lambda_{k+1} + \tilde{E}_k = 0 \\ (c) \ K^* \ni \lambda_{k+1} \perp \tilde{C} z_{k+1} + \tilde{D} \lambda_{k+1} + \tilde{F}_k \in K, \end{cases}$$

where the last two lines of (3.2) are a MLCP with unknowns z_{k+1} and λ_{k+1} . It is 119 therefore possible to study the well-posedness of this MLCP independently of the first line in (3.2). Letting $\zeta_{k+1} \stackrel{\Delta}{=} -z_{k+1}$, this MLCP can be rewritten equivalently as

follows: 122

$$\begin{cases}
\lambda_{k+1} \in (\mathcal{N}_{K^*} + \tilde{D})^{-1}(-\tilde{C}z_{k+1} - \tilde{F}_k), \\
-\tilde{A}\zeta_{k+1} + \tilde{E}_k \in -\tilde{B} (\mathcal{N}_{K^*} + \tilde{D})^{-1}(\tilde{C}\zeta_{k+1} - \tilde{F}_k), \\
x_{k+1} = (P_I^h)^{-1}(P_I x_k + h B_I z_{k+1} + h B_I \lambda_{k+1} + h E_{I,k}).
\end{cases}$$

Also the following calculations stem from the complementarity conditions in (3.2): 124

$$\lambda_{k+1} \in -\mathcal{N}_{K}(\tilde{C}z_{k+1} + \tilde{D}\lambda_{k+1} + \tilde{F}_{k}) \quad \Leftrightarrow -\lambda_{k+1} \in \mathcal{N}_{\Gamma_{k}}(\tilde{C}z_{k+1} + \tilde{D}\lambda_{k+1})$$

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

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

- where the inversion of set-valued mappings is used in the second equivalence, and 126
- where $\sigma_{\Gamma_k}(\cdot)$ is the support function of the set $\Gamma_k(t_k, x_k) \stackrel{\Delta}{=} K \tilde{F}_k(t_k, x_k)$ (which 127
- is a closed convex set, but not a cone as long as $\tilde{F}_k \neq 0$ which is the generic case). 128
- Therefore equivalently we can rewrite (3.3) as: 129

130 (3.5)
$$\begin{cases} (a) \ \lambda_{k+1} \in -(\partial \sigma_{\Gamma_k} + \tilde{D})^{-1}(\tilde{C}z_{k+1}) \\ (b) \ -\tilde{A}z_{k+1} + \tilde{B}(\partial \sigma_{\Gamma_k} + \tilde{D})^{-1}(\tilde{C}z_{k+1}) \ni \tilde{E}_k \\ (c) \ x_{k+1} = (P_I^h)^{-1}(P_I x_k + h B_I z_{k+1} + h \bar{B}_I \lambda_{k+1} + h E_{I,k}). \end{cases}$$

- The second line in (3.3) is a generalized equation (GE) with unknown ζ_{k+1} . Its well-131
- posedness depends on the properties of the operator $\mathcal{F}: \zeta \mapsto -A\zeta + B$ (\mathcal{N}_{K^*}) 132
- \tilde{D})⁻¹($\tilde{C}\zeta \tilde{F}_k$). Similarly for (3.5) with the operator $\mathcal{G}: \zeta \mapsto -\tilde{A}\zeta + \tilde{B}(\partial\sigma_{\Gamma_k} + \sigma_{\kappa_k})$ 133
- $(C\zeta)$, and (3.5) (b) is: $\mathcal{G}(z_{k+1}) \ni E_k$. Various results have been published to 134
- characterize such operators [16, 23, 6] and to guarantee that they are single-valued 135
- Lipschitz continuous [16, Propositions 1, 2, 3, Corollary 1], or maximal monotone 136
- (possibly set-valued) [23, 6]. The difficulty in the analysis of DALCS is that they 137
- involve the composition of two operators as seen in the OSNSP for advancing the 138
- scheme from step k to step k+1: 139

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

- Clearly the scheme can be advanced with a unique solution if both λ_{k+1} and z_{k+1} 141
- are uniquely defined (the "if and only if" holds with $\bar{B}_I \lambda_{k+1}$ and $B_I z_{k+1}$ uniquely 142
- defined). Thus next we are going to characterize the GEs in (3.3) and (3.5) in order 143
- to analyse the scheme (3.6). 144
- Two main classes of results are provided in the sequel of the article: first the well-145
- posedness (existence and uniqueness of solutions) for the OSNSP associated with the 146
- time-discretizations in (3.1) (Propositions 4.3 and 4.4, Corollary 4.6) and (3.2) (Propo-147
- sitions 4.7 and 4.8); second, the convergence of the discrete-time solutions is analysed, 148
- and the limit functions are shown to be solutions of the DALCS (Propositions 5.2 and 149
- 150 5.9).
- 4. OSNSP well-posedness. The approach chosen in this article consists of 151 considering the system as an ODE (2.1) (a) with a nonsmooth constraint (2.1) (b) (c): either (3.1) (b) (c) that is a mixed LCP with unknowns z_{k+1} , λ_{k+1} and parameter 153 x_{k+1} (section 4.2), or (3.2) (b) (c) (equivalently (3.5) (a) (b)) that is a mixed LCP 154 with z_{k+1} , λ_{k+1} (section 4.1). As will become clear later, the results rely heavily on 155 the positive (semi) definiteness of B_{II} , and on specific passivity couplings between 156 z in the complementarity constraint (matrix C_{II}) and λ in the algebraic constraint 157 (matrix \bar{B}_{II}). When $\hat{D} = 0$ in (3.2) (c), let us state the following:

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

The alternative option when D = 0 in (3.1) (c):

ASSUMPTION 3. There exists a $p \times p$ matrix $X = X^{\top} \succ 0$ such that $X\bar{B}_{II} = C_{II}^{\top}$.

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Let us notice that Assumption 2 implies that $\tilde{X}\bar{B}_{II} = C_{II}^{\top} + \mathcal{O}(h)$, while Assumption 163 3 implies that $X\tilde{B} = \tilde{C}^{\top} + \mathcal{O}(h)$. The influence of the $\mathcal{O}(h)$ terms may not be easy 164 to characterize precisely for finite h > 0. Assumption 3 is a model assumption, while 165 Assumption 2 is a property of the numerical method. One sees that Assumption 3 166 corresponds to a passivity constraint on the triple $(B_{II}, \bar{B}_{II}, C_{II})$ [18]. Assumption 167 2 imposes an input/output passivity-like constraint on the triple $(\tilde{A}, \tilde{B}, \tilde{C})$, which is 168 an $\mathcal{O}(h)$ approximation of $(B_{II}, \bar{B}_{II}, C_{II})$. When $D \neq 0$ and $D \neq 0$, extensions of 169 Assumptions 2 and 3 are as follows: 170

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

and the counterpart of Assumption 3:

ASSUMPTION 5. The quadruple $(B_{II}, \bar{B}_{II}, C_{II}, D)$ 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.

Lemma 4.1. The following statements are true:

- 1. Let $(B_{II}, \bar{B}_{II}, C_{II}, D)$ be strictly passive. Then $(\tilde{A}, \tilde{B}, \tilde{C}, \tilde{D})$ is strictly passive for small enough h > 0.
- 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 state passive if: a) h > 0 is small enough, $b) D + D^{\top} = 0$ and $C_I(P_I^h)^{-1}\bar{B}_I + (C_I(P_I^h)^{-1}\bar{B}_I)^{\top} \geq 0$, $c) XA_{II}(P_I^h)^{-1}\bar{B}_I B_I^{\top}(P_I^h)^{-\top}C_I^{\top} = 0$, where X is a solution of the passivity LMI for $(B_{II}, \bar{B}_{II}, C_{II}, D)$.
- 3. Let h > 0 and $(B_{II}, \bar{B}_{II}, C_{II}, D)$ be passive. Then $(\tilde{A}, \tilde{B}, \tilde{C}, \tilde{D})$ is passive if: b) and c) in item 2 hold, and a) $XA_{II}(P_I^h)^{-1}B_I + B_I^\top(P_I^h)^{-\top}A_{II}^\top X \leq 0$, where X is a solution of the passivity LMI for $(B_{II}, \bar{B}_{II}, C_{II}, D)$.

Proofs

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

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

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

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

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

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

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The converse implication in item 1 does not hold, since one may have $\tilde{M} \succ 0$ for some 210 h>0, but $M\neq 0$. It is possible to state necessary and sufficient conditions, relying 211 on [10, Proposition 8.2.4, Fact 8.21.22] or [18, Lemma A.70]. However the obtained 212 213 conditions are cumbersome and hardly usable. Lemma 4.1 can be used to determine the conditions under which Assumption 4 is verified in the results presented in the next 214 section. The conditions b) and c) in Lemma 4.1 item 2 hold if: b) $C_I = 0$ or $\bar{B}_I = 0$, c) $(A_{II} = 0 \text{ or } \bar{B}_I = 0) \text{ and } (B_I = 0 \text{ or } C_I = 0).$ Condition a) in Lemma 4.1 item 3 holds 216 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 $X = (-\frac{1}{2}Q + F)B_{II}^{-1}$, where $F = -F^{\top}$ is unique [60, Theorems 1 and 2]. Let us denote the symmetric 218 part of a square matrix M as sym(M) and its skew-symmetric part as skw(M). This 220 221 allows us to state the following necessary conditions.

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

226 (4.3)
$$\operatorname{sym}(\frac{1}{2}QB_{II}^{-1}A_{II}(P_I^h)^{-1}\bar{B}_I + B_I^{\top}(P_I^h)^{-\top}C_I^{\top}) = F\operatorname{skw}(B_{II}^{-1}A_{II}(P_I^h)^{-1}\bar{B}_I)F,$$

227 for any pair (Q, F), and

228 (4.4)
$$\operatorname{Ker}(C_I(P_I^h)^{-1}B_I) = \{0\}.$$

229 2) Let $(B_{II}, \bar{B}_{II}, C_{II}, D)$ be passive with $XB_{II} + B_{II}^{\top}X = -Q$, $Q = Q^{\top} \geq 0$ and B_{II} 230 invertible. Then $X = (-\frac{1}{2}Q + F)B_{II}^{-1}$, for some $F = -F^{\top}$, and condition a) item 3 231 in Lemma 4.1 holds only if

232 (4.5)
$$\operatorname{sym}(-QB_{II}^{-1}A_{II}(P_I^h)^{-1}B_I) + 2F\operatorname{skw}(B_{II}^{-1}A_{II}(P_I^h)^{-1}B_I)F \leq 0,$$

233 for any pair (Q, F).

Proof: 1) Condition c) yields $(-\frac{1}{2}Q + F)B_{II}^{-1}A_{II}(P_I^h)^{-1}\bar{B}_I = B_I^{\top}(P_I^h)^{-\top}C_I^{\top}$, 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 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.

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

- 1. Let Assumption 4 hold. Suppose that either $\operatorname{Im}(\tilde{C}) \cap \operatorname{int}(\operatorname{Im}(\partial \sigma_{\Gamma_k} + \tilde{D})) \neq \emptyset$ or \tilde{B} is square full rank. Suppose also that there is ζ_0 such that $\tilde{C}\zeta_0 \in \operatorname{Im}(\partial \sigma_{\Gamma_k} + \tilde{D})$. 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 $h \geq 0$, $k \geq 0$ (maximal monotonicity holds for $\tilde{D} \succcurlyeq 0$). 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 $\tilde{R}\tilde{A}\tilde{R}^{-1} \leq 0$. Let us consider the set:

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

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 $VI(-\tilde{R}\tilde{A}\tilde{R}^{-1}, \tilde{q}_k, \tilde{f})$: Find $\xi_{k+1} = \tilde{R}z_{k+1}$ such that

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

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

- (c) If $\hat{\xi}_{k+1}^1$ and $\hat{\xi}_{k+1}^2$ are two solutions of the $VI(-\tilde{R}\tilde{A}\tilde{R}^{-1}, \tilde{q}_k, \tilde{f})$, then $\hat{\xi}_{k+1}^1 \hat{\xi}_{k+1}^2 \in \text{Ker}(\tilde{R}\tilde{A}\tilde{R}^{-1} + \tilde{R}^{-1}\tilde{A}^{\top}\tilde{R})$.
- (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 $\tilde{R}\tilde{A}\tilde{R}^{-1}$ be symmetric. Then any solution of the $VI(-\tilde{R}\tilde{A}\tilde{R}^{-1},\tilde{q}_k,\tilde{f})$ is also a solution to the problem: $\min_{\xi\in\tilde{\Gamma}_k}-\frac{1}{2}\xi^{\top}\tilde{R}\tilde{A}\tilde{R}^{-1}\xi+\langle\tilde{q}_k,\xi\rangle$.
- (f) Let $-\tilde{R}\tilde{A}\tilde{R}^{-1} \succ 0$ and be symmetric. Then $\xi_{k+1} = \operatorname{proj}_{-\tilde{R}\tilde{A}\tilde{R}^{-1}}[\tilde{\Gamma}_k; -\tilde{q}_k]$.

Proof:

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- 1. Follows from [6, Theorem 1], which is itself inspired from [23] and [16, Theorem 2], and which states that the negative feedback interconnection of a passive system with a maximal monotone operator, defines another maximal monotone 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 $\mathrm{dom}(\varphi) = \mathrm{dom}(\tilde{f}) = \{z \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}z \in \Gamma_k\} = \{z \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}z + \tilde{F}_k \in K\} = \{z \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}z + \tilde{G}\tilde{F}_k \geq 0\}$. Hence $(\mathrm{dom}(\varphi))_{\infty} = P(\tilde{G}\tilde{C}\tilde{R}^{-1}, 0)$ from Proposition A.1 e). Now we have $\tilde{f}(\xi) = \psi_{\Gamma_k}(\tilde{C}\tilde{R}^{-1}\xi) = \psi_{\tilde{\Gamma}_k}(\xi)$ with $\tilde{\Gamma}_k = \{\xi \in \mathbb{R}^p \mid \tilde{C}\tilde{R}^{-1}\xi \in K - \tilde{F}_k\} = \{\xi \in \mathbb{R}^p \mid \tilde{G}\tilde{C}\tilde{R}^{-1}\xi + \tilde{G}\tilde{F}_k \geq 0\} = P(\tilde{G}\tilde{C}\tilde{R}^{-1}, -\tilde{G}\tilde{F}_k)$. Thus $\tilde{f}_{\infty}(\cdot) = \psi_{(\tilde{\Gamma}_k)_{\infty}}(\cdot)$ and $\mathrm{dom}(\tilde{f}_{\infty}) = P(\tilde{G}\tilde{C}\tilde{R}^{-1}, 0)$. Now we have to calculate the dual cone $P(\tilde{G}\tilde{C}\tilde{R}^{-1}, 0)^*$. Using [31, Example 15], one finds that $P(\tilde{G}\tilde{C}\tilde{R}^{-1}, 0)^* = \{w \in \mathbb{R}^{n_2} \mid w = \sum_{i=1}^m \lambda_i (\tilde{G}\tilde{C}\tilde{R}^{-1})_{i\bullet}^{\uparrow}, \ \lambda_i \geq 0\}$, where $(\tilde{G}\tilde{C}\tilde{R}^{-1})_{i\bullet}$ denotes the ith row of the matrix $\tilde{G}\tilde{C}\tilde{R}^{-1} \in \mathbb{R}^{m \times n_2}$. Hence it follows that $K(-\tilde{R}\tilde{A}\tilde{R}^{-1}, \tilde{f}) = \{\xi \in \mathbb{R}^{n_2} \mid \tilde{R}\tilde{A}\tilde{R}^{-1}\xi = -\sum_{i=1}^m \lambda_i (\tilde{G}\tilde{C}\tilde{R}^{-1})_{i\bullet}^{\uparrow}, \ \lambda_i \geq 0\}$. The third set in \tilde{S} is equal to $\ker(\tilde{R}\tilde{A}\tilde{R}^{-1} + \tilde{R}^{-1}\tilde{A}^{-1}\tilde{R})$, from $\ker(\mathbf{M} + \mathbf{M}^{-1})$ in Proposition A.2 a). Then the results in items (a) to (f) are a consequence of Proposition A.2, where $\varphi_{\infty}(v) = 0$ in item (b). Item (f) follows from item e) in Proposition 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 $\bar{B}_I(\mathcal{N}_{K^*})^{-1}(-\tilde{C}z_{k+1}-\tilde{F}_k)$ 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 $\tilde{D} \geq 0$.

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

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

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 > 0 be given.

324 Define
$$\Gamma_{2,k} = K_2 - \tilde{F}_{2,k}$$
. Consider the VI: find $\xi_{k+1} \in \text{dom}(\tilde{f}_2)$ such that (4.9)

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

- 326 where $\tilde{f}_2(\cdot) = \psi_{\Gamma_{2,k}} \circ \tilde{C}_2 \tilde{R}_2^{-1}(\cdot)$.
 - 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 Lip-schitz continuous mapping, is maximal monotone:

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

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

335 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:

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

From (4.10) it is deduced that:

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

Using (4.11) it is inferred:

$$\tilde{A}z_{k+1} - \tilde{B}_{1}(\mathcal{N}_{K_{1}} + \tilde{D}_{11})^{-1}(\tilde{C}_{1}z_{k+1} + \tilde{F}_{1,k}) \\
-\tilde{B}_{2}\mathcal{N}_{K_{2}-\tilde{F}_{2,k}}(\tilde{C}_{2}z_{k+1}) + \tilde{E}_{k} = 0$$

$$\updownarrow \\
-\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1}\xi_{k+1} + \tilde{R}_{2}\tilde{B}_{1}(\mathcal{N}_{K_{1}^{*}} + \tilde{D}_{11})^{-1}(\tilde{C}_{1}z_{k+1} + \tilde{F}_{1,k}) - \tilde{R}_{2}\tilde{E}_{k} \\
\in -\tilde{R}_{2}\tilde{B}_{2}\mathcal{N}_{K_{2}-\tilde{F}_{2,k}}(\tilde{C}_{2}\tilde{R}_{2}^{-1}\xi_{k+1}),$$

- 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).
 - 3. The first assertion is a consequence of [9, Corollary 24.4] (for the maximal 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 23.39]. The second assertion follows from [29, Theorem 2.3.3], noting that the

mapping in the left-hand side of (4.9) is maximal monotone [9, Corollary 24.4 (ii)] and strongly monotone since $-\tilde{R}_2\tilde{A}\tilde{R}_2^{-1}\succ 0$. This can also be proved using [9, Corollary 23.37]. Single-valuedness and Lipschitz continuity stem from $\tilde{D}_{11}\succ 0$. \boxtimes

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 $\tilde{G}_2 : \xi_{k+1} \mapsto \tilde{R}_2\tilde{B}_1(\tilde{D}_{11} + \partial \psi_{K_1^*})^{-1}(\tilde{C}_1\tilde{R}_2^{-1}\xi_{k+1} + \tilde{F}_{1,k})$ 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 $\tilde{\Gamma}_k = \mathbb{R}_+^m$ one may rely on P-matrices and LCPs properties, when it is a closed convex cone one may use copositive matrices [29, Theorem 2.5.10].

374 The next corollary completes the existence and uniqueness results for the OSNSP.

COROLLARY 4.6. The next results hold true:

- (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 à ≺ 0, then z_{k+1} is unique.
- 2. (set of solutions λ_{k+1}) Let \hat{C} Dom $(\mathcal{G}) \subseteq \text{Im}(\partial \sigma_{\Gamma_k} + \hat{D})$, 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} .
- 3. Let Assumption 4 hold with strict passivity, and the conditions of item 1 in Proposition 4.3 be true. Then the GE in (3.5) (or the GE in (3.3)) has a unique solution (λ_{k+1}, z_{k+1}) for any \tilde{E}_k . Consequently x_{k+1} in (3.2) is 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.

Proof:

- 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

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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})
\tilde{D}) the result follows.
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- 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.
- 4. Under these conditions $Dom(\mathcal{G})$ is bounded (compact). Hence from [9, Proposition 23.36 (iii) existence of a solution holds true. From [29, Corollary 2.2.5] the compactness of the set of solutions follows. This could also be inferred 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), 411 in which case the \in can be replaced by equality. Concerning item 2 in Corollary 4.6, 412 an important step is to characterize both $\text{Dom}(\mathcal{G})$ and $\text{Im}(\partial \sigma_{\Gamma_k} + \hat{D})$. If $\hat{D} = \hat{D}^{\top} \geq 0$, 413 one can apply [15, Theorem 3], and infer that $\operatorname{Im}(\partial \sigma_{\Gamma_k} + \tilde{D}) \simeq \operatorname{Im}(\partial \sigma_{\Gamma_k}) + \operatorname{Im}(\tilde{D}) =$ 414 $\operatorname{Dom}(\partial \psi_{\Gamma_k}) + \operatorname{Im}(\hat{D}) \subset \operatorname{dom}(\psi_{\Gamma_k}) + \operatorname{Im}(\hat{D}) = \Gamma_k + \operatorname{Im}(\hat{D})$ (see [15] for the exact 415 meaning of \simeq in this context). Similar calculations may be done for in item 1 of 416 Proposition 4.3. It is inferred that the basic condition in Proposition 4.3 item 1, 417 $i.e., \operatorname{Im}(\hat{C}) \cap \operatorname{int}(\operatorname{Im}(\partial \sigma_{\Gamma_k} + \hat{D})) \neq \emptyset$, or in Corollary 4.6 item 2, i.e., $\hat{C} \operatorname{Dom}(\mathcal{G}) \subseteq$ 418 $\operatorname{Im}(\partial \sigma_{\Gamma_k} + D)$, always are satisfied when D > 0. For the sake of briefness the conditions 419 420 for the existence and uniqueness of λ_{k+1} in (4.12) are not stated, similar tools as for Corollary 4.6 can be used. 421
- 422 **4.2.** Analysis of the OSNSP for (3.1). It is also interesting to investigate 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 $\Gamma(t_k, x_{k+1}) = \{\xi \in \mathbb{R}^m \mid \xi + C_I x_{k+1} + F_k \in \mathbb{R}^m \mid \xi + C_I x_{k+1$ 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 > 0, hold. For 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:
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430 (4.14)
$$-RB_{II}R^{-1}\xi(t) - RA_{II}x(t) - RE_{II}(t) \in -\partial \psi_{\Phi(t,x(t))}(\xi(t)),$$

where
$$\xi = -Rz$$
, $R^2 = X$, $R = R^{\top} > 0$, $\Phi(t, x) = \{ \xi \in \mathbb{R}^p \mid C_{II}R^{-1}\xi + C_Ix + F(t) \in K \}$, and the set-valued term is obtained using the chain rule, under the Convex

Analysis chain rule qualification condition [49, Theorem 23.9]. The first step is to 433 434

analyse the existence and uniqueness of solutions to the VI:

435 (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} 436 and the other terms and variables are considered as data for this VI. Similar results 437 438 to those in item 3 of Proposition 4.3 can be obtained.

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

442 in (4.15). Assume that
$$-RB_{II}R^{-1} \geq 0$$
. Consider the set: (4.16)

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

444 Then:

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- 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 solution.
- 2. If $S \neq \{0\}$, and if there exists $\xi_0 \in \Phi(t_k, x_{k+1})$ such that $\langle -RA_{II}x_{k+1} RE_{II,k} + RB_{II}R^{-1}\xi_0, v > 0$ for all $v \neq 0$ and $v \in \mathcal{S}$, then it follows that the $VI(-RB_{II}R^{-1}, -RA_{II}x_{k+1} - RE_{II,k}, \psi_{\Phi(t_k, x_{k+1})})$ has a unique solution.
- **Proof:** The set S in (4.16) is constructed as follows. The chain rule holds from 450 [49, Theorem 23.9] and the fact that $\psi_{K-C_Ix_{k+1}+F_k}(\cdot)$ is polyhedral. Hence the 451 VI in (4.15) is equivalent to (3.1) (b) (c). Here $\varphi = \psi_{\Phi(t_k, x_{k+1})}$, hence $\operatorname{dom}(\varphi) =$ 452
- $\Phi(t_k, x_{k+1})$ and $(\text{dom}(\varphi))_{\infty}$. Assume that K = P(G, 0) (see Proposition A.1 (e)), 453
- 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, $\varphi_{\infty} = 0$ 454
- $\psi_{(\Phi(t_k, x_{k+1}))_{\infty}} = \psi_{P(GC_{II}R^{-1}, 0)}, \text{ hence } \text{dom}(\varphi_{\infty}) = P(GC_{II}R^{-1}, 0) \text{ and } \text{dom}(\varphi_{\infty})^{\star} = (P(GC_{II}R^{-1}, 0))^{\star} = \{w \in \mathbb{R}^{n_2} \mid w = -\sum_{i=1}^{m} \lambda_i (GC_{II}R^{-1})_{i\bullet}^{\top}, \lambda_i \geq 0\}, \text{ where } [31, 15]$ 456
- Example 15] was used as in the proof of Proposition 4.3. Thus $\mathcal{K}(-RB_{II}R^{-1},\varphi) =$ 457
- $\{\xi \in \mathbb{R}^{n_2} \mid -RB_{II}R^{-1}\xi = -\sum_{i=1}^m \lambda_i (GC_{II}R^{-1})_{i\bullet}^{\top}, \ \lambda_i \geq 0\}.$ Finally in Proposition A.2: $\text{Ker}(\mathbf{M} + \mathbf{M}^{\top}) = \text{Ker}(RB_{II}R^{-1} + (RB_{II}R^{-1})^{\top})$ and $R = R^{\top}$. Items 1 and 2 459
- are a direct consequence of Proposition A.2 (a) and (b). \boxtimes 460
- For the sake of briefness we do not state similar results as in Proposition 4.3 (c) to 461 (f). To enlarge the analysis let us now state the counterpart of Assumption 6: 462

463 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 464
- B_{II} and C_{II} , and of the cone K (see the comments after Assumption 6, in particular 465
- if $D_{11} = 0$ the case D = 0 is recovered with $G_2(\cdot, \cdot)$ in (4.17) which vanishes, and the set $\Phi_2 = \Phi$), let us state the following:
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- Assumption 9. There exists a $(p \times p)$ -matrix $X_2 = X_2^{\top} \succ 0$ such that $X_2 \bar{B}_{II}^2 =$ 468 469 $(C_{II}^2)^{\top}$.
- Let us recall that p is the number of equality constraints in (2.1). Let us define
- $R_2 = R_2^{\top} \succ 0, \ R_2^2 = X_2, \ \xi_{k+1} = R_2 z_{k+1}, \ \Phi_2(t_k, x_{k+1}) = \{\xi \in \mathbb{R}^p \mid C_{II}^2 R_2^{-1} \xi + C_I^2 x_{k+1} + F_{2,k} \in K_2\}.$ The next proposition is the counterpart of Proposition 4.4. It 472
- states conditions that guarantee that z_{k+1} can be calculated from (3.1) (b) (c). 473
- Proposition 4.8. Let h > 0, and Assumptions 8 and 9 hold true. Consider the 474 VI: $find \xi_{k+1} \in dom(f_2)$ such that 475

$$\begin{array}{ccc}
-R_{2}B_{II}R_{2}^{-1}\xi_{k+1} - R_{2}A_{II}x_{k+1} - R_{2}E_{II,k} \\
+ \underbrace{R_{2}\bar{B}_{II}^{1}(D_{11} + \partial\psi_{K_{1}^{\star}})^{-1}(C_{II}^{1}R_{2}^{-1}\xi_{k+1} - F_{1,k})}_{\triangleq G_{2}(t_{k},\xi_{k+1})} \in -\partial f_{2}(\xi_{k+1})
\end{array}$$

477 with
$$f_2(\cdot) = (\psi_{K_2 - C_I^2 x_{k+1} - F_{2,k}} \circ C_{II}^2 R_2^{-1})(\cdot) = \psi_{\Phi_2(t_k, x_{k+1})}(\cdot).$$

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 478 equivalent to the MLCP in (3.1) (b) (c). 479

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

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

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

489 $R_2A_{II}x_{k+1} - R_2E_{II,k} \ni 0$, and the assumption guarantees that $H_2(\cdot)$ is maximal 490 monotone [9, Corollary 24.4]. The result follows from [9, Proposition 23.39].

491 The same as Lemma 4.5 can be stated to guarantee the maximal monotonicity in items 3 and 4. It is noteworthy that the main discrepancy between Proposition 492 4.4 and Proposition 4.8, is that the solution ξ_{k+1} depends on x_{k+1} in the latter, 493 while it depends on x_k in the former. This is a difficulty with the application of 494 Proposition 4.8 in some cases, for instance when $C_{II} = 0$ (examples in sections 6.2) 495 and 6.4 will demonstrate that this can occur). Actually, the condition $\operatorname{Im}(C_{II}^2) \cap (K_2 C_I^2 x_{k+1} - F_{2,k} \neq \emptyset$ (that is a form of constraint qualification), or its counterparts in 497 Propositions 4.3, 4.4, 4.7, happens to be crucial in the setting of this work. It states that the algebraic variable z has to enter the complementarity constraints through 499 a C_{II} matrix which has a big enough range. This may occur in some systems, see section 6.5. 501

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

516 (5.1)
$$P_I^h x_{k+1} = P_I x_k + h \bar{E}_k(t_k, x_k)$$

for some $\bar{E}_k(t_k, x_k)$ which is a linear function of $E_{I,k}$, $E_{II,k}$, λ_{k+1} and z_{k+1} . The time-dependence of \bar{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:

521 (5.2)
$$x_k = ((P_I^h)^{-1}P_I)^k x_0 + h \sum_{j=0}^{k-1} ((P_I^h)^{-1}P_I)^{k-1-j} (P_I^h)^{-1} \bar{E}_j(t_j, x_j), \ k \ge 1,$$

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

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

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

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

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

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

526 (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, 527 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 529 530 bounded for all $k \geq 0$, 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 $\mathcal{Z}: \zeta \mapsto \eta =$ $(\partial \sigma_{\Gamma_k} + \tilde{D})^{-1}(\zeta)$ is under the conditions of Proposition 4.3 item 2, a well-defined $\mathbb{R}^m \to \mathbb{R}^m$ 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 the 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 $\lambda_{\min}(\tilde{D} + \tilde{D}^{\top}) > 0$ uniformly in h > 0 only if D > 0. Thus the operator \mathcal{G} : $\zeta \mapsto -\tilde{A}\zeta + \tilde{B}(\partial\sigma_{\Gamma_k} + \tilde{D})^{-1}(\tilde{C}\zeta)$ is also single-valued Lipschitz continuous (being the composition of Lipschitz continuous mappings) with a Lipschitz constant that does not depend on x_k , and which is bounded for bounded h. Let us now assume the following which reinforces Assumption 4 and guarantees that $D \succ 0$:

Assumption 10. The system $(\tilde{A}, \tilde{B}, \tilde{C}, \tilde{D})$ is strictly passive, uniformly in h > 0.

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 4.3 item 1, maximal monotonicity follows from [6, Theorem 1]. Strong monotonicity follows from the fact that under Assumption 10, the quadruple $(\ddot{A} + \epsilon I_{n_2}, \ddot{B}, \ddot{C}, \ddot{D})$ is passive for a small enough $\epsilon > 0$, so that $(\mathcal{G} - \epsilon I_{n_2})(\cdot)$ is maximal monotone for such ϵ . Thus if Assumption 10 and the basic conditions in item 1 of Proposition 4.3 are true, $\mathcal{G}: \mathbb{R}^p \to \mathbb{R}^p$ is well-defined and single-valued, Lipschitz continuous and 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 Lipschitz, and the solution z_{k+1} of the generalized equation $\tilde{E}_k(t_k, x_k) \in \mathcal{G}(z_{k+1}) \Leftrightarrow$ $z_{k+1} = \mathcal{G}^{-1}(E_k(t_k, x_k))$ in (3.5) is Lipschitz continuous in x_k and t_k . Thus λ_{k+1} is also Lipschitz continuous in x_k , and we conclude that the same applies to x_{k+1} in (3.2).

Remark 5.1. Assumption 10 is stringent. Actually the strict state passivity of $(\hat{A}, \hat{B}, \hat{C}, \hat{D})$ 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 $\tilde{D} > 0$, 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 $\bar{B}_I = 0$).

5.1.2. Boundedness and convergence analysis. Under the above conditions, 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 single-valued and Lipschitz continuous in x_k and t_k (provided $E_I(\cdot)$, $E_{II}(\cdot)$ and $F(\cdot)$ are Lipschitz continuous functions of t, which is a mild requirement on the data), since \tilde{E}_k and \tilde{F}_k are linear in the argument x_k . We have $||\bar{E}(t_k^1, x_k^1) - \bar{E}(t_k^2, x_k^2)|| \le L_t |t_k^1 - t_k^2| + L_x ||x_k^1 - x_k^2||$ 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 $||\bar{E}(t_k, x_k^1) - \bar{E}(t_k, x_k^2)|| \le L_x ||x_k^1 - x_k^2||$ and $||\bar{E}(t_k, x_k) - \bar{E}(t_k, 0)|| \le L_x ||x_k||$. Notice that due to the Lipschitz continuity both z_{k+1} and λ_{k+1} are uniformly bounded in 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 $||\bar{E}(t_k, x_k)|| \le ||\bar{E}(t_k, 0)|| + L_x ||x_k|| \le m + L_x ||x_k||$. Now using (5.2) we get: (5.5)

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

where we used the Hölder induced matrix norm $||\cdot||_{2,2} = \sigma_{\text{max}}(\cdot)$ and [10, Proposition 9.4.9, Corollary 9.4.12]. Let us make the following:

Assumption 11. One has
$$\sigma_{\max}((P_I^h)^{-1}P_I) \leq 1$$
 for all $h \geq 0$.

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

589 (5.6)
$$||x_k|| \le ||x_0|| + h\sigma_{\min}(P_I^h) \sum_{i=1}^k (m + L_x||x_{j-1}||),$$

Assume now that $x_0, x_1, ..., x_{k-1}$ are uniformly bounded by some x_{\max} , then from (5.6) one has $||x_k|| \le x_{\max} + h\sigma_{\min}(P_I^h)k(m + L_x x_{\max}) \le x_{\max} + h\sigma_{\min}(P_I^h)n(m + L_x x_{\max}) = x_{\max} + h\sigma_{\min}(P_I^h)\frac{T}{h}(m + L_x x_{\max}) = x_{\max} + \sigma_{\min}(P_I^h)T(m + L_x x_{\max})$. 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:

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

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 the family $\{x^h\}_{h\geq 0}$ is equi-Lipschitz continuous. Thus from Arzela-Ascoli Theorem

 $\{x^h(\cdot)\}\$ stays in a compact subset of $\mathcal{C}^0([0,T];\mathbb{R}^{n_1})$ and there exist subsequences 601 which converge uniformly in $C^0([0,T];\mathbb{R}^{n_1})$ towards a Lipschitz continuous limit $x(\cdot)$. 602 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-603 schitz continuity of $E_I(\cdot)$, $E_{II}(\cdot)$, and $F(\cdot)$, the sequences of functions $\{\tilde{E}^h(\cdot)\}_{h\geq 0}$ and $\{\tilde{F}^h(\cdot)\}_{h>0}$ also converge uniformy in $\mathcal{C}^0([0,T];\mathbb{R}^p)$ and $\mathcal{C}^0([0,T];\mathbb{R}^m)$ to limits $\tilde{E}(\cdot) = E_{II}(\cdot) + A_{II}x(\cdot)$ and $\tilde{F}(\cdot) = F(\cdot) + C_{I}x(\cdot)$ as $h \to 0$, respectively. From (5.4) it 606 is deduced that $\dot{x}^h(\cdot)$ is uniformly bounded. Therefore from the Banach-Alaoglu The-607 orem, it follows that $\dot{x}^h(\cdot)$ converges weakly* to a limit $v(\cdot)$. Since the domain [0,T] is 608 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-609 quently $\dot{x}^h(\cdot)$ converges weakly in $\mathcal{L}^2([0,T];\mathbb{R}^{n_1})$. Since $x^h(t)=x^h(0)+\int_0^t \dot{x}^h(\tau)d\tau$, 610 we infer that $v(\cdot) = \dot{x}(\cdot)$ almost everywhere. 611 Now notice that using (3.2) the following holds: 612

$$\begin{cases}
(a) P_{I}\dot{x}^{h}(t) = A_{I}x_{h}^{\star}(t) + B_{I}z_{h}^{\star}(t) + \bar{B}_{I}\lambda_{h}^{\star}(t) + E_{I}^{h}(t) \\
(b) \tilde{A}z_{h}^{\star}(t) + \tilde{B}\lambda_{h}^{\star}(t) + \tilde{E}_{II}^{h}(t) + A_{II}x^{h}(t) = 0 \\
(c) 0 \leq \lambda_{h}^{\star}(t) \perp \tilde{C}z_{h}^{\star}(t) + \tilde{D}\lambda_{h}^{\star}(t) + F^{h}(t) + C_{I}x^{h}(t) \geq 0,
\end{cases}$$

for all $t \in [t_k, t_{k+1}), k \geq 0, x_h^{\star}(t) = x_{k+1}$ for all $t \in [t_k, t_{k+1})$ and similarly for the other terms are step-functions. From the boundedness of $\|\dot{x}_k\| \stackrel{\Delta}{=} \|\frac{x_{k+1}-x_x}{h}\| \le C$ for some C > 0, we infer that $x_{\star}^h(\cdot)$ converges strongly in $\mathcal{L}^2([0,T];\mathbb{R}^{n_1})$ towards $x(\cdot)$. 615 616 617 (5.9)

$$\|x^{h} - x_{\star}^{h}\|_{\mathcal{L}_{2}([0,T];\mathbb{R}^{n_{1}})}^{2} \leq \sum_{k=0}^{n-1} \int_{t_{k}}^{t_{k+1}} (t_{k+1} - t_{k})^{2} \|\dot{x}_{k}\|^{2} dt$$

$$\leq C^{2} \sum_{k=0}^{n-1} (t_{k+1} - t_{k})^{3} = \frac{C^{2} n h^{3}}{3} = \frac{C^{2} T h^{2}}{3} \longrightarrow 0 \text{ as } h \to 0.$$

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

Proposition 5.2. Let the conditions of item 2 in Corollary 4.6 hold true uni-635 formly in $h (\Rightarrow D \succ 0)$. Assume that Assumptions 10 and 11 hold, $E(\cdot)$ and $F(\cdot)$ are 636 bounded Lipschitz continuous functions, initial data are bounded. Then the DALCS 637 in (2.1) possesses for each admissible bounded initial data, a continuous solution 638 $(x,z,\lambda):[0,T]\to \mathbb{R}^{n_1+n_2+m},\ T>0.$ Moreover the Euler approximations converge 639 to a solution as $h \to 0$.

5.1.3. The case blockdiagonal $\tilde{D} = \tilde{D}^{\top} \succeq 0$. This section is dedicated to 641 analyse an alternative to the conditions stated in sections 5.1.1 and 5.1.2. Let us 642 consider the case of Proposition 4.4 with Assumption 6. Then it follows from the 643 complementarity conditions that $\lambda_{1,k+1} = (\tilde{D}_{11} + \partial \psi_{K_1^*})^{-1}(-\tilde{C}_1 z_{k+1} - \tilde{F}_{1,k})$, where the mapping $(\tilde{D}_{11} + \partial \psi_{K_1^*})^{-1}(\cdot)$ is maximal monotone (being the sum of two maximal 645 monotone mappings satisfying [9, Corollary 24.4 (ii)]), single-valued, well-defined and 646 Lipschitz continuous, see Proposition B.1. Thus if $\bar{B}_I = (\bar{B}_I^1, 0)$ with $\bar{B}_I^1 \in \mathbb{R}^{n_1 \times m_1}$ 647 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 righthand side of (3.5) (c). 649

650 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 $\tilde{X}_2 \succ 0$ for any $h \geq 0$.

This secures that the variable change $\xi_{k+1} = -\tilde{R}_2^{-1} z_{k+1}$ remains well-defined as $h \to 0$. It is also important that item 1 in Proposition 4.4 holds, so that the VI in (4.9) is equivalent to the MLCP in (3.2) (b) (c).

LEMMA 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 $||x_k||$.

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

$$\begin{split} & -\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1}\xi_{k+1} + \tilde{R}_{2}\tilde{E}_{k} + \tilde{G}_{2}(t_{k},\xi_{k+1}) \in -\partial\tilde{f}_{2}(\xi_{k+1}) \\ \Leftrightarrow & -\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1}(\xi_{k+1} - \xi_{0}) + \tilde{R}_{2}\tilde{E}_{k} + \tilde{G}_{2}(t_{k},\xi_{k+1}) + \zeta_{0} - \tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1}\xi_{0} = -\zeta_{k+1} + \zeta_{0} \\ & \text{with } \zeta_{k+1} \in \partial\tilde{f}_{2}(\xi_{k+1}), \zeta_{0} \in \partial\tilde{f}_{2}(\xi_{0}) \\ \Rightarrow & \langle -\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1}(\xi_{k+1} - \xi_{0}), \xi_{k+1} - \xi_{0}\rangle + \langle \tilde{R}_{2}\tilde{E}_{k}, \xi_{k+1} - \xi_{0}\rangle \\ & + \langle \tilde{G}_{2}(t_{k},\xi_{k+1}) - \tilde{G}_{2}(t_{k},\xi_{0}), \xi_{k+1} - \xi_{0}\rangle + \langle \zeta_{0} - \tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1}\xi_{0}, \xi_{k+1} - \xi_{0}\rangle \\ & + \langle \tilde{G}_{2}(t_{k},\xi_{0}), \xi_{k+1} - \xi_{0}\rangle = -\langle \zeta_{k+1} - \zeta_{0}, \xi_{k+1} - \xi_{0}\rangle \leq 0, \end{split}$$

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 $||\tilde{R}_2\tilde{B}_1(\tilde{D}_{11}+\partial\psi_{K_1^{\star}})^{-1}(z_1)-\tilde{R}_2\tilde{B}_1(\tilde{D}_{11}+\partial\psi_{K_1^{\star}})^{-1}(z_2)|| \leq \eta \ ||z_1-z_2||$ 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)

$$\begin{array}{c} \left\langle -\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1}(\xi_{k+1}-\xi_{0}), \xi_{k+1}-\xi_{0} \right\rangle \leq \left\langle -\tilde{R}_{2}\tilde{E}_{k}, \xi_{k+1}-\xi_{0} \right\rangle - \left\langle \tilde{G}_{2}(t_{k},\xi_{0}), \xi_{k+1}-\xi_{0} \right\rangle \\
+ \left\langle \zeta_{0} + \tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1}\xi_{0}, \xi_{k+1}-\xi_{0} \right\rangle + \eta' \left| |\xi_{k+1}-\xi_{0}| \right|^{2}
\end{array}$$

which implies (5.12)

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

Now one has:
$$||\tilde{E}_{k}|| \leq ||E_{II,k} + hA_{II}E_{I,k}|| + ||A_{II}(P_{I}^{h})^{-1}P_{I}|| ||x_{k}||, \quad \tilde{G}_{2}(t_{k}, \xi_{0}) = \tilde{R}_{2}\tilde{B}_{1}(\tilde{D}_{11} + \partial \psi_{K_{1}^{*}})^{-1}(\tilde{C}_{1}\tilde{R}_{2}^{-1}\xi_{0} - \tilde{F}_{1,k}), \quad ||\tilde{G}_{2}(t_{k}, \xi_{0}) - \tilde{G}_{2}(t_{0}, \xi_{0})|| = ||\tilde{R}_{2}\tilde{B}_{1}(\tilde{D}_{11} + \partial \psi_{K_{1}^{*}})^{-1}(\tilde{C}_{1}\tilde{R}_{2}^{-1}\xi_{0} - \tilde{F}_{1,k}) - \tilde{R}_{2}\tilde{B}_{1}(\tilde{D}_{11} + \partial \psi_{K_{1}^{*}})^{-1}(\tilde{C}_{1}\tilde{R}_{2}^{-1}\xi_{0} - \tilde{F}_{1,0})|| \leq \eta \mid|-\tilde{F}_{1,k}| + \eta \mid|C_{I}(P_{I}^{h})^{-1}P_{I}|| \cdot ||x_{k}||. \quad \text{Let us denote } \alpha_{k} = \tilde{R}_{1,0}||\hat{R}_{1}|| \cdot ||x_{k}|| \cdot ||x$$

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$$||\tilde{R}_2|| ||E_{II,k} + hA_{II}E_{I,k}|| + \eta ||F_k + hC_I(P_I^h)^{-1}E_{I,k}||, \beta_0 = ||\tilde{G}_2(t_0, \xi_0)|| + ||\zeta_0 + \tilde{R}_2\tilde{A}\tilde{R}_2^{-1}\xi_0||, \text{ and } \gamma = \eta ||C_I(P_I^h)^{-1}P_I|| + ||\tilde{R}_2|| ||A_{II}(P_I^h)^{-1}P_I||. \text{ Then:}$$

(5.13)

$$\langle (-\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1} - \eta' I_{n_{2}})(\xi_{k+1} - \xi_{0}), \xi_{k+1} - \xi_{0} \rangle \leq (\alpha_{k} + \beta_{0} + \gamma ||x_{k}||)||\xi_{k+1} - \xi_{0}|| \\ \leq \frac{1}{2}\alpha_{k}^{2} + \frac{1}{2}\beta_{0}^{2} + \frac{1}{2}\gamma^{2} ||x_{k}||^{2} + \frac{1}{2}||\xi_{k+1} - \xi_{0}||^{2}$$

678 which implies:

(5.14)

679
$$\langle (-\tilde{R}_2 \tilde{A} \tilde{R}_2^{-1} - \eta' I_{n_2} - \frac{1}{2} I_{n_2})(\xi_{k+1} - \xi_0), \xi_{k+1} - \xi_0 \rangle \leq \frac{1}{2} \alpha_k^2 + \frac{1}{2} \beta_0^2 + \frac{1}{2} \gamma^2 \|x_k\|^2.$$

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

$$(5.15) \qquad \begin{aligned} ||\xi_{k+1} - \xi_{0}|| &\leq \frac{1}{\sqrt{\lambda_{\min}(-\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1} - \eta'I_{n_{2}} - \frac{1}{2}I_{n_{2}})}} \sqrt{\frac{1}{2}\alpha_{k}^{2} + \frac{1}{2}\beta_{0}^{2} + \frac{1}{2}\gamma^{2} ||x_{k}||^{2}} \\ &\leq \frac{1}{\sqrt{2}\sqrt{\lambda_{\min}(-\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1} - \eta'I_{n_{2}} - \frac{1}{2}I_{n_{2}})}} \left(\sqrt{\alpha_{k}^{2} + \beta_{0}^{2}} + \gamma ||x_{k}||\right) \end{aligned}$$

- 682 which ends the proof. \boxtimes
- Therefore the term hB_Iz_{k+1} in the right-hand side of (3.5) (c), admits also an affine
- upper bound in $||x_k||$. A further characterization of the iterates is given now.
- 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
- 686 $\tilde{G}_2(\cdot)$ is monotone (see Lemma 4.5 for conditions), and that $E_I(\cdot)$, $E_{II}(\cdot)$ and $F(\cdot)$
- are uniformly bounded and Lipschitz continuous. Then for any k and $p \geq 0$ and any
- 688 solution of the VI in (4.9), one has:

689 (5.16)
$$\|\xi_k - \xi_p\| \le \alpha + \beta \|x_{k-1} - x_{p-1}\|$$

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

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

with $\zeta_k \in \partial f_2(\xi_k)$, $\zeta_n \in \partial f_2(\xi_n)$, it follows that

$$\langle -\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1}(\xi_{k}-\xi_{p}), \xi_{k}-\xi_{p}\rangle + \langle \tilde{R}_{2}(\tilde{E}_{k-1}-\tilde{E}_{p-1}), \xi_{k}-\xi_{p}\rangle$$

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

$$= -\langle \zeta_{k}-\zeta_{p}, \xi_{k}-\xi_{p}\rangle$$

695 which implies

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

697 One has:

(5.20)

$$\langle -\tilde{R}_{2}(\tilde{E}_{k-1} - \tilde{E}_{p-1}), \xi_{k} - \xi_{p} \rangle$$

$$\leq \langle -\tilde{R}_{2}(E_{II,k-1} - E_{II,p-1}), \xi_{k} - \xi_{p} \rangle - \langle \tilde{R}_{2}(hA_{II}(E_{I,k-1} - E_{I,p-1})), \xi_{k} - \xi_{p} \rangle$$

$$-\langle A_{II}(P_{I}^{h})^{-1}P_{I}(x_{k-1} - x_{p-1}), \xi_{k} - \xi_{p} \rangle$$

$$\leq \frac{1}{2} \|\tilde{R}_{2}((E_{II,k-1} - E_{II,p-1}) + hA_{II}(E_{I,k-1} - E_{I,p-1}))\|^{2} + \frac{1}{2} |\xi_{k} - \xi_{p}\|^{2}$$

$$+ \frac{1}{2} \|A_{II}(P_{I}^{h})^{-1}P_{I}\|^{2} \|x_{k-1} - x_{p-1}\|^{2} + \frac{1}{2} \|\xi_{k} - \xi_{p}\|^{2}$$

$$\leq \|\tilde{R}_{2}\|^{2} \|E_{II,k-1} - E_{II,p-1}\|^{2} + \|\xi_{k} - \xi_{p}\|^{2} + h\|A_{II}\|^{2} \|E_{I,k-1} - E_{I,p-1}\|^{2}$$

$$+ \frac{1}{2} \|A_{II}(P_{I}^{h})^{-1}P_{I}\|^{2} \|x_{k-1} - x_{p-1}\|^{2} \leq \alpha' + \beta' \|x_{k-1} - x_{p-1}\|^{2} + \|\xi_{k} - \xi_{p}\|^{2},$$

for some $\alpha' \geq 0$ 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 $\alpha'' \geq 0$ proportional to $(t_{k-1} - t_{p-1})^2$ (since $E(\cdot)$ and $F(\cdot)$ are Lipschitz continuous), and $\beta'' \geq 0$. Hence using (5.19): 702 703

704 (5.22)
$$\|\xi_{k} - \xi_{p}\| \leq \frac{1}{\sqrt{\lambda_{\min}(-\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1} - \frac{3}{2}I_{p})}} \sqrt{\alpha' + \alpha'' + (\beta' + \beta'') \|x_{k-1} - x_{p-1}\|^{2}}$$
$$\leq \frac{1}{\sqrt{\lambda_{\min}(-\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1} - \frac{3}{2}I_{p})}} (\sqrt{\alpha' + \alpha''} + \sqrt{\beta' + \beta''} \|x_{k-1} - x_{p-1}\|)$$

705 and the lemma is proved. \boxtimes

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COROLLARY 5.5. Let h > 0. Assume that the conditions and assumptions of 707 Lemma 5.4 are fulfilled. For all $k \ge 1$ and any solution of the VI in (4.9), one has: 708

709 (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 $\beta \geq 0$. 710
- **Proof:** Let us consider the proof of Lemma 5.4, replacing k by k+1 and p by k-1. 711
- Then α in (5.16) is proportional to h. Redoing the calculations in (5.20) and (5.22), 712
- one obtains (5.23). 713
- Clearly for all h > 0 (5.23) can be written equivalently as: $\|\xi_{k+1} \xi_k\| \le h\alpha +$ 714
- $\beta ||x_k x_{k-1}||$. The bounds calculated for the variable ξ_k also hold for the variable 715
- $z_k = -\tilde{R}_2 \xi_k$. Let us now analyse λ_{k+1} in (3.5) (a). As seen above, $\lambda_{1,k+1} = (D_{11} + C_{12})$ 716
- 717
- $\partial \psi_{K_1^*})^{-1}(-\tilde{C}_1z_{k+1}-\tilde{F}_{1,k})$ and $\lambda_{2,k+1}\in -\mathcal{N}_{K_2}(\tilde{C}_2z_{k+1}+\tilde{F}_{2,k})$. 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 718
- continuous with respect to the variables z_{k+1} , x_k and F_k . In a more general setting 719
- one obtains: 720

721 (5.24)
$$P_{I}^{h}x_{k+1} \in P_{I}x_{k} + hB_{I}z_{k+1} + h\bar{B}_{I}^{1}(\tilde{D}_{11} + \partial\psi_{K_{1}^{\star}})^{-1}(-\tilde{C}_{1}z_{k+1} - \tilde{F}_{1,k}) - h\bar{B}_{I}^{2}\mathcal{N}_{K_{2}}(\tilde{C}_{2}z_{k+1} + \tilde{F}_{2,k}).$$

- The difficulty is that the last term in the right-hand side of (5.24) is set-valued in 722
- general, and this is intrinsically due to the structure of the implicit scheme in (2.1). 723
- A solution is to consider the VI in (4.9). Its right-hand side matches with the last 724
- term in the right-hand side of (5.24) and it is set-valued. However provided the VI 725
- is well-posed, we know that one element of this set is equal to the left-hand side of 726
- (4.9). We may assume that the scheme in (5.24) is implemented with this element. 727
- 728
- This element appears to define a Lipschitz continuous function of ξ_{k+1} (let us recall that the Lipschitz constant of the third term is equal to $\frac{1}{\lambda_{\min}(\tilde{D}_{11})}$, see Proposition 729
- B.1). If this element is chosen then the numerical scheme (5.24) becomes: 730

731 (5.25)
$$P_{I}^{h}x_{k+1} = P_{I}x_{k} + hB_{I}z_{k+1} + h\bar{B}_{I}^{1}(\tilde{D}_{11} + \partial\psi_{K_{1}^{*}})^{-1}(\tilde{C}_{1}z_{k+1} - \tilde{F}_{1,k}) + h\bar{B}_{I}^{2}(-\tilde{R}_{2}\tilde{A}\tilde{R}_{2}^{-1}\xi_{k+1} - \tilde{R}_{2}\tilde{E}_{k} + \tilde{G}_{2}(t_{k}, \xi_{k+1})).$$

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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 $C^0([0,T];\mathbb{R}^{n_1})$ towards a Lipschitz continuous limit denoted as $x(\cdot)$.

Proof: In the proof, use is made of positive bounded constants denoted generically as k_i . Consider (5.25). First recall that the operator $\zeta \mapsto (\tilde{D}_{11} + \partial \psi_{K_1^*})^{-1}(\zeta)$ is Lipschitz continuous with Lipschitz constant $\frac{1}{\lambda_{\min}(\tilde{D}_{11})}$. Hence we obtain:

$$||(P_{I}^{h})^{-1}\bar{B}_{I}^{1}(\tilde{D}_{11}+\partial\psi_{K_{1}^{\star}})^{-1}(-\tilde{C}_{1}z_{k+1}-\tilde{F}_{1,k}) - (P_{I}^{h})^{-1}\bar{B}_{I}^{1}(\tilde{D}_{11}+\partial\psi_{K_{1}^{\star}})^{-1}(-\tilde{C}_{1}z_{p+1}-\tilde{F}_{1,p})||$$

$$\leq k_{1}||\tilde{C}_{1}(z_{k+1}-z_{p+1}+\tilde{F}_{1,k}-\tilde{F}_{1,p}||\leq k_{2}||z_{k+1}-z_{p+1}||+k_{3}||\tilde{F}_{1,k}-\tilde{F}_{1,p}||$$

$$\leq k_{4}||\xi_{k+1}-\xi_{p+1}||+k_{4}||F_{1,k}-F_{1,p}||+k_{5}||E_{I,k}-E_{I,p}||+k_{6}||x_{k}-x_{p}||$$

$$\leq k_{4}(\alpha+\beta||x_{k}-x_{p}||)+k_{7}|t_{k}-t_{p}|+k_{6}||x_{k}-x_{p}||\leq k_{8}+k_{9}||x_{k}-x_{p}||$$

$$+k_{7}|t_{k}-t_{p}|,$$

for any k and $p \ge 0$, where we used Lemma 5.4 and the Lipschitzness property. In a 745 quite similar way one has $||(P_I^h)^{-1}\bar{B}_I^2(\tilde{G}_2(\xi_{k+1})-\tilde{G}_2(\xi_{p+1}))|| \le k_8+k_9||x_k-x_p||+$ 746 $k_7|t_k-t_p|$, where k_7 , k_8 and k_9 depend on the Lipschitz constants of $\tilde{G}_2(\cdot,\cdot)$, $F(\cdot)$, and 747 on $\|(P_I^h)^{-1}\bar{B}_I^2\|$. One also has $\|(P_I^h)^{-1}(B_I+\bar{B}_I^2(\tilde{R}_2\tilde{A}\tilde{R}_2^{-1})(z_{k+1}-z_{p+1}))\| \le k_{10}(\alpha+1)$ 748 $|\beta||x_k - x_p|| \le k_{11} + k_{12}||x_k - x_p||$, and $||\bar{B}_I^2 \tilde{R}_2 (\tilde{E}_k - \tilde{E}_p)|| \le k_{13}|t_k - t_p| + k_{14}||x_k - x_p||$. 749 It is deduced that $||x_{k+1} - x_{p+1}|| \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k - x_p|| + hk_{16}|t_k - t_p| + hk_{17}$, 750 where the term k_{17} stems from $k_8 + k_{11} = \alpha(k_4 + k_{10})$. Here we are led to make the 751 following which is more stringent that Assumption 11: 752

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

Obviously for h = 0 and the Hölder-induced norm one has $||(P_I^h)^{-1}P_I||_{2,2} + hk_{15} = 1$. 755Let us choose p = 0, it follows that $||x_{k+1} - x_1|| \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k - x_1|| \le (||(P_I^h)^{-1}P_$ 756 $|x_0|| + hk_{16}T + hk_{17}$, hence $||x_{k+1}|| \le (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k|| + hk_{16}T + hk_{17} + hk_{17}$ $||x_1|| + (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_0||$. Now it has to be proved that $||x_1||$ is bounded. 758 From Lemma 5.3 and (5.15), $||\xi_1||$ and $||z_1||$ are upperbounded by an affine function of initial data $||x_0||$ and $||z_0||$. The operator $(\tilde{D}_{11} + \partial \psi_{K_1^*})^{-1}(\cdot)$ is well-defined and single-valued. Taking k = 0 in (5.25) one infers that indeed $||x_1||$ is bounded for bounded 760 761 initial data. Consequently it follows that $||x_{k+1}|| \leq (||(P_I^h)^{-1}P_I|| + hk_{15}) ||x_k|| + hk_{16}T + hk_{17} + hf(x_0, z_0) + ||(P_I^h)^{-1}P_I|| ||x_0||$ for some function $f(\cdot)$ that is bounded 762 763 for bounded initial data. To simplify the notations let us set $h_1 \stackrel{\Delta}{=} ||(P_I^h)^{-1}P_I|| + hk_{15}$ 764 and $h_2 \stackrel{\Delta}{=} k_{16}T + k_{17} + f(x_0, z_0)$. We obtain for all $k \geq 1$: 765

766 (5.27)
$$||x_k|| \leq h_1^k ||x_0|| + \sum_{j=0}^{k-1} h_1^j (hh_2 + ||(P_I^h)^{-1} P_I|| ||x_0||) \\ \leq h_1^k ||x_0|| + (hh_2 + ||(P_I^h)^{-1} P_I|| ||x_0||) \frac{1 - h_1^{k+1}}{1 - h_1}.$$

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

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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 (||(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}) \frac{||x_k - x_{k-1}||}{h} + hk_{16} + k_{17}. Now using Lemma 5.4 and the Lipschitzness it follows that \alpha, hence k_{17}, are proportional to k_E^2 h^2. We may therefore write the inequality as \frac{||x_{k+1} - x_k||}{h} \le (||(P_I^h)^{-1}P_I|| + hk_{15}) \frac{||x_k - x_{k-1}||}{h} + hk_{16} + k_{18}h^2 for some k_{18} \ge 0. Proceeding as for (5.27) one infers that if \frac{||x_{1} - x_{0}||}{h} is bounded then \frac{||x_{k+1} - x_{k}||}{h} is bounded for all k \ge 1. Evaluating (5.25) at k = 0 and using the boundedness of x_1 as shown above, proves that this is indeed the case. The proposition is proved using the Arzela-Ascoli Theorem. \boxtimes
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Clearly in practice one needs to compute the constants k_i , for instance k_{15} is needed to determine h_{max} in Assumption 13.

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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 $C^0([0,T];\mathbb{R}^{n_2})$ towards a Lipschitz continuous limit denoted as $z(\cdot)$ (respectively $\xi(\cdot)$).

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

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Let us now pass to the analysis of the multiplier convergence. Recall that $\lambda_{1,k+1} = (\tilde{D}_{11} + \partial \psi_{K_1^{\star}})^{-1}(-\tilde{C}_1 z_{k+1} - \tilde{F}_{1,k})$ 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 $t \geq 0$, $\lambda_{1,\star}^h(t) = (\tilde{D}_{11} + \partial \psi_{K_1^{\star}})^{-1}(-\tilde{C}_1 z_{\star}^h(t) - \tilde{F}_{1,\star}^h(t))$ 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 $\dot{\xi}_k \triangleq \|\frac{\xi_{k+1} - \xi_k}{h}\| \leq C$ for some C > 0, hence:

794 (5.28)
$$\|\xi^{h} - \xi_{\star}^{h}\|_{\mathcal{L}_{2}([0,T];\mathbb{R}^{n_{2}})}^{2} \leq \sum_{k=0}^{n-1} \int_{t_{k}}^{t_{k+1}} (t_{k+1} - t_{k})^{2} \|\dot{\xi}_{k}\|^{2} dt \\ \leq C^{2} \sum_{k=0}^{n-1} (t_{k+1} - t_{k})^{3} = \frac{C^{2} n h^{3}}{3} = \frac{C^{2} T h^{2}}{3}.$$

Consequently $\|\xi^h - \xi_{\star}^h\|_{\mathcal{L}_2([0,T];\mathbb{R}^{n_2})} \to 0$ as $h \to 0$. Since $\|\xi^h - \xi_{\star}^h\|_{\mathcal{L}_2([0,T];\mathbb{R}^{n_2})} = 1$ $\|\xi^h \pm \xi - \xi_{\star}^h\|_{\mathcal{L}_2([0,T];\mathbb{R}^{n_2})} \to \|\xi - \xi_{\star}^h\|_{\mathcal{L}_2([0,T];\mathbb{R}^{n_2})} \to 0$, we conclude that the sequence $\{\xi_{\star}^h\}_{h\geq 0}$ converges strongly in $\mathcal{L}_2([0,T];\mathbb{R}^{n_2})$ towards the limit $\xi(\cdot)$. The same result holds for the sequence $\{x_{\star}^h\}_{h\geq 0}$ with the limit $x(\cdot)$. Therefore the sequence $\{-\tilde{C}_1z_{\star}^h - \tilde{F}_{1,\star}^h\}_{h\geq 0}$ also converges strongly in \mathcal{L}_2 towards a limit, and we infer that the same conclusion holds for both $\{\lambda_{1,\star}^h(\cdot)\}_{h\geq 0}$ and $\{\lambda_{2,\star}^h(\cdot)\}_{h\geq 0}$. 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.

LEMMA 5.8. Let the conditions of Lemmata 5.3, 5.4, Corollary 5.5 be satisfied for all h > 0. The sequences $\{\xi_{\star}^h\}_{h \geq 0}$, $\{x_{\star}^h\}_{h \geq 0}$, $\{\lambda_{1,\star}^h(\cdot)\}_{h \geq 0}$ and $\{\lambda_{2,\star}^h(\cdot)\}_{h \geq 0}$ converge 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.

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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 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 the limits verify (2.1) (b) (c). Using the same arguments as those developed just above (5.8), it is deduced that $\{\dot{x}^h(\cdot)\}$ converges weakly* to a limit $v(\cdot)$ where $v(\cdot) = \dot{x}(\cdot)$ almost everywhere. Consequently $P_I\dot{x}^h(t) = A_Ix_h^{\star}(t) + B_Iz_h^{\star}(t) + \bar{B}_I\lambda_h^{\star}(t) + E_I^h(t) \longrightarrow$ $P_I\dot{x}(t) = A_Ix(t) + B_Iz(t) + \bar{B}_I\lambda(t) + E_I(t)$ weakly in $\mathcal{L}^2([0,T];\mathbb{R}^{n_1})$ as $h \to 0$.

- **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 $D \geq 0$ as in Proposition 4.8. Indeed, since the strict passivity of $(B_{II}, B_{II}, C_{II}, D)$ implies Assumption 10 for small enough h > 0, it is unnecessary to repeat the developments in sections 5.1.1 and 5.1.2. However, it is not possible in general to determine the implications between Assumption 7 and Assumption 9, nor the relations between $\lambda_{\min}(A)$ and $\lambda_{\min}(B_{II})$ for h > 0. Therefore let us deal with the VI in (4.17) where x_{k+1} is replaced by $x_{k+\theta} = \theta x_k + (1-\theta)x_{k+1}, \ \theta \in [0,1].$ The following holds, which is the counterpart of Lemma 5.3.
- LEMMA 5.10. Assume that $\lambda_{\min}(-R_2B_{II}R_2^{-1})$ (= $\lambda_{\min}(-B_{II})$) is large enough, $E_I(\cdot)$, $E_{II}(\cdot)$ and $F(\cdot)$ are uniformly bounded and Lipschitz continuous functions of time. Then any solution $\xi_{k+1} = R_2^{-1}z_{k+1}$ of the VI in (4.17) is upperbounded by an 832 834 affine function of $||x_{k+\theta}||$. 835
- The proof follows the same steps as the proof of Lemma 5.3 and it is therefore omitted. 836 Lemma 5.10 implies that the term $B_I z_{k+1}$ in the right-hand side of (3.1) (a), is 837 upperbounded by an affine function of $||x_{k+\theta}||$. The coefficients of the affine function 838 depend only on the system's data (constant matrices, Lipschitz constants) and on the 839 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, 841 single-valued and Lipschitz continuous mapping $G_2(\cdot)$ is monotone, and that $E_I(\cdot)$, 842 $E_{II}(\cdot)$ and $F(\cdot)$ are Lispchitz continuous functions of time. Then for any k and $p \geq 0$ 843 and any solution of the VI in (4.17), one has: 844

845 (5.29)
$$\|\xi_k - \xi_p\| \le \alpha + \beta \|x_{k-1+\theta} - x_{p-1+\theta}\|$$

- for all $k \geq 1$ and some bounded constants $\alpha \geq 0$ (proportional to $|t_{k-1} t_{p-1}|$) and 846 $\beta \geq 0$.
- **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

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

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

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$$\langle -R_{2}B_{II}R_{2}^{-1}(\xi_{k+1} - \xi_{p+1}), \xi_{k+1} - \xi_{p+1} \rangle$$

$$\leq \langle R_{2}A_{II}(x_{k+\theta} - x_{p+\theta}) + R_{2}(E_{II,k} - E_{II,p}), \xi_{k+1} - \xi_{p+1} \rangle$$

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

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

$$\leq \|R_{2}A_{II}\|^{2} \|x_{k+\theta} - x_{p+\theta}\|^{2} + \|R_{2}\|^{2}k_{E}^{2}(t_{k} - t_{p})^{2}$$

$$+ \frac{1}{2}k_{G_{2}}^{2} \|F_{k} - F_{p}\|^{2} + \|\xi_{k+1} - \xi_{p+1}\|^{2}$$

$$\leq \|R_{2}A_{II}\|^{2} \|x_{k+\theta} - x_{p+\theta}\|^{2} + \|R_{2}\|^{2}k_{E}^{2}(t_{k} - t_{p})^{2} + \frac{1}{2}k_{G_{2}}^{2}k_{F}^{2}(t_{k} - t_{p})^{2}$$

$$+ \|\xi_{k+1} - \xi_{p+1}\|^{2}.$$

854 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: (5.32)

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

859 equivalently: (5.33)

$$(P_{I}^{h} + h(1 - \theta)\bar{B}_{I}^{2}R_{2}A_{II})x_{k+1} = P_{I}x_{k} + hB_{I}z_{k+1} + h\bar{B}_{I}^{1}(D_{11} + \partial\psi_{K_{1}^{*}})^{-1}(C_{II}^{1}R_{2}^{-1}\xi_{k+1} + F_{1,k}) + h\bar{B}_{I}^{2}(-R_{2}B_{II}R_{2}^{-1}\xi_{k+1} - \theta R_{2}A_{II}x_{k} - R_{2}E_{II,k} + G_{2}(t_{k}, \xi_{k+1})),$$
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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 $(B_{II}, \bar{B}_{II}, C_{II}, D)$ (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 an LCS by solving the equality constraint for z. This yields the LCS quadruple $(P_I^{-1}(A_I - B_IB_{II}^{-1}A_{II}), P_I^{-1}(\bar{B}_I - B_IB_{II}^{-1}\bar{B}_{II}), C_I - C_{II}B_{II}^{-1}A_{II}, D - C_{II}B_{II}^{-1}\bar{B}_{II})$. It is unclear how to relate in general the passivity of this quadruple, to the above assumptions which focus on the MLCP part (2.1) (b) (c) of the system. Especially no specific assumptions are made on the matrices A_I , B_I , \bar{B}_I , C_I in the above analyses. For instance Assumption 3 implies that the feedthrough matrix of the LCS is $D - C_{II}B_{II}^{-1}X^{-1}C_{II}^{\top}$, which may be indefinite in general. It is also noteworthy that transforming the MLCP into an LCP may be undesirable [59].

5.3. Recapitulation. Let us recapitulate the various assumptions upon which the results in this article rely. Assumptions 2, 4, 7, 10, 12 (respectively Assumptions 3, 5, 9), characterize the passivity properties of an auxiliary operator defined by the 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" $w_{k+1} = \tilde{C}z_{k+1} + \tilde{D}\lambda_{k+1}$ (respectively λ and $w = C_{II}z + D\lambda$). Assumptions 6 and 8 have practical importance (because some systems possess such D matrix), they are

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4.4 cannot be used.

also related to passivity in the sense that passive systems possess positive semidefinite 885 886 throughout matrices. Thus all these assumptions are fundamental ones for the developments in this article, which basically deals with the analysis of two types of gener-887 alised equations (GE): $0 \in A(u) + B(u)$ in (4.7) and (4.15), and $0 \in A(u) + B(u) + C(u)$ in (4.9) and (4.17), where the mappings $A(\cdot)$, $B(\cdot)$ and $C(\cdot)$ are maximal monotone 889 (notice in passing that due to time-dependency, these GEs do not fit with those in [28, Chapter 2], so that in particular [28, Theorem 2F.6] cannot be applied in general). 891 Most importantly, the convergence results (Proposition 5.2, Lemmata 5.3 and 5.4, 892 Corollary 5.5, Proposition 5.6, Corollary 5.7, Lemmata 5.8 and 5.10) rely upon the 893 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 894 and guarantees strong monotonicity properties, allowing us to characterize the prop-895 896 erties of the solutions to the GEs. It would be worth relaxing it in future studies. One path could be to use least norm elements [34]. Finally, it is worth noting that 897 many practical systems like circuits with nonlinear resistors and inductances contain 898 nonlinearities [19, section 4.1]. An extension of the above developments, consists of 899 replacing the linear term A_Ix in (2.1), by a nonlinear term $A_I(x)$, where $A_I(\cdot)$ is 900 globally Lipschitz continuous, and setting $A_I(x_k)$ in (3.1). The matrices \tilde{A} , \tilde{B} , \tilde{C} , \tilde{D} keep the same structure with P_I^{-1} instead of $(P_I^h)^{-1}$, while both \tilde{E}_k and \tilde{F}_k contain 901 902 $A_I(x_k)$. 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 from circuits, mechanics and state-dependent-switch DAE. They aim at illustrating 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:

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

with $x_1(t) = \int_0^t i_1(s)ds$, $x_2(t) = \int_0^t i_2(s)ds$, $x_3 = i_3 = z$, $\lambda = (u_{D_1} \ u_{D_2})^{\top}$ (the voltages across the diodes), $K = K^* = \mathbb{R}_+^2$, $n_1 = 2$, $p = n_2 = 1$, m = 2. One has $D = \begin{pmatrix} \frac{1}{R} & 0 \\ 0 & 0 \end{pmatrix}$, $P_I = I_2$, $A_I = \begin{pmatrix} \frac{-2}{RC} & \frac{1}{RC} \\ \frac{1}{2RC} & \frac{-1}{RC} \end{pmatrix}$, $B_I = \begin{pmatrix} 0 \\ 0 \end{pmatrix}$, $\bar{B}_I = \begin{pmatrix} \frac{1}{R} \\ 0 \end{pmatrix}$ 913 $A_{II} = (\frac{1}{2\mathbf{C}} - \frac{1}{\mathbf{C}}), \ B_{II} = -R, \ \bar{B}_{II} = (0 \quad \frac{3}{2}), \ C_{I} = \begin{pmatrix} \frac{-2}{RC} & \frac{1}{RC} \\ 0 & 0 \end{pmatrix}, \ C_{II} = \begin{pmatrix} 0 \\ 1 \end{pmatrix}.$ 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 $\text{Im}(C_{II}^2) = R$. Since $\bar{B}_{II}^1 = 0$ 915 916 917 the mapping $G_2(\cdot)$ is zero, and item 3 in Proposition 4.8 applies. The quadruple ($B_{II}, \bar{B}_{II}, C_{II}, D$) 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 918 919 920 real [18, Theorem 2.48]. The conditions of Lemmata 5.10 and 5.11 are verified by 921 taking R large enough (notice that all other parameters depend on $\frac{1}{R}$). Here $\hat{A} =$ $B_{II} = -R$, but Assumptions 2 and 6 do not hold, hence the material in Proposition 923

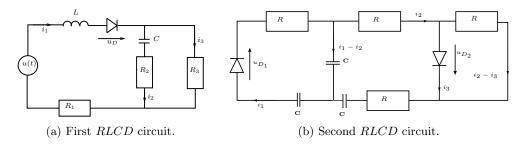


Fig. 1: Two circuits with ideal diodes.

6.2. Example 2 (circuit). Let us consider now the circuit in Figure 1a. Using 925 926 Kirchhoff laws its dynamics may be written as:

927 (6.2)
$$\begin{cases} \dot{x}_1(t) = -\frac{R_1}{L}x_1(t) - \frac{R_3}{L}z(t) + \frac{1}{L}u(t) + \frac{1}{L}u_D(t) \\ \dot{x}_2(t) = \frac{R_3}{R_2}z(t) - \frac{1}{R_2\mathbf{C}}x_2(t) \\ 0 = x_1(t) + \frac{1}{R_2\mathbf{C}}x_2(t) - (1 + \frac{R_3}{R_2})z(t) \\ 0 \le u_D(t) \perp x_1(t) \ge 0, \end{cases}$$

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), $K = K^* = IR_+$, $n_1 = 2$, $p = n_2 = 1$, m = 1. Let apply the feedback control $u = 2R_3z$, 928 929 which modifies B_I . We get: $A_I = \begin{pmatrix} -\frac{R_1}{L} & 0 \\ 0 & -\frac{1}{R_2C} \end{pmatrix}$, $B_I = \begin{pmatrix} \frac{R_3}{L} \\ \frac{R_3}{R_2} \end{pmatrix}$, $\bar{B}_I = \begin{pmatrix} \frac{1}{L} \\ 0 \end{pmatrix}$, 930 $A_{II} = (1 \ \frac{1}{R_2 \mathbf{C}}), B_{II} = -(1 + \frac{R_3}{R_2}), \bar{B}_{II} = 0, C_I = (1 \ 0), C_{II} = 0, D = 0, E_I = \begin{pmatrix} \frac{1}{L}u \\ 0 \end{pmatrix},$ $E_{II} = 0, F(t) = 0, P_I = I_2, P_I^h = \begin{pmatrix} 1 + h \frac{R_1}{L} & 0 \\ 0 & 1 + h \frac{1}{R_2 C} \end{pmatrix} \succ 0 \text{ for all } h \geq 0.$ Thus $\tilde{B} = \frac{h}{hR_1 + L}, \ \tilde{C} = \frac{hR_3}{L + hR_1}, \ \tilde{D} = \frac{h}{L + hR_1}, \ \tilde{A} = -\frac{R_2 + R_3}{R_2} + hR_3 \left(\frac{1}{L + hR_1} + \frac{1}{R_2(h + R_2 \mathbf{C})}\right).$ 933 Assumption 2 holds with $\tilde{X} = R_3$, $\tilde{D} > 0$ for all h > 0, and Assumption 4 holds for h > 0 small enough such that $\tilde{A} < 0$. For h > 0 small enough the set $\tilde{S} = \{0\}$ 934 935 in (4.6) since $\tilde{R}\tilde{A}\tilde{R}^{-1} = \tilde{A} < 0$, hence the VI in (4.7) always has a solution z_{k+1} by 936 Proposition 4.3 item 3 (a). By item (c) (or (f)) of the same proposition, the solution is 937 unique. When $\tilde{A} < 0$, the quadruple $(\tilde{A}, \tilde{B}, \tilde{C}, \tilde{D})$ is strictly passive for all h > 0, hence items 1 and 3 of Corollary 4.6 apply. However Assumption 10 does not hold since 939 $\tilde{D} \to 0$ as $h \to 0$. Consider now the quadruple $(B_{II}, \bar{B}_{II}, C_{II}, D) = (B_{II}, 0, 0, 0)$. 940 It is passive since (C.1) is verified, it is even strictly state passive but not strictly 941 passive since D = 0. Assumptions 3 and 9 are trivially satisfied since both sides of 942 the passivity equality are null. Provided that $\{0\} \in \mathbb{R}_+ -x_{1,k+1} \Leftrightarrow x_{1,k+1} \geq 0$, item 1 943 in Proposition 4.7 applies and the OSNSP for (3.1) has a unique solution (but clearly 944 there is no guarantee that this inequality is satisfied along the integration process, 945 946 rendering Proposition 4.7 unapplicable). Proposition 4.8 item 3 also applies here for the VI in (4.17), however one has to check the equivalence stated in item 1 of the 947 948 same proposition. As noted above, the matrix A_I intervenes only in Assumption 1, hence its modification

949 does not alter the results. Setting $u = 2R_3z + R_1x_1 + aLx_1 + v(t)$ yields the modified

951
$$A_I = \begin{pmatrix} a & 0 \\ 0 & -\frac{1}{R_2 \mathbf{C}} \end{pmatrix}$$
, and B_I as above. This yields (see Remark 5.13): $P_I^{-1}(A_I -$

952
$$B_I B_{II}^{-1} A_{II}$$
) = $\begin{pmatrix} a + \frac{R_2 R_3}{L(R_2 + R_3)} & \frac{R_3}{L\mathbf{C}(R_2 + R_3)} \\ \frac{R_3}{R_2 + R_3} & -\frac{1}{(R_2 + R_3)\mathbf{C}} \end{pmatrix}$, which is unstable if $a + \frac{R_2 R_3}{(R_2 + R_3)L} >$ 953 $\frac{1}{(R_2 + R_3)\mathbf{C}}$. The other elements of the reduced LCS quadruple are given by $P_I^{-1}(\bar{B}_I - R_1)$

954
$$B_I B_{II}^{-1} \bar{B}_{II} = \begin{pmatrix} \frac{1}{L} \\ 0 \end{pmatrix}$$
, $C_I - C_{II} B_{II}^{-1} A_{II} = (1 \ 0)$, $D - C_{II} B_{II}^{-1} \bar{B}_{II} = 0$. Consequently

- the reduced LCS is not passive with the above feedback controller. Also $P_I^h = \text{diag}(1 -$
- 956
- $ha, 1 + \frac{h}{R_2 \mathbf{C}}) \succ 0$ 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

959 (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 DAE). 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

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

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 965 rewritten as: 966

967 (6.5)
$$0 = a\lambda_1(t) + a\lambda_2(t) + b\lambda_3(t) + b\lambda_4(t) - cz(t) + E_{II}(t)$$
$$0 \le \begin{pmatrix} \lambda_1(t) \\ \lambda_2(t) \\ \lambda_3(t) \\ \lambda_4(t) \end{pmatrix} \perp \begin{pmatrix} \lambda_1(t) \\ \lambda_2(t) \\ \lambda_3(t) \\ \lambda_4(t) \end{pmatrix} + \begin{pmatrix} -x(t) \\ x(t) \\ 0 \\ 0 \end{pmatrix} + \begin{pmatrix} 0 \\ 0 \\ -z(t) \\ z(t) \end{pmatrix} \ge 0.$$

- Therefore (6.4) (a) (6.5) fits with (2.1) with $A_{II} = 0$, $B_{II} = -c$, $\bar{B}_{II} = (a\ a\ b\ b)$, $D = I_4$, $C_I = (-1\ 1\ 0\ 0)^{\top}$, $C_{II} = (0\ 0\ -1\ 1)^{\top}$, F(t) = 0, $\bar{B}_I = 0$. Let us check the 968
- 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
- $M \stackrel{\triangle}{=} \begin{pmatrix} 2cA & -aA & -aA & -bA 1 & -bA + 1 \\ -aX & 2 & 0 & 0 \\ -aX & 0 & 2 & 0 & 0 \\ -bX 1 & 0 & 0 & 2 & 0 \end{pmatrix} \succ 0,$
- for some X > 0. Let us denote M = P + N with $P = \operatorname{diag}(2cX, 2I_4) > 0$, and 973
- N_{00} the matrix obtained for a=b=0. Applying [25, Theorem 2.11] or [10, Fact 974
- 8.19.2], it follows that $P + N_{00} > 0$ if $\sigma_{\max}^2(N_{00}) = 2 < \lambda_{\min}(2cX)\lambda_{\min}(2I_4) = 4cX$.
- Using continuity of the eigenvalues [10, Fact 10.11.8], it is deduced that given c > 0, 976
- X > 0, there exists a and b small enough such that (6.6) holds (or, given a, b and 977
- X > 0, there exists c > 0 large enough). Thus, item 1 in Lemma 4.1 holds, item 2 978
- in Proposition 4.3 holds for h>0 small enough (guaranteeing $D \succ 0$), Assumption

10 holds true for small enough h > 0 (see the proof of item 1 in Lemma 4.1), and 980 the arguments stated at the beginning of section 5.2 about the convergence analysis 981 led in sections 5.1.1 and 5.1.2, hold. We have: $\tilde{A} = -c$, $\tilde{B} = (a\ a\ b\ b)$, $\tilde{C} = (0\ 0\ -1\ 1)^{\top} + h(-1\ 1\ 0\ 0)^{\top}(1-hA_I)^{-1}B_I = (\frac{-hB_I}{1-hA_I}\ \frac{hB_I}{1-hA_I}\ 0\ 0)^{\top}$, $\tilde{D} = D$. 982 983 Integration can proceed using (3.2). 984

Let us now pass to a second example of state-dependent-switch DAE: 985

986 (6.7)
$$\begin{cases} (a) \dot{x}_1(t) = a_1 x_1(t) + b_1 x_2(t) + c_1 z(t) + e_1(t) \\ (b) \dot{x}_2(t) = a_2 x_1(t) + b_2 x_2(t) + c_2 z(t) + e_2(t) \\ (c) 0 = a_3 x_1(t) + b_3 x_2(t) + c_3 z(t) + e_3(t) \text{ if } x_1(t) < 0 \\ (d) 0 = a_4 x_1(t) + b_3 x_2(t) + c_3 z(t) + e_3(t) \text{ if } x_1(t) > 0, \end{cases}$$

where the switching conditions in (c) and (d) are equivalently rewritten as: 987

988 (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:
$$A_I = \begin{pmatrix} a_1 & b_1 \\ a_2 & b_2 \end{pmatrix}$$
, $B_I = (c_1 c_2)^{\top}$, $\bar{B}_I = (0 0)^{\top}$, $E_I = (e_1 e_2)^{\top}$, $\bar{B}_{II} = (-a_3 a_4)^{\top}$, $A_{II} = (0 b_3)$, $B_{II} = c_3$, $E_{II} = e_3$, $C_I = \begin{pmatrix} 1 & 0 \\ -1 & 0 \end{pmatrix}$, $C_{II} = (0 0)^{\top}$, $D_I = (-a_1 a_2)^{\top}$, $D_I = (-a_1 a_2)^{\top$

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observed in Figure 2 that the order is approximately equal to one for small enough

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

step size. Figure 3 depicts the time evolution of various variables, demonstrating that 1013 the complementarity constraints are respected. Clearly the multipliers are continuous but non differentiable when the complementarity switches from one mode to another.

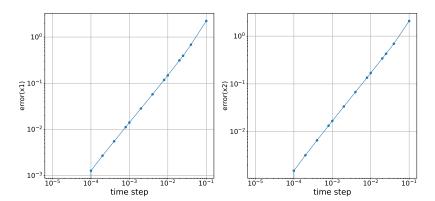


Fig. 2: Log-log error diagram: x_1 (left) and x_2 (right).

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

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

- where $M = M^{\top} \succ 0$, $D \succ 0$, $x = (q^{\top} \dot{q}^{\top})^{\top}$, $q(t) \in \mathbb{R}^{\frac{n_1}{2}}$, $z = \lambda_b \in \mathbb{R}^{n_2}$, $\lambda = \lambda_u \in \mathbb{R}^m$, $G \in \mathbb{R}^{n_2 \times \frac{n_1}{2}}$, $f : \mathbb{R}^+ \to \mathbb{R}^{n_2}$ 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

$$1024 \quad (6.10) \quad 0 = G\ddot{q}(t) - \ddot{f}(t) = GM^{-1}H^{\top}\lambda(t) + GM^{-1}G^{\top}z(t) - GM^{-1}(K_p K_v)x(t) - \ddot{f}(t),$$

- 1025
- 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) = $p=n_2\Rightarrow n_1\geq 2n_2$), $A_{II}=GM^{-1}(K_p\ K_v)$, $\bar{B}_{II}=-GM^{-1}H^{\top}$, $B_I=\begin{pmatrix} 0\\ -M^{-1}G^{\top} \end{pmatrix}$, and $C_{II}=0$ (similarly to the example in
- section 6.2, one major difference being that here D > 0). Let us examine the OSNSP 1028
- for (6.9) on a particular example as depicted in Figure 4, whose dynamics are given
- 1030

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,

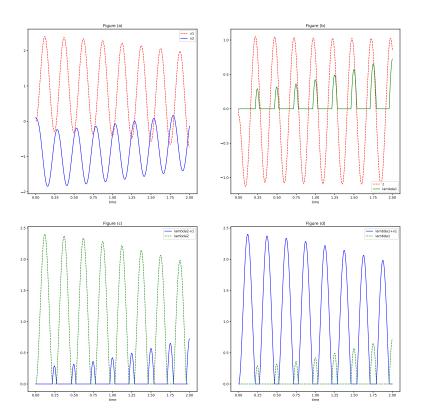


Fig. 3: (a) $x_1(t)$ (dashed) and $x_2(t)$, (b) z(t) (dashed) and $\lambda_1(t)$, (c) $\lambda_2(t) - x_1(t)$ and $\lambda_2(t)$ (dashed), (d) $\lambda_1(t) + x_1(t)$ and $\lambda_1(t)$ (dashed).

```
1034 \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 = 1035 \begin{pmatrix} 1 & 0 \\ -1 & 0 \end{pmatrix}, \ M = \mathrm{diag}(m,m), \ D = I_2, \ A_{II} = 0, \ B_{II} = \frac{-1}{m}, \ \bar{B}_{II} = 0, \ E_{II}(t) = \ddot{f}(t), \ 1036 \ C_I = \begin{pmatrix} k & 0 & 0 & 0 \\ -k & 0 & 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}, \ 1037 \ 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}, \ 1038 \ \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 > 0 \ \text{if } h > 0 \ \text{is small} \ 1039 \ \text{enough}. Let us check Assumption 4. The passivity LMI is \begin{pmatrix} \frac{2X}{m} & \tilde{C} \\ \tilde{C}^{\top} & \tilde{D} + \tilde{D}^{\top} \end{pmatrix} > 0, \ 1040 \ X > 0, which is verified for h > 0 small enough and X > 0 large enough. Hence under these conditions the quadruple (\tilde{A}, \tilde{B}, \tilde{C}, \tilde{D}) is strictly passive. It can be checked that items 1 and 2 in Proposition 4.3 hold, as well as items 1, 2 and 3 in Corollary 4.6.
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Therefore the dynamics can be integrated using (3.2) with $q_{2,0} + f_0 = 0$, $\dot{q}_{2,0} + f_0 = 0$. 1043 We note also that the quadruple $(B_{II}, \bar{B}_{II}, C_{II}, D) = (\frac{-1}{m}, 0, 0, I_2)$ is strictly passive (Assumption 5 is verified), hence item 1 in Lemma 4.1 holds. Assumption 10 is 1045 also verified (for h>0 small enough but uniformly in $h\geq 0$). Let us analyse the convergence as in section 5.1. Assumption 11 may not be verified in this case, however 1048 as noted after this assumption it is easy to modify slightly the numerical scheme so that it is satisfied with $P_I = I_4$. Hence it is inferred that Proposition 5.2 applies. 1049

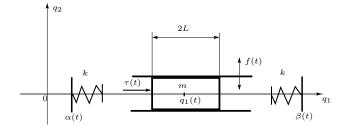


Fig. 4: A mechanical DALCS.

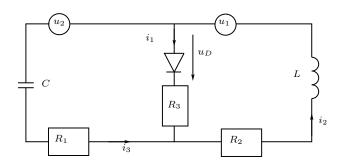


Fig. 5: An *RLCD* passive circuit.

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 1051 diode, $u_1(\cdot)$ and $u_2(\cdot)$ are two voltage sources. One way to write its dynamics is: 1052

1053 (6.12)
$$\begin{cases} \dot{x}_1(t) = -\frac{R_2}{L} x_1(t) + \frac{1}{L} \lambda(t) + \frac{1}{L} u_1(t) \\ \dot{x}_2(t) = x_1(t) - z(t) \\ 0 = \frac{1}{C} x_2(t) + \lambda(t) + u_2(t) \\ 0 \le \lambda(t) \perp z(t) \ge 0. \end{cases}$$

Assume that the feedback $u_2 = -az + u_3$ is applied, so that the equality constraint becomes $0 = \frac{1}{\mathbf{C}}x_2(t) + u_D(t) - az(t) + u_3(t)$. We have: $A_{II} = (0 \frac{1}{\mathbf{C}}), C_I = 0$, $C_{II} = 1, B_{II} = -a, \bar{B}_{II} = 1, D = 0, \bar{B}_{I} = (\frac{1}{L} \ 0)^{\top}, B_{I} = (0 \ -1)^{\top}, \text{ and } \tilde{D} = D = 0,$ $\tilde{C} = C_{II} = 1, \tilde{A} = -a + hA_{II}(P_{I}^{h})^{-1}B_{I}, \tilde{B} = 1 + hA_{II}(P_{I}^{h})^{-1}\bar{B}_{I}.$ The condition $\text{Im}(C_{II}) \cap (K - C_{I}x_{k+1} + F_{k}) = \mathbb{R} \cap \mathbb{R}_{+} - \frac{1}{C}x_{2,k+1} \neq \emptyset$ is verified always. The set S in (4.16) is equal to $\{0\}$ hence Proposition 4.7 item 1 applies and the VI in (4.15) always 1056 1059 has a solution. Proposition 4.8 items 1, 2, 3 apply, with $D_{11} = 0$ in Assumption 8.

The system $(B_{II}, \bar{B}_{II}, C_{II}) = (-a, 1, 1)$ is passive and Assumptions 3 and 5 hold true. Also the system $(\tilde{A}, \tilde{B}, \tilde{C})$ is passive for h > 0 small enough and Assumptions 2 and 4 hold true. The set S in (4.6) is equal to $\{0\}$ so Proposition 4.3 item 3 (a) applies. For h > 0 small enough, $\tilde{A} < 0$ so Corollary 4.6 items 1 and 2 apply. Assumption 12 is verified for $h \geq 0$ small enough with $\tilde{X} = \tilde{B}^{-1}$. It is inferred that if a > 0 is large enough, the convergence results in sections 5.1.3 and 5.2 apply.

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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 one-step nonsmooth problem is studied first, then the convergence of the discrete-time solutions is analysed. This article only scratches the surface of singular nonsmooth systems analysis. The above results could be extended in various directions. Most importantly, it would be worth extending the results to more general indices and relative 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. Some basic notions of Convex Analysis and Complementarity Theory are briefly introduce now [11, 29, 49, 31]. Let $f: \mathbb{R}^n \to \mathbb{R} \cup \{+\infty\}$ be a proper convex and lower semicontinuous (lsc) function, we denote by $\operatorname{dom}(f) \stackrel{\Delta}{=} \{x \in \mathbb{R}^n | f(x) < +\infty\}$ the domain of the function $f(\cdot)$. The Fenchel transform $f^*(\cdot)$ of $f(\cdot)$ is the proper, convex and lsc 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-differential $\partial f(x)$ of $f(\cdot)$ at $x \in \mathbb{R}^n$ is defined by: $\partial f(x) = \{\omega \in \mathbb{R}^n | f(v) - f(x) \ge \langle \omega, v - x \rangle, \forall v \in \mathbb{R}^n \}$. We denote by $\operatorname{Dom}(\partial f) \stackrel{\Delta}{=} \{x \in \mathbb{R}^n | \partial f(x) \ne \emptyset \}$ the domain of the subdifferential operator $\partial f: \mathbb{R}^n \to \mathbb{R}^n$.

Let x_0 be any element in the domain $\operatorname{dom}(f)$, 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]:

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

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

1090 Important properties of the recession function and recession cone are recalled now:

Proposition A.1. [11, Proposition 1.4.8] The following statements hold:

1092 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 1093 lsc functions. Suppose that $f_1 + f_2$ is proper. Then for all $x \in \mathbb{R}^n$: $(f_1 + f_2)_{\infty}(x) = 1094$ $(f_1)_{\infty}(x) + (f_2)_{\infty}(x)$.

1095 b) Let $f: \mathbb{R}^n \to \mathbb{R} \cup \{+\infty\}$ be a proper, convex and lsc function, and let K be a nonempty closed convex set, such that $f + \Psi_K$ is proper (equivalently $dom(f) \cap K$ is 1097 non empty). Then for all $x \in \mathbb{R}^n$: $(f + \Psi_K)_{\infty}(x) = f_{\infty}(x) + (\Psi_K)_{\infty}(x)$.

- c) Let $K \subset \mathbb{R}^n$ be a nonempty, closed and convex set. Then for all $x \in \mathbb{R}^n$:
- $(\Psi_K)_{\infty}(x) = \Psi_{K_{\infty}}(x)$. Moreover for all $x \in K$ and $e \in K_{\infty}$: $x + e \in K$.
- 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 $K_{\infty} = P(A, 0) = \{x \in \mathbb{R}^n | Ax \ge 0\}$.
- Let us now concatenate [4, Theorem 3, Corollaries 3 and 4]. They concern variational
- inequalities (VIs) of the form: Find $u \in \mathbb{R}^n$ such that 1104
- $\langle \mathbf{M}u + \mathbf{q}, v u \rangle + \varphi(v) \varphi(u) \ge 0$, for all $v \in \mathbb{R}^n$ 1105(A.1)
- where $\mathbf{M} \in \mathbb{R}^{n \times n}$ is a real matrix, $\mathbf{q} \in \mathbb{R}^n$ a vector and $\varphi : \mathbb{R}^n \to \mathbb{R} \cup \{+\infty\}$ 1106
- a proper convex and lsc function. The VI in (A.1) is equivalent to the inclusion
- $\mathbf{M}u + \mathbf{q} \in -\partial \varphi(u)$, which is a generalized equation (GE) of the form: $0 \in A(u) + B(u)$
- where $A(\cdot)$ and $B(\cdot)$ are two maximal monotone operators (since it will be assumed 1109
- $\mathbf{M} \geq 0$. 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
- 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
- $\mathcal{K}(\mathbf{M}, \varphi) = \{x \in \mathbb{R}^n | \mathbf{M}x \in (\text{dom}(\varphi_{\infty}))^* \}.$ 1116 (A.2)
- Note that $(\operatorname{dom}(\varphi_{\infty}))^{\star}$ 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
- $dom(\varphi)$. 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 $\mathbf{M} \geq 0$. 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
- $VI(\mathbf{M}, \mathbf{q}, \varphi)$ has at least one solution. 1123
- 1124 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 $x_0 \in$

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

- $dom(\varphi)$ such that 1125 (A.3)
- 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 $u_1 u_2 \in \operatorname{Ker}\{\mathbf{M} +$ 1128
- 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 $\langle \mathbf{q}, u_1 u_2 \rangle = \varphi(u_2) \varphi(u_1)$. 1130
- 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
- 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
- If $K \subseteq \mathbb{R}^n$ is a set, then $K^* = \{z \in \mathbb{R}^n | \langle z, x \rangle \geq 0 \text{ for all } x \in K\}$ is 1135
- its dual cone. Let K be a nonempty closed convex cone, then: 1136
- $K^* \ni x \perp y \in K \iff x \in -\mathcal{N}_K(y) \iff y \in -\mathcal{N}_{K^*}(x).$ (B.1)

Let $M = M^{\top} \succ 0$, x and y two vectors, then [29] 1138

1139 (B.2)
$$M(x-y) \in -\mathcal{N}_K(x) \Leftrightarrow x = \operatorname{proj}_M[K; y] \Leftrightarrow x = \min_{z \in K} \frac{1}{2} (z-y)^\top M(z-y).$$

- 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 and Lipschitz continuous with modulus $L = \frac{2}{\lambda_{\min}(D+D^{\top})}$. 1145
- 1146
- It was proposed in [16, Proposition 1] when $\mathcal{M} = \partial \phi$ for some convex proper lsc $\phi(\cdot)$,
- extended to maximal monotone $\mathcal{M}(\cdot)$ in [5]. 1148
- **Appendix C. Passive Systems.** A quadruple (A, B, C, D) is said passive if: 1149

1150 (C.1)
$$\begin{pmatrix} -XA - A^{\top}X & -XB + C^{\top} \\ -B^{\top}X + C & D + D^{\top} \end{pmatrix} \succcurlyeq 0, \quad X = X^{\top} \succ 0.$$

- for some X. It is said *strictly* passive if the first matrix inequality is satisfied with ≥ 0 1151
- (this implies that D > 0). This is not to be confused with the *strict state* passivity [18, 1152
- Definition 4.54 (strict passivity implies strict state passivity). In fact strict passivity
- is directy related to strongly SPR transfer matricess [18, p.62 and section 3.12.2], 1154
- while SPR transfer matrices with minimal realization are strictly state passive [18,
- Theorem 4.73 [43]. Strict state passivity holds when $-XA A^{\top}X > 0$ in (C.1), with 1156
- $D+D^{\top} \geq 0$. The two inequalities in (C.1) make a linear matrix inequality (LMI) with 1157
- unknown X, equivalently a Riccati inequality using the Schur complement Theorem 1158
- 1159 [18, Theorem A.65]. Numerical algorithms exist that efficiently solve LMIs, see [13],
- and can be used to test whether or not a quadruple is passive. 1160

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