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1 AN ALGORITHM SOLVING COMPRESSIVE SENSING PROBLEMS 2 BASED ON MAXIMAL MONOTONE OPERATORS*

3 YOHANN TENDERO^T, IGOR CIRIL[‡], JÉRÔME DARBON[§], AND SUSANA SERNA[¶]

Abstract. The need to solve ℓ^1 regularized linear problems can be motivated by various com-4 pressive sensing and sparsity related techniques for data analysis and signal or image processing. 5 These problems lead to non-smooth convex optimization in high dimensions. Theoretical works pre-6 dict a sharp phase transition for the exact recovery of compressive sensing problems. Our numerical experiments show that state-of-the-art algorithms are not effective enough to observe this phase tran-8 9 sition accurately. This paper proposes a simple formalism that enables us to produce an algorithm 10 that computes an ℓ^1 minimizer under the constraints $A\boldsymbol{u} = \boldsymbol{b}$ up to the machine precision. In addition, a numerical comparison with standard algorithms available in the literature is exhibited. The 11 12 comparison shows that our algorithm compares advantageously with other state-of-the-art methods, 13 both in terms of accuracy and efficiency. With our algorithm, the aforementioned phase transition 14 is observed at high precision.

15 **Key words.** Sparse solution recovery, Compressive sensing, Inverse scale space, ℓ^1 minimization, 16 Non-smooth optimization, Maximal monotone operator, Phase transition.

17 **AMS subject classifications.** 34A60 ,49M29 ,90C06 ,90C25.

1. Introduction. Compressive sensing and sparsity-related paradigms have gained enormous interest in the last decade and can be used for, e.g., data analysis, signal and image processing, inverse problems or acquisition devices. Indeed, in many cases the unknowns of an under-determined system can be obtained by finding the sparsest (or simplest) solution to a linear system

$$23 \quad (1.1) \qquad \qquad A\boldsymbol{u} = \boldsymbol{b}.$$

With this formulation **b** is the observed data, $A \in \mathcal{M}_{m \times n}(\mathbb{R})$, $m \ll n$ and the columns 24 of A represent a suitable frame or dictionary able to sparsely encode or observe $\boldsymbol{u} \in \mathbb{R}^n$. 2526However, finding a minimizer of the ℓ^0 pseudo-norm under the constraints (1.1) is a highly non-convex and non-smooth optimization problem. Hence, methods [19, 25, 27 30, 35, 39, 3, 14, 28] that aim at tackling ℓ^0 pseudo norm minimization guarantee 28 an optimal solution only with high probability and for a specific class of matrices A. 29Another class of methods consists in using an ℓ^1 relaxation. The problem therefore 30 31 becomes

32
$$(P_{\ell^1})$$

$$\begin{cases} \inf_{\boldsymbol{u}\in\mathbb{R}^n} & \|\boldsymbol{u}\|_{\ell^1} \\ \text{s.t.} & A\boldsymbol{u} = \mathbf{b}. \end{cases}$$

It turns out that under various assumptions, the minimizers remain the same if one replaces the ℓ^0 pseudo-norm by the ℓ^1 norm (see, e.g., [11, 12, 16, 17] and the references

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therein). Problem (P_{ℓ^1}) is a convex albeit non-smooth optimization problem in high 35 36 dimension (n can be thought as the number of pixels of an image for instance). For these reasons developing efficient algorithmic solutions is still a challenge in many 37 cases. For instance, the CVX system "is not meant for very large problems" [20, 38 Sec. 1.3, p.3] that arise from signal/image processing applications [24, 37]. Hence, 39 many algorithms have been proposed to solve ℓ^1 minimization problems, see, e.g., [34, 40 21, 2, 8, 9, 44, 45, 43, 42, 18, 40]. In this paper, we propose a simple algorithm 41 that can be employed to solve these ℓ^1 minimization problems up to the machine 42 precision. Indeed, it is only assumed that the matrix A has full row rank. This 43paper also exhibits a numerical comparison with several classic algorithms in the 44 literature. These comparisons illustrate that our algorithm compares advantageously: 4546 the theoretically predicted phase transition, see e.g. [29, 10], is empirically observed with a higher accuracy. 47

To design our algorithm, we required that: i) the method computes a solution to (P_{ℓ^1}) up to the machine precision, and that ii) the method requires few computations involving vectors of length n.

The first requirement can be thought as guaranteeing the quality of the solution or the fidelity to the problem. The second requirement can be thought as promoting the numerical efficiency. Indeed, computations with vectors of length $m \ll n$ require less memory than the memory needed for vectors of the primal. (We recall that the unknown \boldsymbol{u} lives in a high dimensional space, while the observed data \boldsymbol{b} lives in a space of dimension $m \ll n$). It seems unrealistic to find a minimizer to (P_{ℓ^1}) up to the machine precision with a direct method. Consequently, the approach we employ is iterative and can be summarized as follows.

To the best of our knowledge, the most similar approach to the one developed in this paper is the AISS [7] method. AISS iterates over two variables: a primal one 60 that belongs to \mathbb{R}^n and a dual one in \mathbb{R}^m . Instead, we compute *one* finite discrete 61 sequence λ_k for k = 1, ..., K in \mathbb{R}^m . The last iterate, namely λ_K , is an solution to 62 63 the dual problem of (P_{ℓ^1}) up to the machine precision. Given λ_K a simple formula allows us to compute a solution $\bar{\boldsymbol{u}}$ to (P_{ℓ^1}) up to the machine precision. This last 64 computation is the only one that requires vectors of the high dimensional space. Our 65 main assumption throughout this paper is that $\exists u$ such that Au = b, i.e., (P_{ℓ^1}) has at 66 least one solution. This can be guaranteed if one assumes, as we shall do hereinafter, 67 that A has full rank. 68

69

Outline of the paper. The paper is organized as follows. Section 2 gives a very 70compact, yet self-contained, presentation of the numerical computations needed to 71implement the algorithm proposed in this paper (see algorithm 2.1 on page 4). Sec-73 tion 3 on page 5 proves the mathematical validity of this algorithm. In other words, we shall prove that the solution computed by algorithm 2.1 is exact (and numerically, 74 up to the machine precision). The convergence (in finite time) of algorithm 2.1 to a solution to (P_{ℓ^1}) is mathematically guaranteed. Section 4 on page 9 proposes a numer-76 ical evaluation and comparison of algorithm 2.1 with some state-of-the-art solutions 77 78 solving (P_{ℓ^1}) . We show in this section that our method has a higher probability of success to reconstruct solutions with high precision compared to other state-of-the-art 79 80 methods, i.e., the phase transition is observed with a high precision. Discussions and conclusions are summarized in section 5 on page 14. The appendix 6 on page 16 con-81 tains several proofs used throughout this paper. A glossary containing the notations 82 and basic definitions is in appendix 7 on page 24. In the sequel, Latin numerals refer 83 to the glossary of notations on page 24. Appendix 8 on page 25 contains general 84

85 results on convex analysis used in this paper.

2. An Algorithm Solving (P_{ℓ^1}) . This section presents the algorithm proposed 86 in this paper. As usual in the literature on compressive sensing, we shall assume that 87 $A \in \mathcal{M}_{m,n}(\mathbb{R})$ with $m \ll n$. The algorithm we shall develop in this paper begins by 88 computing a solution to the dual problem associated to (P_{ℓ^1}) then computes a solution 89 to the primal. The first step involves the computation of a finite and piecewise affine 90 trajectory, or more precisely the positions λ_k where the trajectory changes of slope. 91 The second step relies on the computation of a solution to a constrained least square problem. The construction leads to algorithm 2.1 (page 4). 93 Consider the Lagrangian $\mathcal{L}: \mathbb{R}^n \times \mathbb{R}^m \to \mathbb{R}$ of (P_{ℓ^1}) namely 94

95 (2.1)
$$\mathcal{L}(\boldsymbol{u},\boldsymbol{\lambda}) := J(\boldsymbol{u}) + \langle \boldsymbol{\lambda}, A\boldsymbol{u} \rangle + \langle \boldsymbol{\lambda}, -\boldsymbol{b} \rangle$$

96 where $J(\cdot) = \|\cdot\|_{\ell^1}$. Consider also the function $g: \mathbb{R}^m \to \mathbb{R} \cup \{+\infty\}$ defined by

97
$$g(\boldsymbol{\lambda}) := -\inf_{\boldsymbol{u} \in \mathbb{R}^n} \mathcal{L}(\boldsymbol{u}, \boldsymbol{\lambda}) = -\inf_{\boldsymbol{u} \in \mathbb{R}^n} \left\{ J(\boldsymbol{u}) - \left\langle -A^T \boldsymbol{\lambda}, \boldsymbol{u} \right\rangle \right\} - \left\langle \boldsymbol{\lambda}, -\boldsymbol{b} \right\rangle$$

98 (2.2)
$$= J^* \left(-A^T \boldsymbol{\lambda} \right) + \langle \boldsymbol{\lambda}, \boldsymbol{b} \rangle = \chi_{\mathbf{B}_{\infty}} \left(-A^T \boldsymbol{\lambda} \right) + \langle \boldsymbol{\lambda}, \boldsymbol{b} \rangle$$

where J^* denotes the Lengendre-Fenchel transform of J (see (**xvi**)) and $\chi_{\mathbf{B}_{\infty}}$ denotes the convex characteristic function of ℓ^{∞} (see (**vii**)) unit ball $B_{\infty} \subset \mathbb{R}^n$ (see (**xi**)). (We recall that hereinafter Latin numerals refer to the glossary of notations on page 24.)

102 Consider further the optimization problem

103
$$(D_{\ell^1})$$
 $\inf_{\boldsymbol{\lambda} \in \mathbb{R}^m} g(\boldsymbol{\lambda})$

where g is given by (2.2). As we shall see, under classic assumptions problems (P_{ℓ^1}) and (D_{ℓ^1}) have at least one solution (see proposition 3.4 on page 5). We now give a strategy to solve (D_{ℓ^1}) . The trajectory $[0, +\infty) \ni t \mapsto \lambda(t)$ explicitly given, for every $t \ge 0$, by

108 (2.3)
$$\begin{cases} \frac{\mathrm{d}^{+}\boldsymbol{\lambda}}{\mathrm{dt}}(t) = -\Pi_{\partial g(\boldsymbol{\lambda}(t))}(\mathbf{0}) \\ \boldsymbol{\lambda}(0) = \boldsymbol{\lambda}_{0} \end{cases}$$

converges for some finite time $t_K \in [0, +\infty)$ to a solution to (D_{ℓ^1}) . The main idea 109 of (2.3) is that it generalizes the usual steepest Euclidean descent for non-smooth con-110 vex functions. When the function is not differentiable, then (2.3) selects the smallest 111 velocity in the ℓ^2 sense among all possible velocities that corresponds to the subdif-112 ferential of the function at a non-differentiable point. Note that the subdifferential 113 always only contains one element, which is the gradient, when the function is differen-114 tiable. Formula (2.3) formalizes an evolution equation governed by the (multi-valued) 115maximal monotone operator ∂g (see, for instance, [1, Eq. 2, p. 158]). In (2.3), 116 $\Pi_{\partial q(\boldsymbol{\lambda}(t))}$ denotes the Euclidean projection (xviii) on $\partial g(\boldsymbol{\lambda}(t))$ and $\boldsymbol{\lambda}_0 \in \text{dom } g$ is 117 some initial state. We always set $\lambda_0 = 0$ in our experiments. For any $\lambda \in \text{dom } g$ the 118119multi-valued monotone operator ∂g is given by the non-empty convex cone

120 (2.4)
$$\partial g(\boldsymbol{\lambda}) = \left\{ \boldsymbol{b} + \sum_{i \in \mathrm{S}(\boldsymbol{\lambda})} \eta_i A \tilde{\boldsymbol{e}}_i : \eta_i \ge 0, i \in \mathrm{S}(\boldsymbol{\lambda}) \right\},$$

121 where the set $S(\lambda)$ is defined by

122 (2.5)
$$S(\boldsymbol{\lambda}) := \{i \in \{1, \dots, 2n\} : \langle \boldsymbol{\lambda}, A\tilde{\boldsymbol{e}}_i \rangle = 1\}$$
 and $\tilde{\boldsymbol{e}}_i = \begin{cases} \boldsymbol{e}_i & \text{for } i = \{1, \dots, n\} \\ -\boldsymbol{e}_{i-n} & \text{for } i = \{n+1, \dots, 2n\} \end{cases}$

123 In (2.5) and everywhere else, e_i denotes the *i*-th canonical vector of \mathbb{R}^n .

In addition, the trajectory given by (2.3) is piecewise affine. This means that the next iterate λ_{k+1} produced by the algorithm is computed from the current iterate λ_k , the scalar $(t_{k+1} - t_k)$ and the direction $d_k = -\prod_{\partial g(\lambda_k)}(\mathbf{0})$. We now detail the computation of the scalar $(t_{k+1} - t_k)$. For any $k \in \mathbb{N}$, we define

128 (2.6)
$$S^+(d_k) := \{i \in \{1, \dots, 2n\} : \langle d_k, A\tilde{e}_i \rangle > 0\}$$

129 (2.7) and we have
$$\overline{\Delta t}_k := (t_{k+1} - t_k) = \min\left\{\frac{1 - \langle A\tilde{\boldsymbol{e}}_i, \boldsymbol{\lambda}_k \rangle}{\langle A\tilde{\boldsymbol{e}}_i, \boldsymbol{d}_k \rangle}, \ i \in \mathrm{S}^+(\boldsymbol{d}_k)\right\}.$$

Note that (2.6) and (2.7) are easy to compute since these quantities are given 130explicitly and only involves computations of inner products. Therefore, from (2.3) we 131observe that it remains to compute the direction $d_k = -\prod_{\partial q(\lambda_k)}(\mathbf{0})$ which corresponds 132133to computing the projection on a non-empty closed convex cone given by $\partial g(\boldsymbol{\lambda}_k)$. Note that this subdifferential has an explicit formula given by (2.4). One can use a 134constrained least square solver, available in Matlab, to compute the solution. (See also 135 remark 2.1 below.) To sum up, to compute a solution to (D_{ℓ^1}) one can compute the 136limit of the trajectory $\lambda(t)$ given by (2.3) using the update rules (2.6) and (2.7). This 137138 limit is attained after finitely many updates (see also proposition 3.15). It remains to compute a solution to (P_{ℓ^1}) given $\bar{\lambda}$ solution to (D_{ℓ^1}) . 139

Given λ solution to (D_{ℓ^1}) , one can compute a solution $\bar{\boldsymbol{u}}$ to (P_{ℓ^1}) by solving the constrained least square problem

142 (2.8)
$$\begin{cases} \min_{\boldsymbol{u}\in\mathbb{R}^n} & \|A\boldsymbol{u}-\boldsymbol{b}\|_{\ell^2} \\ \text{s.t.} & u_i \ge 0 \text{ if } \langle \bar{\boldsymbol{\lambda}}, A\boldsymbol{e}_i \rangle = -1, u_i \le 0 \text{ if } \langle \bar{\boldsymbol{\lambda}}, A\boldsymbol{e}_i \rangle = 1 \text{ and } u_i = 0 \text{ otherwise.} \end{cases}$$

143 We are now in position to state the entire algorithm.

144

Algorithm 2.1 Algorithm computing $\bar{\boldsymbol{u}}$ solution to (P_{ℓ^1}) .Input: Matrix A, bOutput: $\bar{\boldsymbol{u}}$ solution to (P_{ℓ^1}) Set k := 0 and $\lambda_k := \boldsymbol{0} \in \mathbb{R}^m$ repeat1. Compute $S(\lambda_k)$ (see (2.5));2. Compute d_k as $d_k := -\Pi_{\partial g(\lambda_k)}(\boldsymbol{0})$ (see remark 2.1)3. Compute $S^+(d_k)$ (see (2.6)) then Δt_k (see (2.7))4. Set $\lambda_{k+1} := \lambda_k + \Delta t_k d_k$;5. Set k = k + 1 and set $\boldsymbol{\lambda} := \frac{\boldsymbol{\lambda}}{\|A^T \boldsymbol{\lambda}\|_{\ell^{\infty}}}$ if $\|A^T \boldsymbol{\lambda}\|_{\ell^{\infty}} > 1$;until $d_k = \boldsymbol{0}$ (see remark 2.1);Compute $\bar{\boldsymbol{u}}$ using (2.8).

145 Remark 2.1. To compute d_k we define $G := \left\{ \sum_{i \in S(\lambda_k)} \eta_i A \tilde{e}_i : \eta_i \ge 0, i \in S(\lambda_k) \right\}$. We

- 146 have that $d_k := -\Pi_{\partial g(\lambda_k)}(\mathbf{0}) = -\Pi_G(-\mathbf{b}) \mathbf{b}$ (see lemma 6.5 on page 21) can be com-
- 147 puted from a constrained least square problem similar to (2.8). We refer to [15,

149 precision) numerical algorithms solving the above constrained least square problem.

150 For instance, one can use the *lsqnonneg* Matlab routine albeit we used an implementa-

151 tion based on [31] that is supposedly faster than the Matlab routine. The termination

152 condition, namely $d_k = 0$, was replaced by $||d_k||_{\ell^2} \vee ||\overline{\Delta t_k} d_k||_{\ell^2} < 10^{-10}$ in all of our 153 experiments. The projection in step 5 is unnecessary if the precision of numbers is

high enough. However, we empirically observed that it increased the performance of

155 the method for the Matlab implementation.

156 **3. From Maximal Monotone Operator to** ℓ^1 **Solutions of Linear Prob-**157 **lems.** This section justifies the mathematical validity of algorithm 2.1 presented in 158 section 2.

We recall that to solve (P_{ℓ^1}) , we first solve the dual (D_{ℓ^1}) then compute a solution 159to the primal problem (P_{ℓ^1}) . Hence, we first give the assumptions that justify the 160 existence of solutions to problems (P_{ℓ^1}) and (D_{ℓ^1}) and give a closed formula that 161 allows us to compute the solution to (P_{ℓ^1}) from a solution to (D_{ℓ^1}) . This is done in 162proposition 3.4. We then briefly justify the fact that the trajectory we used in the 163 previous section converges to a solution to the dual. This is done in proposition 3.6. 164This proposition translates into algorithm 2.1 on page 4 and is illustrated numerically 165 in section 4 on page 9. 166

167 PROPOSITION 3.1 (and definition). We assume that $A \in \mathcal{M}_{m,n}(\mathbb{R})$ has full row 168 rank and that $J(\cdot) = \|\cdot\|_{\ell^1}$. We consider the functions

169 (3.1)
$$\forall \boldsymbol{u} \in \mathbb{R}^n, \quad f(\boldsymbol{u}) := J(\boldsymbol{u}) + \chi_{\{\boldsymbol{b}\}}(A\boldsymbol{u});$$

170 (3.2)
$$\forall \boldsymbol{\lambda} \in \mathbb{R}^m, \quad g(\boldsymbol{\lambda}) := J^* (-A^T \boldsymbol{\lambda}) + \langle \boldsymbol{\lambda}, \boldsymbol{b} \rangle = \chi_{\mathbf{B}_{\infty}} (-A^T \boldsymbol{\lambda}) + \langle \boldsymbol{\lambda}, \boldsymbol{b} \rangle.$$

- 171 We have $f \in \Gamma_0(\mathbb{R}^n)$ and $g \in \Gamma_0(\mathbb{R}^m)$ (see (\mathbf{x})).
- 172 *Proof.* See appendix 16 on page 18.

173 Remark 3.2. The assumptions of proposition 3.1 allow to cover the case of com-174 pressive sensing problems. Note that one can relax the assumption that A is full row 175 rank, we just need that $\mathbf{b} \in \text{span } A$. For instance if, for some specific application, the 176 observed \mathbf{b} 's belong to a subspace B then we just need span $A \supset B$.

We recall that we wish to solve (P_{ℓ^1}) using a solution to (D_{ℓ^1}) . To this aim the following definition and proposition are needed.

179 DEFINITION 3.3 (Active set). For any $\lambda \in dom \ g \ we \ define$

180 (3.3)
$$S(\lambda) := \{i \in \{1, ..., 2n\} : \langle \lambda, A\tilde{e}_i \rangle = 1\}$$
 and $\tilde{e}_i = \begin{cases} e_i & \text{for } i = \{1, ..., n\} \\ -e_{i-n} & \text{for } i = \{n+1, ..., 2n\}, \end{cases}$

181 and e_i denotes the *i*-th canonical vector of \mathbb{R}^n .

- 182 PROPOSITION 3.4 (Existence of solutions and computation of a solution to (P_{ℓ^1})). 183 We posit the same assumptions as in proposition 3.1. We have
- 184 1. Problems (P_{ℓ^1}) and (D_{ℓ^1}) have at least one solution.
- 185 2. Let $\bar{\lambda}$ be a solution to (D_{ℓ^1}) . Consider the coefficients \tilde{u}_i such that $\tilde{u}_i = 0 \ \forall i \in$
- 186 $\{1,\ldots,2n\}\setminus S(-\bar{\lambda}) \text{ and } \tilde{u}_i \ge 0 \text{ for } i \in S(-\bar{\lambda}), \text{ of the Euclidean projection of}$

1

187 **b** onto

188 (3.4)
$$\left\{ \boldsymbol{y} : \boldsymbol{y} = \sum_{\substack{\tilde{u}_i \ge 0 \ \forall i \in \mathrm{S}(-\bar{\boldsymbol{\lambda}}) \\ \tilde{u}_i = 0 \ \forall i \in \{1, \dots, 2n\} \setminus \mathrm{S}(-\bar{\boldsymbol{\lambda}})} \tilde{u}_i A \tilde{\boldsymbol{e}}_i \right\},$$

189 190

where
$$S(-\lambda)$$
 and e_i are defined by (3.3).
We have that the vector \bar{u} obtained from the above coefficients \tilde{u}_i

191 (3.5)
$$\bar{\boldsymbol{u}} := \sum_{i=1}^{n} u_i \boldsymbol{e}_i, \text{ with } u_i := \tilde{u}_i - \tilde{u}_{i+n}$$

192 is a solution to (P_{ℓ^1}) .

Note that (3.4) is equivalent to formula (2.8) given in section 2. Indeed, in (3.4) at least one of the coefficients \tilde{u}_i or \tilde{u}_{i+n} is zero.

1 1

(0,0)

195 Proof. See appendix 17 on page 18.

 $\alpha(\bar{\mathbf{x}})$

Remark 3.5. The reconstruction formula given by (2.8) is different from the re-196construction methods that can sometimes be found in the literature (see, e.g., [32, al-197 gorithm 6, p. 11]). However, for matrices satisfying compressive sensing assumptions 198(see, e.g., [12, 16]), the signal can be obtained from an unconstrained least-squares 199 solution to Au = b. Indeed, the support constraint issued form λ boils down to 200 201 solving, in the least squares sense, $B\boldsymbol{u} = \boldsymbol{b}$, where B is a sub-matrix formed from A by removing appropriate columns. Note that in this case, there is no sign constraint 202on u_i contrarily to (2.8). In addition, in many cases, the unconstrained least squares 203 solution can be computed using a Moore-Penrose pseudo inverse formula. However, 204 the least squares solution and (2.8) will, in general, differ: they have same ℓ^0 pseudo 205norms but different ℓ^1 norms. 206

To solve (D_{ℓ^1}) we rely on a specific trajectory of feasible points for (D_{ℓ^1}) governed by the maximal monotone operator ∂g (see, e.g., [1]). The main properties of this trajectory are summarized in the next proposition.

210 PROPOSITION 3.6 (Properties of the trajectory $\lambda(t)$ [1, 5]). We posit the same 211 assumptions as in proposition 3.1. Consider the evolution equation explicitly given, 212 for every $t \in [0, +\infty)$, by

213 (3.6)
$$\begin{cases} \frac{\mathrm{d}^+ \boldsymbol{\lambda}(t)}{\mathrm{d}t} = -\Pi_{\partial g(\boldsymbol{\lambda}(t))}(\mathbf{0}) \\ \boldsymbol{\lambda}(0) = \boldsymbol{\lambda}_0, \end{cases}$$

where $\lambda_0 \in dom \ \partial g$. We have that the solution $\lambda : [0, +\infty) \ni t \mapsto \lambda(t) \in \mathbb{R}^m$ to (3.6) satisfies:

- 1. for every $t \in [0, +\infty)$, $\lambda(\cdot)$ is continuous, right-differentiable and belongs to dom ∂g ;
- 218 2. the limits of $g(\lambda(t))$ and $\lambda(t)$ when $t \to +\infty$ exist;
- 219 3. $\lim_{t\to+\infty} g(\boldsymbol{\lambda}(t)) = \min_{\boldsymbol{\lambda}\in\mathbb{R}^m} g(\boldsymbol{\lambda}) \text{ and } \lim_{t\to+\infty} \boldsymbol{\lambda}(t) \in \arg\min_{\boldsymbol{\lambda}\in\mathbb{R}^m} g(\boldsymbol{\lambda}).$
- 220 *Proof.* See appendix 18 on page 18.

The proposition above means that the limit of the trajectory $\lambda(t)$ is a solution to (D_{ℓ^1}) .

In the sequel, we shall prove that the limit is attained for a finite time $t \ge 0$. It is

worth noticing the similarity between (3.6) and inverse scale space methods (see, e.g., [7, 32]). To compute $\lambda(t)$ one could rely on an Euler scheme to approximate the trajectory for instance. However, a numerical computation of the trajectory $\lambda(t)$ up to the machine precision is doable. This is the goal of the next paragraph.

Computation of the trajectory $\lambda(t)$ given by (3.6). We recall that to obtain 227 an algorithm we need to compute a solution λ to (D_{ℓ^1}) . To do so, we recall that we 228compute the positions where $\lambda(t_k)$ changes of slope. Since dom $g \neq \mathbb{R}^m$ we cannot 229 recur to classic textbooks such as, e.g., [22, Chap. VIII]. Thus, some work is needed. 230Proposition 3.13 (on page 8) proves that $\lambda(t)$ defined by (3.6) is piecewise affine. 231 In other words, $\lambda(t)$ is made of pieces of straight lines. Hence, the computation of 232 $\lambda(t)$ boils down to the detection of "kicks", i.e., positions where $\lambda(t)$ changes slope 233234 and the computation of these slopes. The computation of these slopes is obtained from (3.6) and lemma 3.8. Propositions 3.10 and 3.14 yield a direct and optimal 235 numerical method to detect kick times, i.e., times t such that $\lambda(t)$ and $\lambda(t+\varepsilon)$ don't 236 have the same slope for some $\varepsilon > 0$. Propositions 3.11 and 3.15 give the termination 237238 condition and proves that $\lambda(t)$ converges to a solution to (D_{ℓ_1}) after finitely many kicks. We recall that proposition 3.4 (page 5) directly gives an explicit formula that 240 allows us to compute a solution to (P_{ℓ^1}) given a solution to (D_{ℓ^1}) obtained as the limit of the trajectory $\lambda(t)$. 241

We recall that one of the two main ingredients to compute the trajectory $\lambda(t)$ is the computation of slopes given by a projection onto the closed convex cone $\partial g(\lambda(t))$ (see proposition 3.6 on page 6). Hence, a closed formula for ∂g is needed. This is the goal of the next proposition that leads to lemma 3.8.

PROPOSITION 3.7 (The function g defined by (3.2) is polyhedral). We posit the same assumptions as in proposition 3.1. The function g defined in (3.2) is a polyhedral proper and convex function that satisfies dom $g = C \neq \emptyset$ and we have

249 (3.7)
$$g(\boldsymbol{\lambda}) = \langle \boldsymbol{\lambda}, \boldsymbol{b} \rangle + \chi_{\mathbf{C}}(\boldsymbol{\lambda}), \text{ where } C := \{ \boldsymbol{\lambda} \in \mathbb{R}^m : \langle \boldsymbol{\lambda}, A\tilde{\boldsymbol{e}}_i \rangle \leq 1, i \in \{1, \dots, 2n\} \}$$

- 250 and $\tilde{\boldsymbol{e}}_i$ is defined in (3.3).
- 251 *Proof.* See appendix 6.5 on page 19.
- 252 We now give a formula for the subdifferential of g.

LEMMA 3.8 (Subdifferential formula for g). We posit the same assumptions as in proposition 3.1. We have dom $\partial q = dom q = C \neq \emptyset$ and, for any $\lambda \in C$,

- 255 (3.8) $\partial g(\boldsymbol{\lambda}) = \{\boldsymbol{b}\} + \operatorname{co}\{A\tilde{\boldsymbol{e}}_i : i \in \mathrm{S}(\boldsymbol{\lambda})\},\$
- 256 where $\tilde{\boldsymbol{e}}_i$, S($\boldsymbol{\lambda}$) are given by (3.3) and co by (\mathbf{v}).
- 257 Proof. See appendix 6.6 on page 19.

With the above formula it is easily seen that one can compute the slope of $\lambda(t)$ for any $t \ge 0$. It remains to compute the kick times, i.e., times t when the slope of the trajectory $\lambda(t)$ changes. This is the goal of the next three propositions and lemma.

261 PROPOSITION 3.9 (and definition: descent direction). We posit the same setup 262 as in proposition 3.1. We say that a direction $\mathbf{d} \in \mathbb{R}^m \setminus \{\mathbf{0}\}$ is a descent direction for g263 at $\boldsymbol{\lambda} \in \text{dom } g \text{ iff } (\boldsymbol{\lambda} + t\mathbf{d}) \in \text{dom } g \text{ and } g(\boldsymbol{\lambda} + t\mathbf{d}) < g(\boldsymbol{\lambda}) \text{ for some } t > 0$. Moreover, 264 we have that a direction $\mathbf{d} \neq \mathbf{0}$ is a descent direction for g at $\boldsymbol{\lambda}$ iff \mathbf{d} satisfies

- 265 (3.9) $\langle \boldsymbol{d}, A \tilde{\boldsymbol{e}}_i \rangle \leq 0 \ \forall i \in S(\boldsymbol{\lambda}) \ and$
- 266 (3.10) $g'(\boldsymbol{\lambda}, \boldsymbol{d}) = \langle \boldsymbol{d}, \boldsymbol{b} \rangle < 0, \text{ where } \tilde{\boldsymbol{e}}_i \text{ is given by (3.3).}$

267 Proof. See appendix 21 on page 19.

PROPOSITION 3.10 (Kick time computation). We posit the same assumptions as in proposition 3.1 and further assume that $\lambda \in \text{dom } g$ and that d is a direction that satisfies (3.9). Consider \tilde{e}_i given by (3.3), the set $S^+(d)$ defined by

271 (3.11)
$$\mathbf{S}^+(\boldsymbol{d}) := \{i \in \{1, \dots, 2n\} : \langle A\tilde{\boldsymbol{e}}_i, \boldsymbol{d} \rangle > 0\}$$

272 and the scalar $\overline{\Delta t}$ by

273 (3.12)
$$\begin{cases} \overline{\Delta t} := \min\left\{\frac{1 - \langle A\tilde{\boldsymbol{e}}_i, \boldsymbol{\lambda} \rangle}{\langle A\tilde{\boldsymbol{e}}_i, \boldsymbol{d} \rangle} : i \in \mathbf{S}^+(\boldsymbol{d}) \right\} & \text{if } \mathbf{S}^+(\boldsymbol{d}) \neq \emptyset;\\ \overline{\Delta t} := +\infty & \text{otherwise.} \end{cases}$$

We have that $\overline{\Delta t}$ satisfies $\overline{\Delta t} > 0$. In addition, $S^+(d) = \emptyset$ iff $(\lambda + td) \in dom g$ for every $t \ge 0$. Furthermore, we have

276 (3.13)
$$(\boldsymbol{\lambda} + t\boldsymbol{d}) \in dom \ g \ iff \ t \in [0, \overline{\Delta t}];$$

277 (3.14) $\forall t \in [0, \overline{\Delta t}) \quad S(\boldsymbol{\lambda} + t\boldsymbol{d}) \subset S(\boldsymbol{\lambda}) \quad and \quad \partial g(\boldsymbol{\lambda} + t\boldsymbol{d}) \subset \partial g(\boldsymbol{\lambda}).$

278 *Proof.* See appendix 23 on page 20.

279 LEMMA 3.11 (Well possedness of $d := -\prod_{\partial g(\lambda)}(\mathbf{0})$, optimality conditions). We 280 posit the same assumptions as in proposition 3.1. For any $\lambda \in \text{dom } g$, the vector given 281 by

282 (3.15)
$$\boldsymbol{d} := -\Pi_{\partial g(\boldsymbol{\lambda})}(\boldsymbol{0}).$$

is well defined. Consider **d** defined by (3.15) and $\overline{\Delta t}$, $S^+(d)$ defined in proposition 3.10. We have that the three following conditions are equivalent

285 (3.16)
$$d = 0 \Leftrightarrow \overline{\Delta t} = +\infty \Leftrightarrow S^+(d) = \emptyset.$$

In addition, $\lambda \in \text{dom } g$ is a solution to (D_{ℓ^1}) iff the conditions in (3.16) hold true.

287 Proof. See appendix 24 on page 8

PROPOSITION 3.12 $(\prod_{\partial g(\lambda)}(\mathbf{0})$ is constant on time intervals). We posit the same assumptions as in proposition 3.1. Consider any $\lambda \in \text{dom } g$, d defined by (3.15) and $\overline{\Delta t}$ defined in corollary 3.10. We have

291 (3.17)
$$\forall t \in \left[0, \overline{\Delta t}\right) \qquad \Pi_{\partial g(\boldsymbol{\lambda})}(\boldsymbol{0}) \in \partial g(\boldsymbol{\lambda} + t\boldsymbol{d});$$

292 (3.18)
$$\forall t \in [0, \overline{\Delta t}) \qquad \Pi_{\partial g(\lambda)}(\mathbf{0}) = \Pi_{\partial g(\lambda + td)}(\mathbf{0}).$$

293 Proof. See appendix 25 on page 21.

We are now in position to give a mathematical definition of the trajectory computed by the algorithm.

PROPOSITION 3.13 (and definition: piecewise affine trajectory $\lambda(t)$). We posit the same assumptions as in proposition 3.1. Consider $\lambda_0 \in \text{dom } g$ and the sequences $(t_k)_k \subset [0, +\infty], (d_k)_k \text{ and } (\lambda(t_k))_k$ recursively defined by

299 (3.19)
$$\begin{cases} t_0 := 0; \quad \boldsymbol{d}_k := -\Pi_{\partial g(\boldsymbol{\lambda}(t_k))}(\boldsymbol{0}); \quad t_{k+1} := t_k + \overline{\Delta} t_k; \\ \boldsymbol{\lambda}(t_{k+1}) := \boldsymbol{\lambda}(t_k) + (t_{k+1} - t_k) \boldsymbol{d}_k \quad if \ t_{k+1} < +\infty \\ \boldsymbol{\lambda}(t_{k+1}) := \boldsymbol{\lambda}(t_k) \quad otherwise, \end{cases}$$

where Δt_k is obtained from proposition 3.10 (applied with $\boldsymbol{\lambda} := \boldsymbol{\lambda}(t_k)$ and $\boldsymbol{d} := \boldsymbol{d}_k$). Consider also the affine interpolate (continuous) trajectory $\boldsymbol{\lambda} : [0, +\infty] \ni t \mapsto \mathbb{R}^m$ defined by

303 (3.20)
$$\boldsymbol{\lambda}(t) := \boldsymbol{\lambda}(t_k) + (t - t_k)\boldsymbol{d}_k \text{ for any } t \in [t_k, t_{k+1}), \qquad \boldsymbol{\lambda}(t_0) := \boldsymbol{\lambda}_0.$$

We have that the trajectory $\lambda(t)$ given in (3.20) coincides for every $t \ge 0$ with the solution to the evolution equation (3.6). In addition, for every $t \ge 0$ we have $\lambda(t) \in$ dom g.

307 Proof. See appendix 29 on page 23.

To compute $\lambda(t)$ the algorithm relies on the computation of the sequence $(d_k, t_k)_k$ defined by (3.19). The next two propositions prove that $\lambda(t)$ changes of slope at every t_k and that the sequences in (3.19) are finite.

311 PROPOSITION 3.14 (Optimality of the sampling of the trajectory $\lambda(t)$). We posit 312 the same assumptions as in proposition 3.1 and further assume that $\lambda(t_k) \in \text{dom } g$ is 313 not a solution to (D_{ℓ^1}) . For $\lambda(t_{k+1})$ given by proposition 3.13 we have

314 (3.21)
$$\Pi_{\partial g(\boldsymbol{\lambda}(t_k))}(\mathbf{0}) \neq \Pi_{\partial g(\boldsymbol{\lambda}(t_{k+1}))}(\mathbf{0}) \text{ and } \|\Pi_{\partial g(\boldsymbol{\lambda}(t_{k+1}))}(\mathbf{0})\|_{\ell^2} < \|\Pi_{\partial g(\boldsymbol{\lambda}(t_k))}(\mathbf{0})\|_{\ell^2}.$$

315 *Proof.* See appendix 30 on page 23.

PROPOSITION 3.15 ($\lambda(t)$ converges to a minimizer of (D_{ℓ^1}) after finitely many kicks). We posit the same assumptions as in proposition 3.1. Consider the sequences $(t_k)_k, (d_k)_k$ and the trajectory $\lambda(t)$ defined in proposition 3.13. We have that $\exists K \in \mathbb{N}$ such that $\lambda(t) = \lambda(t_K)$ for every $t \ge t_K$. In addition, $\lambda(t_K)$ is a solution to (D_{ℓ^1}) and d_K satisfies $d_K = 0$.

321 *Proof.* See appendix 31 on page 24.

We now briefly justify that the computations in algorithm 2.1 (page 4) end with a 322 solution to (P_{ℓ^1}) after finitely many iterations. We obtained that for any $\lambda_0 \in \text{dom } g$ 323 (see proposition 3.13) the sequence defined in (3.19) converges (see proposition 3.15) 324 after finitely many kicks to a solution to (D_{ℓ^1}) . In algorithm 2.1, the initialization step 325 namely $\lambda_0 = 0$ is valid since $0 \in \text{dom } g$. In addition, it is easily seen that steps 1-5 326 implement (3.19). From proposition 3.15, we deduce the validity of the termination 327 condition. Proposition 3.15 justifies that this termination condition is reached after 328 finitely many iterations. Hence, the while loop ends with some $\bar{\lambda}$ solution to (D_{ℓ^1}) . 329 330 Therefore, the computation of $\bar{\boldsymbol{u}}$ solution to (P_{ℓ^1}) is justified by proposition 3.4. Therefore, the validity of algorithm 2.1 is proved. 331

Remark 3.16. Supplementary material shows that our proposed approach can be extended to handle affine inequality constraints. In addition, supplementary material presents how our proposed algorithm 2.1 can be used to solve the optimization problem with constraints of the form $||Au - b||_{\ell^2} \leq \epsilon$, i.e., when there is Gaussian noise. This approach will be presented in another paper.

4. Experiments. This section proposes an empirical evaluation of the following methods to solve (P_{ℓ^1}) : AISS [7], LARS [18], SPGL1 [40, 41], SeDuMI [38] and algorithm 2.1. Two parameters settings are consdered for SeDuMI: the first version version which is called "standard precision" (SP) uses the standard parameters provided in the CVX package, while the second version which is called "high precision" (HP) uses the option "cvx_precision best". Supplementary material gives the same comparisons

between OMP [35], CoSamp [33] and GISS [32]. Note that OMP, CoSamp and GISS 343 are greedy-based numerical algorithms. LARS, SPGL1, AISS and algorithm 2.1 are 344 ℓ^1 -based numerical algorithms. SeDuMi [38] is a toolbox for linear, second order and 345 semi-definite problems. These methods are compared in terms of a "probability of success" (defined below) and average number of iterations needed. The criterion will 347 be used to observe a so called phase transition that separates cases where algorithms 348 successfully recover the sparsest solution and when they fail. Note that solutions with 349 high precision are required to observe an accurate phase transition because, if the pre-350 cision of the computed solutions is too poor, then any estimation can be considered as a solution (i.e., a "success" in our experiments). Numerically, it seems to be hard 352 to know a priori the desired precision on the solutions to observe phase transitions. 353 354 Therefore, it is of interest to have numerical methods that can achieve reconstructions with high precision, i.e., up to the machine precision. 355

First we describe the experimental setup. In these experiments the sensing matrix 356 A always has 1000 columns. The entries of A are drawn from i.i.d. realizations of 357 a centered Gaussian distribution. Without loss of generality we may normalize the 358 columns of A to unit Euclidean norm. The number of rows of A, i.e., the dimension 359 360 of the ambient space m, vary in $M := \{50, \ldots, 325\}$ with increments of 25. For each number of rows, we vary the sparsity level s between 5% and 40% with increments 361 of 5% and therefore consider the discrete set $S := \{0.05, \ldots, 0.4\}$. The sparsity level 362 is related to the ℓ^0 norm of **u** by " $\|\mathbf{u}\|_{\ell^0} = \text{round} (s \times m)$ " following [10]. The posi-363 tions of the non-zero entries of \boldsymbol{u} are chosen randomly, with uniform probability. The 364 365 non-zero entries of u are drawn from a uniform distribution on |-1,1|. To do so, for each parameter (i.e., sparsity level s and dimension of ambient space m) we repeated 366 the experiments 1,000 times. The implementations of AISS and SPGL1 we used are 367 the ones given by the authors of [7, 32, 41]. For LARS [18], we used the SPAMS 368 toolbox [26]. The implementation of SeDuMi [38] we used can be found at https: 369 //sedumi.ie.lehigh.edu/sedumi/files/sedumi-downloads/SeDuMi_1_3.zip. Default pa-370 371 rameters have been used for all methods. We now give the criteria used for the numerical comparisons of these numerical algorithms. 372

We need to define the "success" of an algorithm. We choose to define "success" 373 as "the output of an algorithm is equal to the source element \boldsymbol{u} ". This choice can 374be justified by several theoretical works, see, e.g., [11, 12, 16, 17]. This criterion, 375 namely the output is equal to the source element, is chosen for the numerical experi-376 ments proposed thereafter. Note that this criterion seems slightly in favor of methods 377 specifically designed for the compressive sensing method compared to methods that 378 propose to solve (P_{ℓ^1}) . Here, this means that the comparisons are slightly biased in 379 favor of [33, 35]. We also need to deal with the finite numerical precision of computa-380 tions. Thus, we define that a reconstruction is a success if the relative error satisfies 381 $\frac{\|\boldsymbol{u}-\boldsymbol{u}_{est}\|_{\ell^2}}{\|\boldsymbol{u}\|_{\infty}} < \varepsilon, \text{ where } \varepsilon = 10^{-10} \text{ or } \varepsilon = 10^{-4}. \text{ Hence, for any } (m,s) \in M \times S, \text{ the } S, \text{ the } S \in \mathbb{R}^{d}$ 382 empirical probability of success is given by 383

384 (4.1)
$$P_{(m,s)} := \frac{1}{\# \text{ of tests}} \sum_{i} \mathbb{1}_{\left\{\frac{\|\mathbf{u}^{i} - \mathbf{u}_{est}^{i}\|_{\ell^{2}}}{\|\mathbf{u}^{i}\|_{\ell^{2}}} < \varepsilon\right\}}(i),$$

where \boldsymbol{u}_{est}^i (resp. \boldsymbol{u}^i) is the estimated signal (resp. source signal). Each method is tested on the same data by using the same random seed. Note that this type of experimental setup has been used before, for instance in [25].

Remark that another choice for defining "success" could be stated as "the output of an algorithm is a solution to (P_{ℓ^1}) ". However, this criterion would be verified for every output of algorithm 2.1. Indeed, algorithm 2.1 ends with some $\bar{\boldsymbol{u}}$ that numerically verifies an optimality condition associated with (P_{ℓ^1}) . Thus, this choice seems uninformative. Therefore, we have decided to not consider this definition of "success" in this paper. We first consider $\varepsilon = 10^{-10}$. Figure 4.1 depicts the empirical probability of success (4.1) for AISS, LARS, SPGL1, SeDuMi and algorithm 2.1.

We also consider the difference of probability of success between algorithm 2.1 and all other methods that is defined as follows

397 (4.2)
$$D_{(m,s)} := P_{(m,s)}^{\text{algorithm 2.1}} - P_{(m,s)},$$

where $m \in M$, $s \in S$, $P_{(m,s)}^{\text{algo 2.1}}$ (resp. $P_{(m,s)}$) denotes the quantity (4.1) obtained with 398 algorithm 2.1 (resp. AISS, LARS and SPGL1). Note that a positive (negative) value 399 in (4.2) means that algorithm 2.1 achieves a higher (lower) probability of success than 400 the compared algorithm. These differences of probability of success are depicted in 401 figure 4.2. We deduce from figure 4.2 that algorithm 2.1 always achieves a higher 402probability of success than AISS and GISS. We observe that LARS, SeDuMi (stan-403 dard precision) and SPGL1 algorithms do not perform well for $\varepsilon = 10^{-10}$ since the 404 probability of success tends to be low, even for problems with very sparse signals. We 405also observe that both SeDuMi (high precision) and our proposed algorithm produce 406 the best results. Table 4.1 gives the main assumptions on A and **b** for LARS [18], 407 SPGL1 [40, 41], AISS [7], SeDuMI [38] and algorithm 2.1 In this table, we also give 408 the empirical probability that at least x% of signals are successfully reconstructed for 409each method. This statistical indicator is defined as follows 410

411 (4.3)
$$P_{\ge x} = \frac{\#\{(m,s) \in M \times S : P_{(m,s)} \ge x\}}{\#M \cdot \#S}$$

where $P_{(m,s)}$ is defined by (4.1) and # denotes the cardinality of a set. Supplementary material presents numerical results in terms of ℓ^1 -norm for ℓ^1 -based methods namely AISS, LARS, SPGL1, SeDuMi and algorithm 2.1. Up to a probability of 0.95 AISS, SeDuMi (HP) and our algorithm give the same best results. For probability 0.99 SeDuMi and our algorithm give the same best results. For higher probabilities Algorithm 2.1 gives the best results.

Table 4.2 presents the time results for AISS, LARS, SPGL1, SeDuMi and Al-418 gorithm 2.1. All experiments are done using a single core of an Intel Core 10600k. 419We observe that our proposed algorithm is very competitive compared to the state-420 of-the art competitors. Indeed, our proposed algorithm outperforms the competitors 421 for sparsity 5/10% and 50/175 rows while the second best algorithm is AISS. The 422 computational time of our proposed algorithm is similar to AISS for sparsity 15/20%423 and 175/300 rows. For sparsity 25/30% and 175/300 rows AISS performs better that 424 our proposed algorithm. We observe that the runtime of LARS [18], SPGL1 [40, 41], 425SeDuMi (SP) [38] remains close to constant when the sparsity is greater or equals 426 20%: this suggests that for these levels of sparsity LARS [18], SPGL1 [40, 41], Se-427 DuMi [38] computed poor solutions as it has been numerically exhibited previously. 428 429Recall that SeDuMi (HP) [38] computes very good results as previously shown but the computational time is significantly larger that our proposed algorithm 2.1 and 430431 AISS except for the case of 30% sparsity with 300 rows.

432 As noted above the numerical results for LARS and SPGL1 show that these 433 two numerical methods are not able to produce good results for the above set of 434 experiments with $\varepsilon = 10^{-10}$. We now present numerical experiments for a higher 435 threshold in (4.1) where we set $\varepsilon = 10^{-4}$. Figure 4.3 depicts the empirical probability



FIG. 4.1. Empirical probability of success (4.1), with $\varepsilon = 10^{-10}$. Panel (a): AISS [7], panel (b): LARS [18], panel (c): SPGL1 [40, 41], panel (d) : SeDuMi (standard precision) [38], panel (e): algorithm 2.1 and panel (f): SeDuMi (high precision) [38]. The non-zero entries of the source element \mathbf{u} are drawn from a uniform distribution on [-1,1]. The entries in A are drawn from i.i.d. realizations of a Gaussian distribution. With their default parameters LARS, SPGL1 and SeDuMi (standard precision) are not able to produce good result for the above set of experiments. However, SeDuMi (high precision) produces good results. We also present results for an higher threshold $\varepsilon = 10^{-4}$, see figure 4.3.



FIG. 4.2. Differences of probability of success (4.2), with $\varepsilon = 10^{-10}$. panel (a): algorithm 2.1-AISS [7] panel panel (b): algorithm 2.1-LARS [18] panel (c): algorithm 2.1-SPGL1 [40, 41], panel (d) algorithm 2.1-SeDuMi (standard precision) [38] and panel (e) algorithm 2.1-SeDuMi (high precision) [38]. A positive value indicates that algorithm 2.1 achieves a higher probability of success than the considered method, a negative value the contrary.

of success (4.1) for AISS, LARS, SPGL1, SeDuMi and algorithm 2.1. Figure 4.4 436 depicts the differences of probability of success. These results for $\varepsilon = 10^{-4}$ show 437that all numerical algorithms have a higher empirical probability of success compared 438 to the results for $\varepsilon = 10^{-10}$. In particular, we note that SPGL1 and LARS that were performing poorly for $\varepsilon = 10^{-10}$ have dramatically improved their performance. 439

TABLE 4.1

Main assumption and statistical indicator of "success" for LARS, SPGL1, AISS, SeDuMi and algorithm 2.1. The numbers without parentheses correspond to $\varepsilon = 10^{-10}$ and those between parentheses correspond to $\varepsilon = 10^{-4}$. Below, R.I.C. stands for restricted isometry constant see, e.g., [33] and S.F.P.D. stands for strong feasibility of primal and dual program.

Algorithm	LARS [18]	SPGL1 [40, 41]	AISS [7]	SeDuMi (SP) [38]	SeDuMi (HP) [38]	Algorithm 2.1
Assumption	R.I.C.	$\exists u : Au = b$	$\exists u : Au = b$	S.F.P.D.	S.F.P.D.	full row rank
$P_{\geq 0.9}$ (4.3)	0(0.4688)	0 (0.0833)	0.4688 (0.4688)	0.104(0.4688)	0.4688 (0.4688)	0.4688 (0.4688)
$P_{\geq 0.95}$ (4.3)	0(0.4375)	0 (0.0625)	0.4375 (0.4375)	0(0.4583)	0.4375 (0.4375)	0.4375 (0.4375)
$P_{\geq 0.99}$ (4.3)	0 (0.4167)	0 (0.0104)	0.3438(0.4167)	0 (0.4167)	0.4167 (0.4167)	0.4167 (0.4167)
$P_{\geq 0.999}$ (4.3)	0(0.3750)	0 (0)	0.1250(0.3646)	0 (0.3333)	0.3229(0.3437)	0.3646 (0.3750)
$P_{\ge 1}$ (4.3)	0(0.3646)	0 (0)	0.0521(0.3646)	0(0.1875)	0.1562(0.1979)	0.3333 (0.3646)

Also, from figure 4.4, we observe that LARS and our proposed algorithm produce very similar results. It seems that LARS works for the considered experiments (see figure 4.4) although it was proved in [6] that LARS may not converge.



FIG. 4.3. Empirical probability of success (4.1), with $\varepsilon = 10^{-4}$. Panel (a): AISS [7], panel (b): LARS [18], panel (c): SPGL1 [40, 41], panel (d) : SeDuMi (standard precision) [38], panel (e): algorithm 2.1 and panel (f): SeDuMi (high precision) [38]. The non-zero entries of the source element **u** are drawn from a uniform distribution on [-1,1]. The entries in A are drawn from i.i.d. realizations of a Gaussian distribution.

5. Conclusion. In this paper, a new algorithm to solve ℓ^1 regularized linear 444 problems up to the machine precision has been proposed. The method is based on i) 445 the numerical computation of a finite sequence that converges to a solution the dual 446 447 problem and ii) an explicit recovery formula -based on a non-negative least squaresto compute a solution to the primal problem. The sequence we employed is driven 448 449 by an evolution equation ruled by a maximal monotone operator. The numerical computations of this algorithm involve: the computation of a projection onto a closed 450convex cone and the evaluation of inner products. The sequence in the dual space 451 lives in a low dimensional space compared to the unknown. Hence, most of the 452453 numerical efforts require fewer memory usage than primal-based method. Numerical

TABLE 4.2

Computational time results for the following methods: algorithm 2.1, AISS [7] and SeDuMi (SP) [38], SPGL1 [40, 41], LARS [18], SeDuMi (HP) [38], the number of columns is set to 1,000 as everywhere else in this paper and various number of rows (NR) and several level of sparsity. Time results are given in seconds and corresponds to the average time of 200 experiments. The variance is also given in parenthesis.

NR	algorithm	sparsity		
		5%	10%	15%
50	Algorithm 2.1	6.2437e-04 (1.1274e-08)	9.9688e-04 (4.6010e-07)	0.0056 (1.2178e-05)
	AISS $[7]$	$0.0046 \ (6.2449e-04)$	0.0053 (5.8065e-04)	0.0135(5.9706e-04)
	SeDuMi (SP) [38]	0.0326 (2.9972e-04)	0.0369 (2.7969e-04)	0.0497 (3.4206e-04)
	SPGL1 [40, 41]	0.0115 (3.5370e-05)	0.0240 (2.2200e-04)	0.0693 (4.0352e-04)
	LARS [18]	0.0059 (1.2202e-06)	0.0069 (3.5498e-06)	0.0100 (7.6192e-06)
	SeDuMi (HP) [38]	0.2029 (5.4966e-04)	0.2130 (5.9799e-04)	$0.2330\ (0.0010)$
175	Algorithm 2.1	0.0033 (1.1617e-07)	0.0061 (5.7905e-07)	$0.01\overline{61} (2.6142e-05)$
	AISS $[7]$	0.0052 (5.6492e-04)	0.0071 (6.6601e-04)	0.0153 (7.7310e-04)
	SeDuMi (SP) [38]	0.1497 (3.6402e-04)	0.1714 (3.4310e-04)	0.1729 (4.1359e-04)
	SPGL1 [40, 41]	0.0105 (2.3802e-05)	0.0192 (3.6582e-05)	0.0372 (2.3437e-04)
	LARS [18]	0.0097 (7.4096e-06)	0.0159 (2.9454e-05)	$0.0198 \ (2.0914e\text{-}05)$
	SeDuMi (HP) [38]	$0.5813 \ (0.0027)$	$0.6406\ (0.0038)$	$0.6787 \ (0.0028)$
300	Algorithm 2.1	0.0081 (1.1997e-07)	0.0162 (1.9130e-06)	$0.0411 \ (6.8016e-05)$
	AISS $[7]$	0.0063 (6.0184e-04)	0.0102 (5.8937e-04)	0.0283 (8.1150e-04)
	SeDuMi (SP) [38]	0.3502(5.3829e-04)	0.3724 (4.9952e-04)	0.3854 (5.9227e-04)
	SPGL1 [40, 41]	$0.0112 \ (2.1614e-05)$	$0.0181 \ (2.6325e-05)$	0.0299 (5.3353e-05)
	LARS [18]	0.0244 (1.7508e-05)	0.0300 (3.3939e-05)	0.0369 (4.4780e-05)
	SeDuMi (HP) [38]	1.2914(0.0152)	$1.4947 \ (0.0179)$	$1.5797 \ (0.0143)$
NR	algorithm		sparsity	
		20%	25%	30%
50	Algorithm 2.1	0.0085 (3.6050e-06)	0.0091 (1.0906e-06)	0.0092 (7.7791e-07)
	AISS $[7]$	0.0186 (7.7498e-04)	0.0197 (7.6886e-04)	0.0195 (7.6616e-04)
	SeDuMi (SP) [38]	0.0547 (3.6486e-04)	0.0539 (3.0153e-04)	0.0538(3.0346e-04)
	SPGL1 [40, 41]	0.0755 (3.2247e-04)	0.0779 (2.6061e-04)	0.0791 (2.8112e-04)
	LARS [18]	0.0056 (1.2995e-06)	0.0059 (1.2474e-06)	$0.0601 \ (1.2694e-06)$
	SeDuMi (HP) [38]	0.2407 (9.3084e-04)	$0.2496\ (0.0010)$	$0.2537 \ (0.0012)$
175	Algorithm 2.1	0.0997 (0.0165)	0.4943(0.0270)	$0.\overline{5521} \ (0.0028)$
	AISS $[7]$	$0.0866\ (0.0083)$	$0.3485\ (0.0115)$	$0.3739\ (0.0024)$
	SeDuMi (SP) [38]	$0.2077 \ (0.0011)$	$0.2491 \ (0.0013)$	0.2534 (4.9534e-04)
	SPGL1 [40, 41]	0.1187 (0.0026)	$0.1421 \ (0.0015)$	0.1334 (9.7981e-04)
	LARS [18]	0.0199 (1.7296e-05)	0.0173 (3.1513e-06)	0.0209 (1.6831e-05)
	SeDuMi (HP) [38]	0.7787 (0.0170)	.9592(0.0147)	$0.9820 \ (0.0054)$
300	Algorithm 2.1	0.1402(0.0035)	0.7100(0.3449)	3.6610 (2.1706)
	AISS $[7]$	$0.1039\ (0.0025)$	0.4710(0.1126)	$2.0890 \ (0.5683)$
	SeDuMi (SP) [38]	0.4069 (7.2509e-04)	0.4403(0.0021)	$0.5739\ (0.0049)$
	SPGL1 [40, 41]	$0.0591 \ (4.4726e-04)$	0.1639(0.0038)	$0.2013 \ (0.0019)$
	LARS [18]	0.0435 (4.1577e-05)	0.0515 (2.4983e-05)	0.0532 (9.1295e-06)
	SeDuMi (HP) [38]	1.7613(0.0235)	1.8634(0.0752)	2.6824(0.1787)

comparisons with other existing state-of-the-art methods is exhibited for noiseless compressive sensing (basis pursuit) problems.

The numerical comparisons above showed that our algorithm compares advantageously with existing methods: the phase transition is observed with a higher accuracy. The algorithm proposed in this paper is parameter-less once a starting point has been chosen. However, the starting point can be tuned to further speed up the method. A future work could study the impact of this choice in terms of convergence speed.

We also leave as future work theoretical and numerical comparisons with approximate path-methods (as opposed to piecewise affine paths such as our approach) such as [27] which corresponds to an approximate discertization of trajectories. In partic-



FIG. 4.4. Differences of probability of success (4.2), with $\varepsilon = 10^{-4}$. Panel (a): algorithm 2.1-AISS [7], panel (b): algorithm 2.1-LARS [18], panel (c): algorithm 2.1-SPGL1 [40, 41], panel (d) algorithm 2.1-SeDuMi (standard precision) [38] and panel (e) algorithm 2.1-SeDuMi (high precision) [38]. A positive value indicates that algorithm 2.1 achieves a higher probability of success than the considered method, a negative value the contrary.

ular, it would be interest to know if it is better to compute an exact trajectory versusan approximate trajectory from a computational point of view.

6. Appendix. This section contains several proofs used throughout this paper and some properties on the projection on polyhedral convex cone.

468 **6.1.** Some properties of functions J, J^*, f and \mathcal{J} .

469 LEMMA 6.1 (Some elementary properties of J and J^*). We posit the same 470 assumptions as in proposition 3.1. We have

471 1. $J \in \Gamma_0(\mathbb{R}^n)$, $dom(J) = \mathbb{R}^n$, $J^* = \chi_{\mathbf{B}_{\infty}} \in \Gamma_0(\mathbb{R}^n)$ and $dom(J^*) = B_{\infty}$; 472 2. (Primal feasibility)

473
$$(6.1) \quad \mathbf{0} \in int(A \ dom \ J - \{\mathbf{b}\}) = A \ \mathbb{R}^n - \{\mathbf{b}\} = \mathbb{R}^m \ (see \ (iv));$$

474 **3.** (Dual feasibility)

475 (6.2)
$$\mathbf{0} \in int\left(A^T \operatorname{dom} \chi^*_{\{\mathbf{b}\}} + \operatorname{dom} J^*\right) = int\left(\operatorname{span} A^T + B_{\infty}\right)$$

476 *Proof.* We sequentially prove the three assertions.

Note that dom $J = \mathbb{R}^n$ and that J is convex. It follows that $J \in \Gamma_0(\mathbb{R}^n)$ and, from theorem 8.10, that $J^* \in \Gamma_0(\mathbb{R}^n)$. Combining lemma 8.2 with proposition 8.4 we

479 obtain that for any $\boldsymbol{u} \in \mathbb{R}^n$ we have $J^*(\boldsymbol{u}) = \chi_{\mathbf{B}_{\infty}}(\boldsymbol{u})$ and dom $J^* = B_{\infty}$.

From dom $J = \mathbb{R}^n$ and the assumption that A has full row rank, we have A dom J =span $A = \mathbb{R}^m$ and (6.1) immediately follows.

482 Applying lemma 8.3 with $C := \{b\}$ we have $\chi^*_{\{b\}}(\cdot) = \langle \cdot, b \rangle \in \Gamma_0(\mathbb{R}^m)$ and also

483 dom $\left(\chi_{\{b\}}^*(\cdot)\right) = \mathbb{R}^m$. Since, in addition, dom $J^* = B_{\infty}$, we have

$$484 \quad (6.3) \qquad \qquad A^T \operatorname{dom} \chi^*_{\{b\}} + \operatorname{dom} J^* = B_{\infty} + \operatorname{span} A^T.$$

485 We have obviously have $B_{\infty} \subset B_{\infty} + \text{span } A^T$, and from (6.3) we deduce (6.2).

486 PROPOSITION 6.2 (and definition : function \mathcal{J}). We posit the same assumptions 487 as in proposition 3.1. Consider the function $\mathcal{J} : \mathbb{R}^m \to \mathbb{R} \cup \{+\infty\}$ defined by

488 (6.4)
$$\forall \boldsymbol{\lambda} \in \mathbb{R}^{m}, \quad \mathcal{J}(\boldsymbol{\lambda}) := J^{*}(-A^{T}\boldsymbol{\lambda}) = \chi_{\mathbf{C}}(\boldsymbol{\lambda})$$

489 where C is defined by (3.7). We have $\mathcal{J} \in \Gamma_0(\mathbb{R}^m)$ and dom $\mathcal{J} = C \neq \emptyset$.

490 Proof. From item 1 of lemma 6.1 we have $J^* \in \Gamma_0(\mathbb{R}^n)$. Note that (6.2) in 491 lemma 6.1 implies that span $A^T \cap \text{dom } J^* \neq \emptyset$. Then, from theorem 8.5 we obtain 492 that $\mathcal{J}(\cdot) := J^*(-A^T \cdot) \in \Gamma_0(\mathbb{R}^m)$. Moreover, for any $\lambda \in \mathbb{R}^m$, we have

493 (6.5)
$$\mathcal{J}(\boldsymbol{\lambda}) = \chi_{\mathbf{B}_{\infty}} \left(A^T \boldsymbol{\lambda} \right) = \chi_{\mathbf{C}}(\boldsymbol{\lambda}).$$

494 The first equality in (6.5) is justified by combining item 1 $(J^* = \chi_{\mathbf{B}_{\infty}})$ of lemma 6.1 495 and that $-A^T \boldsymbol{\lambda} \in B_{\infty} \Leftrightarrow A^T \boldsymbol{\lambda} \in B_{\infty}$. The second equality in (6.5) is justified by the 496 fact that $\boldsymbol{\lambda} \in C \Leftrightarrow A^T \boldsymbol{\lambda} \in B_{\infty}$. Indeed, we have

497 (6.6)
$$\boldsymbol{\lambda} \in C \Leftrightarrow \langle \boldsymbol{\lambda}, A \tilde{\boldsymbol{e}}_i \rangle \leqslant 1 \; \forall i = \{1, \dots, 2n\} \Leftrightarrow A^T \boldsymbol{\lambda} \in B_{\infty},$$

where $\tilde{\boldsymbol{e}}_i$ is defined by (3.3). The first equivalence in (6.6) is obvious from the definition of *C* given by (3.7). The last equivalence follows from the definition of the $\ell^{\infty}(\mathbb{R}^n)$ unit ball (see (vii)). From (6.5) we can verify that dom $\mathcal{J} = C$ and that $\mathbf{0} \in C \neq \emptyset.\square$

501 LEMMA 6.3 (Subdifferential formulas for f, J^* and \mathcal{J}). We posit the same 502 assumptions as in proposition 3.1. We have

503 1. (Subdifferential formula for f)

504 (6.7)
$$\forall \boldsymbol{u} \in dom \ J \cap dom \ \chi_{\{\boldsymbol{b}\}}(A \cdot) \qquad \partial f(\boldsymbol{u}) = \partial J(\boldsymbol{u}) + A^T \partial \chi_{\{\boldsymbol{b}\}}(A\boldsymbol{u});$$

505 2. (subdifferential formula for J^*)

506 (6.8)
$$\forall \boldsymbol{\lambda} \in dom \ g$$
 $\partial J^* \left(-A^T \boldsymbol{\lambda} \right) = N_{B_{\infty}} \left(-A^T \boldsymbol{\lambda} \right) = \operatorname{co}\{ \tilde{\boldsymbol{e}}_i : i \in \mathrm{S}(-\boldsymbol{\lambda}) \};$

507 3. (Subdifferential formula for \mathcal{J})

$$\forall \boldsymbol{\lambda} \in C \qquad \partial \mathcal{J}(\boldsymbol{\lambda}) = AN_{B_{\infty}}(A^T \boldsymbol{\lambda}) = \operatorname{co}\{A\tilde{\boldsymbol{e}}_i : i \in S(\boldsymbol{\lambda})\},\$$

509 where $N_{B_{\infty}}(A^T \lambda)$ is the normal cone to B_{∞} at $A^T \lambda \in \mathbb{R}^n$ (see (vi)), the 510 set $S(\lambda)$ is defined by (3.3) and the 2n vectors $\tilde{\boldsymbol{e}}_i$ of \mathbb{R}^n are defined by (3.3) 511 (page 5).

Proof. We sequentially justify (6.7)-(6.9). Combining (6.1) in lemma 6.1 and theorem 8.16 (with "U = J" and " $V = \chi_{\{b\}}$ ") we immediately obtain (6.7). The first equality in (6.8) is justified by lemma 8.6. The second equality in (6.8) follows form lemma 8.7 applied with $p := 2n, s_i := e_i$ for $i = 1, \ldots, n, s_i := -e_{i-n}$ for i = $n+1, \ldots, 2n, r_i := 1$ for $i = 1, \ldots, 2n$ and $W(-A^T \lambda) = S(-\lambda)$. We now justify (6.9). From (6.4) in proposition 6.2 we have $\mathcal{J}(\cdot) = J^*(A^T \cdot)$. To prove the first equality in (6.9), we need to justify that

519 (6.10)
$$\operatorname{int}(\operatorname{dom} J^*) \cap \operatorname{span} A^T \neq \emptyset.$$

Assuming that (6.10) holds true, combining item 1 in lemma 6.1 $(J^* \in \Gamma_0(\mathbb{R}^n))$ and theorem 8.9 (with " $f = J^*$ ") we obtain that $\partial \mathcal{J}(\boldsymbol{\lambda}) = -A\partial J^*(-A^T\boldsymbol{\lambda})$. We notice that

523 (6.11)
$$-\partial J^* \left(-A^T \lambda \right) = \operatorname{co} \{ -\tilde{\boldsymbol{e}}, i \in \mathrm{S}(-\lambda) \} = \operatorname{co} \{ \tilde{\boldsymbol{e}}, i \in \mathrm{S}(\lambda) \} = \partial J^* \left(A^T \lambda \right).$$

Indeed, the first equality in (6.11) is justified by (6.8). The second equality is obvious from the definition of $S(\lambda)$ in definition 3.3 and the last equality follows. From (6.11) we immediately obtain (6.9). We now justify (6.10). From, again, item 1 in lemma 6.1 we have dom $J^* = B_{\infty}$ and therefore deduce that $\mathbf{0} \in \operatorname{int}(\operatorname{dom} J^*) \cap \operatorname{span} A^T$ which justifies (6.10). This concludes our proof.

529 6.2. Proof of proposition 3.1 on page 5.

From the assumption that $f \in \Gamma_0(\mathbb{R}^n)$ then that $g \in \Gamma_0(\mathbb{R}^m)$ and that (3.2) hold true. From the assumption that A has a full row rank, it follows that span $A \cap$ dom $\chi_{\{b\}} = \text{span } A \cap \{b\} \neq \emptyset$. In addition, $\chi_{\{b\}} \in \Gamma_0(\mathbb{R}^m)$ as the characteristic function of the closed convex set $\{b\}$. Therefore, from proposition 8.5 we deduce that $\chi_{\{b\}}(A \cdot) \in \Gamma_0(\mathbb{R}^n)$. Combining lemma 6.1 and proposition 8.10 we have that $f \in \Gamma_0(\mathbb{R}^n)$ as the sum of the finite valued convex function $J \in \Gamma_0(\mathbb{R}^n)$ and $\chi_{\{b\}}(A \cdot) \in \Gamma_0(\mathbb{R}^n)$.

Now we prove that $g \in \Gamma_0(\mathbb{R}^m)$. From proposition 6.2 the function $\mathcal{J}(\cdot) = J^*(-A^T \cdot) \in \Gamma_0(\mathbb{R}^m)$. From proposition 8.10 we obtain that $g \in \Gamma_0(\mathbb{R}^m)$ as the sum of the finite valued convex function $\langle \boldsymbol{b}, \cdot \rangle$ and $J^*(-A^T \cdot) \in \Gamma_0(\mathbb{R}^m)$. The second equality in (3.2) follows from item 1 $(J^* = \chi_{\mathbf{B}_m})$ of lemma 6.1. This concludes our proof.

6.3. Proof of proposition 3.4 on page 5.

542 *Proof.* The proof is in two steps. Step 1 proves that problems (P_{ℓ^1}) and (D_{ℓ^1}) 543 have at least one solution. Step 2 justifies (3.4) and (3.5).

544 Step 1. Problems (P_{ℓ^1}) and (D_{ℓ^1}) have at least one solution and . Com-545 bining the definitions of function f and g given in proposition 3.1, lemma 6.1, propo-546 sition 8.17 and theorem 8.18 with U := J and $V := \chi_{\{b\}}$ we conclude that prob-547 lems (P_{ℓ^1}) and (D_{ℓ^1}) have at least one solution. This concludes step1. We now turn 548 to step2.

549 **Step2.** Formulas (3.4) and (3.5) hold true. From [1, p. 166-167] applied 550 with "U := J and $V := \chi_{\{b\}}$ " we have that any point $\bar{\boldsymbol{u}}$ in the non-empty closed 551 convex set $S(\bar{\boldsymbol{\lambda}}) = \partial J^*(-A^T \bar{\boldsymbol{\lambda}}) \cap \{\boldsymbol{u} : A\boldsymbol{u} = \boldsymbol{b}\}$ is a solution to (P_{ℓ^1}) . The set $S(\bar{\boldsymbol{\lambda}})$ is 552 non-empty and from step 1 the primal has a solution. Consider $\bar{\boldsymbol{\lambda}}$ solution to (D_{ℓ^1}) . 553 Combining theorem 8.11 and lemma 3.8 we obtain that $\boldsymbol{b} \in \operatorname{co} \{-A\tilde{\boldsymbol{e}}, i \in S(\bar{\boldsymbol{\lambda}})\} =$ 554 $\operatorname{co} \{A\tilde{\boldsymbol{e}}, i \in S(-\bar{\boldsymbol{\lambda}})\}$. This means that \boldsymbol{b} can be written as $\boldsymbol{b} = \sum_{i=1}^{2n} \tilde{u}_i A \tilde{\boldsymbol{e}}_i$, where 555 $\tilde{u}_i \ge 0, \forall i \in (-\bar{\boldsymbol{\lambda}})$ and $\tilde{u}_i = 0 \forall i \in \{1, \dots, 2n\} \setminus S(-\bar{\boldsymbol{\lambda}})$. Consider $\bar{\boldsymbol{u}}$ defined by (3.5). 556 It is easy to see that $\bar{\boldsymbol{u}} \in S(\bar{\boldsymbol{\lambda}})$ and therefore $\bar{\boldsymbol{u}}$ is a solution to (P_{ℓ^1}) . This concludes 557 our proof.

558 6.4. Proof of proposition 3.6 on page 6.

Proof. From proposition 3.1, we have $g \in \Gamma_0(\mathbb{R}^m)$. Hence, from proposition 8.12 we immediately obtain that $\partial g(\cdot)$ is a maximal monotone operator. Item 1 of proposition 3.6 follows from theorem 8.14. Items 2 and 3 of proposition 3.6 follows from theorem 8.15. This concludes our proof. □ 563 6.5. Proof of proposition 3.7 on page 7.

564 *Proof.* Proposition 3.7 is obvious combining proposition 6.2 and (3.2) in proposi-565 tion 3.1. This concludes our proof.

566 **6.6. Proof of lemma 3.8 on page 7.**

- 567 *Proof.* Combining propositions 3.1 and 6.2, we have that g can be written as
- 568 (6.12) $\forall \boldsymbol{\lambda} \in \mathbb{R}^{m}, \quad g(\boldsymbol{\lambda}) = \mathcal{J}(\boldsymbol{\lambda}) + \langle \boldsymbol{\lambda}, \boldsymbol{b} \rangle.$

From, again, proposition 6.2, we deduce that int (dom \mathcal{J}) \cap int (dom $\langle \cdot, \boldsymbol{b} \rangle$) $\neq \emptyset$.

Hence, combining theorem 8.13 and lemma 6.3 we obtain (3.8) and that $\forall \lambda \in \text{dom } g = C$ we have $\mathbf{b} \in \partial q(\lambda) \neq \emptyset$. This concludes our proof.

572 6.7. Proof of proposition 3.9 on page 7.

573 *Proof.* We first establish the following lemma

574 LEMMA 6.4 (Directional derivative and Taylor formula for g).

We posit the same assumptions as in proposition 3.1. For every $\lambda \in \text{dom } g$ and any that satisfies (3.9) we have

577 (6.13) $q'(\boldsymbol{\lambda}, \boldsymbol{d}) = \langle \boldsymbol{d}, \boldsymbol{b} \rangle;$

578 (6.14) $g(\boldsymbol{\lambda} + t\boldsymbol{d}) = g(\boldsymbol{\lambda}) + tg'(\boldsymbol{\lambda}, \boldsymbol{d})$ for some t > 0 small enough.

579 Proof. Let $\lambda \in \text{dom } g$. By assumption we have that d satisfies (3.9). Combining 580 proposition 3.7 and (3.9) we immediately obtain that $(\lambda + td) \in \text{dom } g$ for some t > 0581 small enough. Hence from the definition of g (3.7), we obtain that, for some small 582 enough t > 0,

583 (6.15)
$$g(\boldsymbol{\lambda} + t\boldsymbol{d}) - g(\boldsymbol{\lambda}) = \langle \boldsymbol{\lambda} + t\boldsymbol{d}, \boldsymbol{b} \rangle - \langle \boldsymbol{d}, \boldsymbol{b} \rangle = t \langle \boldsymbol{d}, \boldsymbol{b} \rangle.$$

Formula (6.13) follows (**xiii**). Combining (6.13) and (6.15) we deduce (6.14). This concludes our proof.

We first prove that the conditions (3.9)-(3.10) are necessary then that they are sufficient.

If d is a descent direction for g at $\lambda \in \text{dom } g$ then, from definition 3.9, $(\lambda + td) \in \text{dom } g$ for some t > 0 small enough. From the definition of C (3.7) it follows that $\lambda + td$ satisfies, in particular, $\langle \lambda + td, A\tilde{e}_i \rangle \leq 1$ for every $i \in S(\lambda)$. The definition of $S(\lambda)$ (3.3) and the fact that $\lambda \in \text{dom } g$ imply that necessarily $\langle d, A\tilde{e}_i \rangle \leq 0$ for every $i \in S(\lambda)$ and (3.9) holds true. In addition, from definition 3.9 we have $g(\lambda + td) < g(\lambda)$ for some t > 0 enough small and combining (6.13)-(6.14) we obtain that (3.10) holds true. Hence, (3.9)-(3.10) are necessary conditions. We now turn to the sufficiency.

Conversely, consider $d \in \mathbb{R}^m \setminus \{0\}$ satisfying (3.9)-(3.10). From proposition 3.7 we have that $\lambda \in \text{dom } g$ satisfies $\langle \lambda, A\tilde{e}_i \rangle \leq 1$ for every $i \in \{1, \ldots, 2n\}$. On the one hand, 596from (3.9), for any $i \in S(\lambda)$ we have $\langle d, A\tilde{e}_i \rangle \leq 0$ and therefore $\langle \lambda + td, A\tilde{e}_i \rangle \leq 1$ $\forall t > 0$. On the other hand, from $\lambda \in \text{dom } g$ and the definition of $S(\lambda)$ we deduce that 598 599 for any $i \in \{1, \ldots, 2n\} \setminus S(\boldsymbol{\lambda})$ we have $\langle \boldsymbol{\lambda}, A \tilde{\boldsymbol{e}}_i \rangle < 1$ and therefore that $\langle \boldsymbol{\lambda} + t \boldsymbol{d}, A \tilde{\boldsymbol{e}}_i \rangle \leq 1$ for t > 0 small enough. Thus, $\langle \boldsymbol{\lambda} + t\boldsymbol{d}, -A\tilde{\boldsymbol{e}}_i \rangle \leq 1$ for every $i \in \{1, \ldots, 2n\}$ and 600 t > 0 small enough. It follows that $(\lambda + td) \in \text{dom } g$ for some t > 0 small enough. 601 Combining (3.10) and (6.14) we obtain $g(\lambda + td) < g(\lambda)$ for some t > 0 small enough. 602 603 It follows that d is a direction descent for q at λ . This concludes our proof. Π

604 6.8. Proof of proposition 3.10 on page 8.

605 Proof. We recall that, in the sequel, $\lambda \in \text{dom } g$ and that d is a descent direction 606 for g at $\lambda \in \text{dom } g$. We sequentially consider the three following complementary cases

607 **Case 1.** The case of indexes *i* such that $i \in S(\lambda)$.

608 **Case 2.** The case of $i \in \{1, \ldots, 2n\} \setminus S(\lambda)$ and $\langle A\tilde{e}_i, d \rangle \leq 0$.

609 **Case 3.** The case of $i \in \{1, \ldots, 2n\} \setminus S(\lambda)$ and $\langle A \tilde{e}_i, d \rangle > 0$.

20

611 **Case 1.** From proposition 3.9 equation (3.9), we have that any descent direction 612 *d* for *g* at $\lambda \in \text{dom } g$ satisfies $\langle d, A\tilde{e}_i \rangle \leq 0$ for every $i \in S(\lambda)$. From definition 3.3 613 (page 5), for every $i \in S(\lambda)$, we have $\langle \lambda, A\tilde{e}_i \rangle = 1$ and therefore deduce that

614 (6.16)
$$i \in S(\lambda) \Rightarrow \langle \lambda + td, A\tilde{e}_i \rangle \leq 1 \quad \forall t \geq 0.$$

615 **Case 2.** For any $i \in \{1, ..., 2n\} \setminus S(\lambda)$, from $\lambda \in \text{dom } g$ we deduce that $\langle \lambda, A\tilde{e}_i \rangle < 1$. 616 Hence, if $\langle A\tilde{e}_i, d \rangle \leq 0$ then

617 (6.17)
$$i \in \{1, \ldots, 2n\} \setminus S(\lambda) \text{ and } \langle A\tilde{e}_i, d \rangle \leq 0 \Rightarrow \langle \lambda + td, A\tilde{e}_i \rangle < 1 \quad \forall t \geq 0.$$

618 From (6.17) it is easy to deduce that

$$619 \quad (6.18) \quad i \in \{1, \dots, 2n\} \setminus \mathbf{S}(\boldsymbol{\lambda}) \text{ and } \langle A \tilde{\boldsymbol{e}}_i, \boldsymbol{d} \rangle \leq 0 \implies i \in \{1, \dots, 2n\} \setminus \mathbf{S}(\boldsymbol{\lambda} + t\boldsymbol{d}) \quad \forall t \geq 0.$$

620 Case 3. We begin by noticing that

621 (6.19) $\{i \in \{1,\ldots,2n\} \setminus S(\boldsymbol{\lambda}) : \langle A\tilde{\boldsymbol{e}}_i, \boldsymbol{d} \rangle > 0\} = \{i \in \{1,\ldots,2n\} : \langle A\tilde{\boldsymbol{e}}_i, \boldsymbol{d} \rangle > 0\}.$

Indeed, we recall that from proposition 3.9 any descent direction \boldsymbol{d} for g at $\boldsymbol{\lambda}$ implies that for every $i \in S(\boldsymbol{\lambda})$ we have $\langle \boldsymbol{d}, A\tilde{\boldsymbol{e}}_i \rangle \leq 0$. Hence, if $\langle A\tilde{\boldsymbol{e}}_i, \boldsymbol{d} \rangle > 0$ for some $i \in$ $\{1, \ldots, 2n\}$ then $i \notin S(\boldsymbol{\lambda})$. This means that $\{i \in \{1, \ldots, 2n\} \setminus S(\boldsymbol{\lambda}) : \langle A\tilde{\boldsymbol{e}}_i, \boldsymbol{d} \rangle > 0\} \supset$ $\{i \in \{1, \ldots, 2n\} : \langle A\tilde{\boldsymbol{e}}_i, \boldsymbol{d} \rangle > 0\}$ and therefore proves (6.19). The converse inclusion is trivial. Hence, *Case 3* is, from (3.11), the case defined by $S^+(\boldsymbol{d})$. For any $i \in S^+(\boldsymbol{d})$, it is easy to see that

628 (6.20)
$$i \in S^+(d) \Rightarrow \langle \boldsymbol{\lambda} + t\boldsymbol{d}, A\tilde{\boldsymbol{e}}_i \rangle \leq 1 \text{ iff } t \in \left[0, \frac{1 - \langle A\tilde{\boldsymbol{e}}_i, \boldsymbol{\lambda} \rangle}{\langle A\tilde{\boldsymbol{e}}_i, \boldsymbol{d} \rangle}\right].$$

629 From (6.20) we obviously deduce

630
$$i \in S^+(d) \Rightarrow \langle \boldsymbol{\lambda} + t \boldsymbol{d}, A \tilde{\boldsymbol{e}}_i \rangle < 1 \text{ iff } t \in \left[0, \frac{1 - \langle A \tilde{\boldsymbol{e}}_i, \boldsymbol{\lambda} \rangle}{\langle A \tilde{\boldsymbol{e}}_i, d \rangle}\right]$$

631 and therefore that

632 (6.21)
$$i \in S^+(d) \Rightarrow i \in \{1, \dots, 2n\} \setminus S(\lambda + td) \text{ iff } t \in \left[0, \frac{1 - \langle A\tilde{e}_i, \lambda \rangle}{\langle A\tilde{e}_i, d \rangle}\right).$$

633 In addition, for any $i \in S^+(d)$ we have that $\langle A\tilde{\boldsymbol{e}}_i, \boldsymbol{\lambda} \rangle < 1$ and therefore that $1 - \langle A\tilde{\boldsymbol{e}}_i, \boldsymbol{\lambda} \rangle > 0$. Since, for any $i \in S^+(d)$, we also have $\langle A\tilde{\boldsymbol{e}}_i, \boldsymbol{d} \rangle > 0$ and we deduce

635 (6.22)
$$\forall i \in S^+(d) \text{ we have } \frac{1 - \langle A \tilde{e}_i, \lambda \rangle}{\langle A \tilde{e}_i, d \rangle} > 0.$$

From (6.16), (6.17) and (6.20) we deduce that $S^+(d) = \emptyset$ iff $(\lambda + td) \in \text{dom } g$ for every $t \ge 0$. In addition, from (6.16), (6.17) and (6.20) we deduce that if $S^+(d) \ne \emptyset$

⁶¹⁰

then $(\lambda + td) \in \text{dom } g$ for every $t \in [0, \Delta t]$, where Δt is defined by (3.12). The 638 fact that $\overline{\Delta t} > 0$ follows from (6.22) and, again (3.12). It remains to prove that for 639 any $t \in [0, \overline{\Delta t})$ we have $S(\lambda + td) \subset S(\lambda)$. To this aim we consider an arbitrary 640 $i \in \{1, \ldots, 2n\} \setminus S(\lambda)$. Combining (6.18), (6.21) and the definition of Δt as a mini-641 mum (3.12), we deduce that $i \in \{1, \ldots, 2n\} \setminus S(\boldsymbol{\lambda} + t\boldsymbol{d})$ for any $t \in [0, \Delta t)$. Hence, 642 by considering the complementary set we obtain that for any $t \in [0, \Delta t)$ we have 643 $S(\boldsymbol{\lambda} + t\boldsymbol{d}) \subset S(\boldsymbol{\lambda})$. Furthermore, the fact that $\partial g(\boldsymbol{\lambda} + t\boldsymbol{d}) \subset \partial g(\boldsymbol{\lambda})$ for all $t \in [0, \overline{\Delta t})$ 644immediately follows from lemma 3.8. It is easy to see that for every $t \in [0, \overline{\Delta t}]$ we 645 have $(\boldsymbol{\lambda} + t\boldsymbol{d}) \in \text{dom } q$. П 646

647 **6.9. Proof of lemma 3.11 on page 8.**

648 *Proof.* The proof is in three steps. The first step justifies the well-posedness 649 of (3.15). The second step proves that the conditions in (3.16) are equivalent. The 650 last step proves that $\lambda \in \text{dom } g$ is a solution to (D_{ℓ^1}) iff the conditions in (3.16) hold 651 true.

From lemma 3.8, for any $\lambda \in \text{dom } g$ we have that $\partial g(\lambda) \neq \emptyset$ and obviously closed, convex. Therefore, for any $\lambda \in \text{dom } g$ (3.15) is well posed.

Combining the definitions of d (3.15) and of $S^+(d)$ (3.11) we have that d = 0 implies S⁺(d) = \emptyset . Conversely, if $S^+(d) = \emptyset$ then we obtain that $\langle d, Ae_i \rangle = 0$ for every canonical vector e_i of \mathbb{R}^m . From proposition 3.1 we have that A has full row rank and therefore deduce d = 0. Thus, $S^+(d) = \emptyset$ is equivalent to d = 0. From the definitions of $S^+(d)$ (3.11) and of $\overline{\Delta t}$ (3.12) it is obvious that $S^+(d) = \emptyset$ is equivalent to $\overline{\Delta t} = +\infty$. Thus, the three conditions in (3.16) are equivalent.

From theorem 8.11 we have that $\lambda \in \text{dom } g$ is a solution to (D_{ℓ^1}) iff $\mathbf{0} \in \partial g(\lambda)$. Hence,

661 $\lambda \in \text{dom } g$ is a solution to (D_{ℓ^1}) iff d defined by (3.15) satisfies d = 0. It follows that 662 $\lambda \in \text{dom } g$ is a solution to (D_{ℓ^1}) iff the conditions in (3.16) hold true.

663 6.10. Proof of proposition 3.12 on page 8.

664 *Proof.* We begin to establish the following lemmas that will be useful for the proof 665 of propositions 3.12.

666 LEMMA 6.5 (Technical lemma). Consider a convex set $\emptyset \neq K \subset \mathbb{R}^m$. For any 667 \boldsymbol{x} we have $\Pi_{K+\boldsymbol{x}}(\boldsymbol{0}) = \Pi_K(-\boldsymbol{x}) + \boldsymbol{x}$.

668 *Proof.* From theorem 8.8, we have that a vector $\boldsymbol{y}_{\boldsymbol{x}}$ is the projection of some \boldsymbol{x} 669 on K iff $\langle \boldsymbol{x} - \boldsymbol{y}_{\boldsymbol{x}}, \boldsymbol{y} - \boldsymbol{y}_{\boldsymbol{x}} \rangle \leq 0 \quad \forall \boldsymbol{y} \in K$. Hence, the projection $\boldsymbol{y}_{-\boldsymbol{x}} := \Pi_{K}(-\boldsymbol{x})$ of $-\boldsymbol{x}$ 670 onto K satisfies $\langle -\boldsymbol{x} - \boldsymbol{y}_{-\boldsymbol{x}}, \boldsymbol{y} - \boldsymbol{y}_{-\boldsymbol{x}} \rangle \leq 0$ for all $\boldsymbol{y} \in K$. Therefore we obtain, for all 671 $\boldsymbol{y} \in K - \boldsymbol{x}$, that

672
$$\langle -\boldsymbol{x} - \boldsymbol{y}_{-\boldsymbol{x}}, (\boldsymbol{y} - \boldsymbol{x}) - \boldsymbol{y}_{-\boldsymbol{x}} \rangle \leq 0 \quad \forall \boldsymbol{y} \in K - \boldsymbol{x} \Leftrightarrow \langle \boldsymbol{0} - (\boldsymbol{x} + \boldsymbol{y}_{-\boldsymbol{x}}), \boldsymbol{y} - (\boldsymbol{y}_{-\boldsymbol{x}} + \boldsymbol{x}) \rangle \leq 0$$

and, from theorem 8.8 we obtain that $\boldsymbol{x} + \boldsymbol{y}_{-\boldsymbol{x}}$ is the projection of **0** on $K - \boldsymbol{x}$. In other words, $\Pi_K(-\boldsymbol{x}) + \boldsymbol{x} = \Pi_{K-\boldsymbol{x}}(\boldsymbol{0})$ and the formula is proved.

EEMMA 6.6 $(-\Pi_{\partial g(\lambda)}(\mathbf{0})$ satisfies (3.9)). We posit the same assumptions as in proposition 3.1. For any $\lambda \in \text{dom } g$ consider d defined by (3.15) in lemma 3.11. We have that d satisfies (3.9).

678 Proof. Consider $\lambda \in \text{dom } g$ and $S(\lambda)$ defined by (3.3) (page 5). We wish to 679 prove that $d := -\prod_{\partial g(\lambda)}(\mathbf{0})$ satisfies (3.9). From lemma 6.5 applied with $\boldsymbol{x} = \boldsymbol{b}$ and 680 $K := \operatorname{co} \{A\tilde{\boldsymbol{e}}_i, i \in S(\lambda)\}$ we obtain $\boldsymbol{d} = -\prod_{\partial g(\lambda)}(\mathbf{0}) = -\boldsymbol{b} - \prod_K(-\boldsymbol{b})$. Thus, from 681 theorem 8.8 we have that $\prod_K(-\boldsymbol{b})$ satisfies

682
$$\langle -\boldsymbol{b} - \Pi_K(-\boldsymbol{b}), \boldsymbol{y} - \Pi_K(-\boldsymbol{b}) \rangle \leq 0 \quad \forall \boldsymbol{y} \in K := \operatorname{co} \{A \tilde{\boldsymbol{e}}_i, i \in S(\boldsymbol{\lambda})\}$$

and therefore, since $\boldsymbol{d} = -\boldsymbol{b} - \Pi_K(-\boldsymbol{b})$, we obtain

684 (6.23)
$$\langle \boldsymbol{d}, \boldsymbol{y} - \Pi_K(-\boldsymbol{b}) \rangle \leq 0 \quad \forall \boldsymbol{y} \in K := \operatorname{co} \left\{ A \tilde{\boldsymbol{e}}_i, \ i \in \mathrm{S}(\boldsymbol{\lambda}) \right\}.$$

685 Any coefficients μ_i of $\Pi_K(-\boldsymbol{b})$ onto K satisfy

$$\Pi_K(-oldsymbol{b}) = \sum_{\substack{\mu_i \geqslant 0 \ orall \ i \in \mathrm{S}(oldsymbol{\lambda}) \ \mu_i = 0 \ \mathrm{otherwise}}} \mu_i A ilde{oldsymbol{e}}_i.$$

687 Consider any $j \in S(\lambda)$ and the coefficients α_i given by $\alpha_j = 1 + \mu_j$ where $\alpha_i = \mu_i$ 688 for $i \neq j \in S(\lambda)$ and $\alpha_i = 0$ otherwise. Note that the vector $\sum_i \alpha_i A \tilde{\boldsymbol{e}}_i \in K$ and that 689 $\sum_i \alpha_i A \tilde{\boldsymbol{e}}_i - \prod_K (-\boldsymbol{b}) = A \tilde{\boldsymbol{e}}_j$. Thus, from (6.23) we obtain that \boldsymbol{d} satisfies (3.9).

690 LEMMA 6.7 (Descent direction condition). We posit the same assumptions as 691 in proposition 3.1 and consider **d** defined by (3.15) in lemma 3.11. We have that if 692 $d \neq 0$ then **d** is a descent direction for g at $\lambda \in \text{dom } g$. In addition, for all $\lambda \in \text{dom } g$ 693 we have

694 (6.24)
$$g'\left(\boldsymbol{\lambda}, -\Pi_{\partial g(\boldsymbol{\lambda})}(\mathbf{0})\right) = -\left\|\Pi_{\partial g(\boldsymbol{\lambda})}(\mathbf{0})\right\|_{\ell^{2}}^{2}.$$

Proof. We consider d defined by (3.15) in lemma 3.11. We first establish (6.24) then justify that if $d \neq 0$ then d is a descent direction for g at $\lambda \in \text{dom } g$. From theorem 8.8, we have that

698 (6.25)
$$\langle d, s + d \rangle \leq 0 \quad \forall s \in \partial g(\lambda) \quad \Leftrightarrow \quad \langle d, s \rangle \leq - \|d\|_{\ell^2}^2 \quad \forall s \in \partial g(\lambda).$$

699 We have

700 (6.26)
$$-\|\boldsymbol{d}\|_{\ell^2}^2 \ge \langle \boldsymbol{d}, \boldsymbol{b} \rangle = g'(\boldsymbol{\lambda}, \boldsymbol{d}) = \sup\{\langle \boldsymbol{s}, \boldsymbol{d} \rangle : \boldsymbol{s} \in \partial g(\boldsymbol{\lambda})\} \ge \langle -\boldsymbol{d}, \boldsymbol{d} \rangle = -\|\boldsymbol{d}\|_{\ell^2}^2.$$

Indeed, in (6.26) the first inequality is obtained by choosing $\mathbf{s} = \mathbf{b} \in \partial g(\boldsymbol{\lambda})$ in (6.25). Combining lemma 6.6 and proposition 3.9 we obtain first equality in (6.26). The second equality is justified by the definition of the subdifferential (see (**xvii**)). The second inequality follows from $-\mathbf{d} \in \partial g(\boldsymbol{\lambda})$. The last equality is obvious. Thus, we obtain (6.24). From (6.24), it follows that if $\mathbf{d} \neq 0$ we have that \mathbf{d} satisfies (3.9)-(3.10). Hence, from proposition 3.9 we obtain that \mathbf{d} is a descent direction.

Consider d defined by (3.15) in lemma 3.11. If d = 0 then (3.17)-(3.18) hold true. From now on, we assume that $d \neq 0$. Let $t \in [0, \overline{\Delta t})$, where $\overline{\Delta t}$ is defined in proposition 3.10 (page 8). For any $\lambda' \in \mathbb{R}^m$ we have

710 (6.27)
$$g(\lambda') \ge g(\lambda) + \langle \Pi_{\partial g(\lambda)}(\mathbf{0}), \lambda' - \lambda \rangle$$

711 (6.28)
$$= g(\boldsymbol{\lambda}) + \left\langle \Pi_{\partial g(\boldsymbol{\lambda})}(\mathbf{0}), \boldsymbol{\lambda}' - \boldsymbol{\lambda} - td \right\rangle - t \left\| \Pi_{\partial g(\boldsymbol{\lambda})}(\mathbf{0}) \right\|_{\ell^2}^2$$

712 (6.29)
$$= g(\boldsymbol{\lambda}) + \langle \Pi_{\partial g(\boldsymbol{\lambda})}(\boldsymbol{0}), \boldsymbol{\lambda}' - \boldsymbol{\lambda} - td \rangle + tg'(\boldsymbol{\lambda}, d)$$

713 (6.30)
$$= g(\boldsymbol{\lambda}) + \langle \Pi_{\partial g(\boldsymbol{\lambda})}(\boldsymbol{0}), \boldsymbol{\lambda}' - (\boldsymbol{\lambda} + t\boldsymbol{d}) \rangle + t \langle \boldsymbol{d}, \boldsymbol{b} \rangle$$

714 (6.31)
$$= g(\boldsymbol{\lambda} + t\boldsymbol{d}) + \left\langle \Pi_{\partial g(\boldsymbol{\lambda})}(\boldsymbol{0}), \boldsymbol{\lambda}' - (\boldsymbol{\lambda} + t\boldsymbol{d}) \right\rangle.$$

The inequality in (6.27) is nothing but the definition of $\Pi_{\partial g(\lambda)}(\mathbf{0}) \in \partial g(\lambda)$ (see (xvii)) and (6.28) follows. Lemma 6.7 (we assumed $\mathbf{d} \neq \mathbf{0}$) justifies (6.29). From proposi-

717 tion 3.10 (page 8), for any $t \in [0, \overline{\Delta t})$ we have $(\lambda + td) \in \text{dom } g$ and, from lemma 6.4

we obtain (6.30). Equation (6.31) immediately follows from (3.7) in proposition 3.7.

From (6.27)-(6.31) and (xvii) we obtain (3.17). From theorem 8.8, we have that 719 720 $\Pi_{\partial q(\boldsymbol{\lambda})}(\mathbf{0})$ satisfies

$$\forall \boldsymbol{s} \in \partial g(\boldsymbol{\lambda}), \qquad \left\langle -\Pi_{\partial g(\boldsymbol{\lambda})}(\boldsymbol{0}), \boldsymbol{s} - \Pi_{\partial g(\boldsymbol{\lambda})}(\boldsymbol{0}) \right\rangle \leqslant 0$$

which is equivalent to 722

721

723
$$\forall s \in \partial g(\lambda) \qquad \left\| \Pi_{\partial g(\lambda)}(\mathbf{0}) \right\|_{\ell^2}^2 \leq \left\langle s, \Pi_{\partial g(\lambda)}(\mathbf{0}) \right\rangle.$$

724 Hence, from (3.14) in proposition 3.10 (page 8) we deduce that $\Pi_{\partial q(\lambda)}(\mathbf{0})$ satisfies

725 (6.32)
$$\forall t \in [0, \overline{\Delta t})$$
 $\forall s \in \partial g(\lambda + td)$ $\|\Pi_{\partial g(\lambda)}(\mathbf{0})\|_{\ell^2}^2 \leq \langle s, \Pi_{\partial g(\lambda)}(\mathbf{0}) \rangle$

Combining (3.17), (6.32) and, again, theorem 8.8 we obtain (3.18). 726

6.11. Proof of proposition 3.13 on page 8. 727

Proof. The proof is in two steps. We first justify the well-posedness of (3.19) then 728 justify that (3.20) coincides with the evolution equation (3.6) (see proposition 3.6 on 729page 6). 730

Step 1. Let k = 0. By assumption, $\lambda(t_k) \in \text{dom } g$. From lemma 3.11 we have that d_k 731 is well defined. From proposition 3.10 (page 8) this implies that t_{k+1} is well defined. 732 s In addition, from, again, proposition 3.10 and the definition of t_{k+1} it is easy to see 733 that $\lambda(t) \in \text{dom } g$ for every $t \in [t_k, t_{k+1}]$. The rest of the recursion follows. Thus, 734 we obtain that the trajectory $\lambda(t)$ given in (3.20) is mathematically well-posed. It 735 remains to show that (3.20) coincides with the trajectory given by (3.6). 736

Step 2. From proposition 3.12, the vector $-\Pi_{\partial q(\lambda)}(\mathbf{0})$ that appears in (3.6) is piece-737 wise constant on every intervals $[t_k, t_{k+1})$. In addition, it is easy that the trajectory 738 739 given by (3.20) coincides by construction with the solution to the evolution equation (3.6) for every $t \ge 0$. The fact that for every $t \ge 0$, $\lambda(t) \in \text{dom } g$ follows 740 combining proposition 3.6 and lemma 3.8. This concludes our proof. Π 741

742 6.12. Proof of proposition 3.14 on page 9.

Proof. From (3.17) and the lower semi-continuity of $g \in \Gamma_0(\mathbb{R}^m)$, we obtain 743 $\Pi_{\partial g(\boldsymbol{\lambda}(t_k))}(\mathbf{0}) \in \partial g(\boldsymbol{\lambda}(t_{k+1}))$ and therefore that 744

745 (6.33)
$$\left\| \Pi_{\partial g(\boldsymbol{\lambda}(t_{k+1}))}(\mathbf{0}) \right\|_{\ell^2} \leq \left\| \Pi_{\partial g(\boldsymbol{\lambda}(t_k))}(\mathbf{0}) \right\|_{\ell^2}$$

From proposition/definition 3.13 (page 8), we have that $\lambda(t_{k+1}) \in \text{dom } g$ and there-746

fore, from lemma 3.8 we have $\partial g(\lambda(t_{k+1})) \neq \emptyset$. The uniqueness of the projection of 747 **0** onto the non-empty closed convex set $\partial g(\lambda(t_{k+1}))$ and (3.17) imply that 748

749 (6.34)
$$\|\Pi_{\partial g(\boldsymbol{\lambda}(t_k))}(\mathbf{0})\|_{\ell^2}^2 = \|\Pi_{\partial g(\boldsymbol{\lambda}(t_{k+1}))}(\mathbf{0})\|_{\ell^2}^2 \Leftrightarrow \Pi_{\partial g(\boldsymbol{\lambda}(t_k))}(\mathbf{0}) = \Pi_{\partial g(\boldsymbol{\lambda}(t_{k+1}))}(\mathbf{0}).$$

750 We wish to prove that $\prod_{\partial g(\lambda(t_k))}(\mathbf{0}) \neq \prod_{\partial g(\lambda(t_{k+1}))}(\mathbf{0})$. To do so, we set $d_k :=$ $-\Pi_{\partial g(\boldsymbol{\lambda}(t_k))}(\mathbf{0})$ and $d_{k+1} := -\Pi_{\partial g(\boldsymbol{\lambda}(t_{k+1}))}(\mathbf{0})$ and denote by $\overline{\Delta}t_k$ (resp. $\overline{\Delta}t_{k+1}$) the 751 positive kick times computed, from proposition 3.10, at $\lambda(t_k)$ (resp. $\lambda(t_{k+1})$). 752

By assumption, we have that $\lambda(t_k)$ is not a solution to (D_{ℓ^1}) . From lemma 3.11 753 we have that $d_k \neq 0$. Assume, for the sake of contradiction, that $d_k = d_{k+1}$. From, 754again, proposition 3.10, we would have $(\lambda(t_k) + \Delta t_k d_k) + t d_k = \lambda(t_{k+1}) + t d_k \in \text{dom } g$, 755

for some positive $t \in (0, \overline{\Delta t}_{k+1})$. This is impossible. Indeed, proposition 3.10 applied 756

at $\lambda(t_k)$ with the direction d_k implies that $(\lambda(t_k) + td_k) \in \text{dom } g \text{ iff } t \in [0, \overline{\Delta t_k}].$ 757

6.13. Proof of proposition 3.15 on page 9. 759

Proof. The proof is in two steps. The first step justifies the existence of $K \in \mathbb{N}$ 760 and of t_K such that $\lambda(t) = \lambda(t_K)$ for every $t \ge t_K$. The second step justifies that 761 $\lambda(t_K)$ is a solution to (D_{ℓ^1}) and d_K satisfies $d_K = 0$. 762

Step 1. This part of the proof follows a classic approach that can be found in, 763 e.g., [22, Thm. 3.4.8, p. 382]. From lemma 3.8 (page 7) there are 4^n possible sets 764 $\partial g(\boldsymbol{\lambda})$ for $\boldsymbol{\lambda} \in \text{dom } g$. Each of them is uniquely associated with $\boldsymbol{d}_k = -\prod_{\partial g(\boldsymbol{\lambda})} (\mathbf{0})$. 765 From proposition 3.14 they are all different from each other. This implies that the 766 sequence $(d_k)_k$ has a finite number of terms (is finite) and therefore converges for a 767 768 finite index K. From proposition 3.10 (page 8) it is easy to deduce that the sequence $(t_k)_k$ is also finite. From proposition 3.13 we deduce the existence of K. In other 769 words, we obtained that the trajectory $\lambda(t)$ given in (3.20) satisfies, for some $K \in \mathbb{N}$, 770771 $\lambda(t) = \lambda(t_K)$ for every $t \ge t_K$. We now turn to the second step of the proof.

Step 2. From proposition 3.13, we have that the trajectory $\lambda(t)$ given in (3.20) 772 773 coincides with the trajectory given by (3.6). Thus, the limit of $\lambda(t)$ when $t \to +\infty$ is a solution to (D_{ℓ^1}) . (The fact that (D_{ℓ^1}) has a solution is justified by lemma 3.4 774page 5). From step1, we have that the limit of $\lambda(t)$ is attained for $t = t_K$. Hence, 775776 $\lambda(t_K)$ is a solution to (3.4). From lemma 3.11 we immediately obtain $d_K = 0$.

7. Glossary of notations, definitions. 777

- (i) (Vectors) Throughout the paper the vectors of, e.g., \mathbb{R}^n are denoted in bold 778typeface, e.g., \boldsymbol{x} . Other objects like scalars or functions are denoted in non-779 bold typeface. 780
 - (ii) (Canonical vectors) Throughout the paper the *i*-th canonical vectors of, e.g., \mathbb{R}^n are denoted by \boldsymbol{e}_i .
- (iii) (Inner product) For $x, y \in \mathbb{R}^n$ we denote by $\langle x, y \rangle$ the Euclidean inner product 783 in \mathbb{R}^n . 784
- 785(iv) (Interior) int(E): interior of a set E.
- (v) (Conical hull) co $\{a_1, \ldots, a_p\} := \{\sum_{i=1}^p \mu_i a_i : \mu_i \ge 0, \forall i = 1, \ldots, p\}$ (vi) Normal cone to a convex set $C \ne \emptyset$ at $\lambda \in C$: 786
- 787
- $N_C(\boldsymbol{\lambda}) = \{ \boldsymbol{s} : \langle \boldsymbol{s}, \boldsymbol{s}' \boldsymbol{\lambda} \rangle \leq 0 \ \forall \boldsymbol{s}' \in C \}$ (See, e.g., [22, Def. 5.2.3, p. 136]) 788
- (vii) (the $\ell^{\infty}(\mathbb{R}^n)$ unit ball) $B_{\infty} = \{ \boldsymbol{u} : \langle \boldsymbol{u}, \tilde{\boldsymbol{e}}_i \rangle \leq 1, i = 1, \dots, 2n \}$ where $\tilde{\boldsymbol{e}}_i$ is given 789 790by (3.3).
- (viii) (Effective domain) dom f: The domain of a convex function f is the (convex, 791 possibly empty) set dom $f = \{ \boldsymbol{x} \in \mathbb{R}^n : f(\boldsymbol{x}) \in \mathbb{R} \}$ 792
- (ix) (Convex function) A function $f : \mathbb{R}^n \to \mathbb{R} \cup \{+\infty\}$ is said to be convex if 793 $\forall (\boldsymbol{x}, \boldsymbol{y}) \in \mathbb{R}^n \times \mathbb{R}^n \text{ and } \forall \alpha \in (0, 1) \ f(\alpha \boldsymbol{x} + (1 - \alpha) \boldsymbol{y}) \leq \alpha f(\boldsymbol{x}) + (1 - \alpha) f(\boldsymbol{y})$ 794holds true (in $\mathbb{R} \cup \{+\infty\}$). 795
- (x) (Set $\Gamma_0(\mathbb{R}^n)$) The set of lower semi-continuous, convex functions with 796dom $f \neq \emptyset$ is denoted $\Gamma_0(\mathbb{R}^n)$
- (xi) (Characteristic function of a set) $\chi_E(\mathbf{x}) = 0$ if $\mathbf{x} \in E$ and $\chi_E(\mathbf{x}) = +\infty$ 798 otherwise 799
- (xii) (Polyhedral convex function) f is a polyhedral convex function if f(u) =800 $h(\boldsymbol{u}) + \chi_{\mathbf{C}}(\boldsymbol{u}) \ h(\boldsymbol{u}) = \max_{i=1,\dots,p} \left(\langle \boldsymbol{u}, \boldsymbol{a}_i \rangle - r_i \right)$ and 801 802

$$C = \{ \boldsymbol{u} \in \mathbb{R}^n \langle \boldsymbol{u}, \boldsymbol{\alpha}_i \rangle \leq \rho_i, \ i = 1, \dots, q \}.$$

(xiii) (Directional derivative) For $f \in \Gamma_0(\mathbb{R})$ at $a \in \text{dom } f$ in the direction d is 803 $f'(\boldsymbol{a}, \boldsymbol{d}) := \lim_{t \to 0^+} \frac{f(\boldsymbol{a}+t\boldsymbol{d}) - f(\boldsymbol{a})}{t}$ 804

805 (xiv) (Right derivative)
$$\frac{d^+ \lambda(t)}{dt} := \lim_{h \to 0^+} \frac{\lambda(t+h) - \lambda(t)}{h}$$

806 (xv) (Descent direction) $d \neq 0$ is a descent direction for f at x if $\exists t > 0$ such that

781

782

- 807 $\mathbf{x} + t\mathbf{d} \in \text{dom } f \text{ and } f(\mathbf{x} + t\mathbf{d}) < f(\mathbf{x}) \text{ (See, e.g., [22, Def. 1.1.1, p. 343]).}$ 808 (**xvi**) (Convex conjugate) For any $f \text{ convex that satisfies dom } f \neq \emptyset$, the function 809 $f^* \text{ defined by } \mathbb{R}^n \ni \mathbf{s} \mapsto f^*(\mathbf{s}) := \sup_{\mathbf{x} \in \text{dom } f} \{\langle \mathbf{s}, \mathbf{x} \rangle - f(\mathbf{x})\}.$ (See , e.g., [23, 810 Def. 1.1.1, p. 37]). For any $f \in \Gamma_0(\mathbb{R}^n)$ we have $f^* \in \Gamma_0(\mathbb{R}^n)$ (See, e.g., [22, 811 Thm. 1.1.2., p. 38])
- 812 (xvii) (Sub-differential) For $f \in \Gamma_0(\mathbb{R}^n)$ and $x \in \text{dom } f$ the vector $s \in \mathbb{R}^n$ is a 813 subgradient of f at x if one of the following equivalent assertions is satisfied

814 (7.1)
$$\forall \boldsymbol{y} \in \mathbb{R}^n$$
, $f(\boldsymbol{y}) \ge f(\boldsymbol{x}) + \langle \boldsymbol{s}, \boldsymbol{y} - \boldsymbol{x} \rangle$; or $\forall \boldsymbol{d} \in \mathbb{R}^n$, $\langle \boldsymbol{s}, \boldsymbol{d} \rangle \le f'(\boldsymbol{x}, \boldsymbol{d})$.

- 815 We denote by $\partial f(\boldsymbol{x})$ the closed convex set of vectors $\boldsymbol{s} \in \mathbb{R}^n$ that satisfy (7.1). 816 For $\boldsymbol{x} \notin \text{dom } f$ we set $\partial f(\boldsymbol{x}) := \emptyset$.
- 817 (**xviii**) (Euclidean projection) $\Pi_C(\boldsymbol{x}) = \arg \min_{\boldsymbol{y} \in C} \|\boldsymbol{y} \boldsymbol{x}\|_{\ell^2}$ for $C \neq \emptyset$ closed and 818 convex.
- 819 **8. Mathematical background.** This section contains several propositions and 820 theorems used throughout proofs given in section 6.
- 821 THEOREM 8.1. For any $f \in \Gamma_0(\mathbb{R}^n)$ we have $f^* \in \Gamma_0(\mathbb{R}^n)$.

822 Proof. From $f \in \Gamma_0(\mathbb{R}^n)$ we have $f \neq 0$, and from [22, Pro. 1.2.1, p. 147] there is 823 an affine function minorizing f on \mathbb{R}^n . Applying [23, Thm. 1.1.2, p. 38] we conclude 824 that $f^* \in \Gamma_0(\mathbb{R}^n)$.

LEMMA 8.2 (Conjugate of absolute value). [4, Table 3.1, p. 76] Let $f: x \in \mathbb{R} \mapsto f(x) := |x|$. We have for all $v \in \mathbb{R}$, $f^*(v) = \chi_{-1,1}(v) = \begin{cases} 0 & \text{if } v \in [-1,1] \\ +\infty & \text{otherwise.} \end{cases}$

EEMMA 8.3 (Conjugate of characteristic function). [23, Ex. 1.1.5, p. 39]] The conjugate of the characteristic function of the nonempty convex set set C (see (xi)) is for all $v \in \mathbb{R}^n$, $\chi^*_{\mathbf{C}}(v) = \sup_{x \in C} \langle v, x \rangle$.

830 PROPOSITION 8.4 (Conjugation in Product Spaces). [36, Prop. 11.22, p. 493] 831 Let f_1, \ldots, f_n be in $\Gamma_0(\mathbb{R})$, and $f : \mathbb{R}^n \to \mathbb{R} \cup \{+\infty\}$ given by: $\forall (x_1, \ldots, x_n) \in$ 832 \mathbb{R}^n , $f(x_1, \ldots, x_n) = f_1(x_1) + \ldots + f_n(x_n)$. Then $f^*(v_1, \ldots, v_n) = f_1^*(v_1) + \ldots + f_n^*(v_n)$.

THEOREM 8.5 (Pre-composition With a Matrix). [22, Prop. 2.1.5, p. 159] Let $f \in \Gamma_0(\mathbb{R}^m)$ and $A \in \mathcal{M}_{m \times n}(\mathbb{R})$, and assume that span $A \cap dom f \neq \emptyset$. We have $f(A \cdot) \in \Gamma_0(\mathbb{R}^n)$.

EEMMA 8.6 (Subdifferential of Normal Cone to a Closed Convex Set). [23, Def. 1.1.3, p. 93] The set of normal directions to a closed convex set $C \subset \mathbb{R}^m$ at $\lambda \in C$, is the subdifferential of the characteristic function $\chi_{\mathbf{C}}$ at λ : $N_C(\lambda) := \partial \chi_{\mathbf{C}}$.

EEMMA 8.7. [22, Ex. 5.2.6 b), p.138] Let a closed convex polyhedron defined by $C := \{ \boldsymbol{u} \in \mathbb{R}^n : \langle \boldsymbol{s}_i; \boldsymbol{x} \rangle \leq r_i \text{ for } i = 1, \dots, p \}$ where $\boldsymbol{s}_i \in \mathbb{R}^n$ and $r_i \in \mathbb{R}$ for all $i = 1, \dots, p$. The set of active constraints at $\boldsymbol{u} \in C$ by $W(\boldsymbol{u}) = \{i \in 1, \dots, p : \langle \boldsymbol{s}_i; \boldsymbol{x} \rangle = r_i\}$. Here we have $N_C(\boldsymbol{u}) = \operatorname{co} \{ \boldsymbol{s}_i : i \in W(\boldsymbol{u}) \}$.

843 PROPOSITION 8.8. [22, Thm 3.1.1, p. 117] Let C be a nonempty closed convex 844 set of \mathbb{R}^n . We have that $\mathbf{y}_x \in C$ is the Euclidean projection of some \mathbf{x} onto C if only 845 if $\langle \mathbf{x} - \mathbf{y}_x, \mathbf{y} - \mathbf{y}_x \rangle \leq 0$ for all $\mathbf{y} \in C$.

THEOREM 8.9 (Subdifferential of Pre-composition with a matrix). Let $f \in$ $\Gamma_0(\mathbb{R}^n)$ such that $int(dom f) \neq \emptyset$ and $A \in \mathcal{M}_{m \times n}(\mathbb{R})$. Assume that $int(dom f) \cap$ $span A \neq \emptyset$. Then, any $u \in \mathbb{R}^n$ such that $Au \in dom f$ we have $\partial(f(A \cdot))(u) =$ $A^T \partial f(Au)$. 850 Proof. Since we assumed that int $(\text{dom } f) \neq \emptyset$ we have ri $(\text{dom } f) \cap \text{span } A =$ 851 int $(\text{dom } f) \cap \text{span } A = \emptyset$. Thus, from [23, Thm. 3.2.1, p. 117] applied with $\varepsilon = 0$ 852 and g := f to conclude.

853 PROPOSITION 8.10. [22, Prop. 2.1.1, p. 158] Let $f_1 \in \Gamma_0(\mathbb{R}^n), ..., f_p \in \Gamma_0(\mathbb{R}^n)$ 854 and $t_l, ..., t_p$ be positive numbers. We assume that there is a point where all the f_j are 855 finite. Then the function $\sum_{i=1}^{p} t_i f_i \in \Gamma_0(\mathbb{R}^n)$.

THEOREM 8.11 (Fermat's rule). [36, Thm. 10.1, p. 422] Let $f \in \Gamma_0(\mathbb{R}^n)$. Then f has a global minimum at $\bar{\boldsymbol{u}}$ if only if $\boldsymbol{0} \in \partial f(\bar{\boldsymbol{u}})$.

PROPOSITION 8.12. [1, Prop. 1, p. 159] Let $f \in \Gamma_0(\mathbb{R}^n)$. Then the set-valued map $\mathbb{R}^n \ni \mathbf{u} \mapsto \partial f(\mathbf{u})$ is maximal monotone.

860 THEOREM 8.13 (Subdifferential of sum of Γ_0 -functions). Let $f_1, f_2 \in \Gamma_0(\mathbb{R}^n)$. 861 We assume that $int(dom f_1) \cap int(dom f_2) \neq \emptyset$. Then for all $\mathbf{u} \in dom (f_1 + f_2)$ we 862 have $\partial(f_1 + f_2)(\mathbf{u}) = \partial f_1(\mathbf{u}) + f_2(\mathbf{u})$.

863 Proof. Since int $(\text{dom } f_1) \cap \text{int} (\text{dom } f_2) \neq \emptyset$ we deduce that ri $(\text{dom } f_1) \cap$ 864 ri $(\text{dom } f_2) = \text{int} (\text{dom } f_1) \cap \text{int} (\text{dom } f_2) \neq \emptyset$. Thus, from [23, Cor. 3.1.2, p. 114] 865 applied with $\varepsilon = 0$ to conclude.

THEOREM 8.14. [5, Thm. 3.1, p. 54] Let T be a maximal monotone operator from \mathbb{R}^m to \mathbb{R}^m and dom (T) be its domain. Consider the problem $\frac{d\lambda(t)}{dt} \in -T(\lambda(t))$ with $\lambda(0) = \lambda_0$. For all $\lambda_0 \in dom$ (T), there exists a unique solution $\lambda(\cdot) : [0, +\infty) \rightarrow \mathbb{R}^m$ such that :

870 1. $\lambda(t) \in dom(T)$ for all t > 0, and $\lambda(0) = \lambda_0$;

871 2. the function $\lambda(\cdot)$ is continuous on $[0, +\infty)$;

872 3. the function $\lambda(\cdot)$ admits a right derivative $\frac{d^+\lambda(t)}{dt}$ at all $t \ge 0$, given by 873 $\frac{d^+\lambda(t)}{dt} = -\prod_{T(\lambda(t))}(\mathbf{0})$ for all $t \in [0, +\infty);$

874 4. the function $\frac{d^+}{dt} \lambda(\cdot)$ is continuous from the right on $[0, +\infty)$.

THEOREM 8.15. [1, Thm. 2, p. 160] Let $g \in \Gamma_0(\mathbb{R}^m)$, and assume that g achieves its minimum at some point. Then, for all $\lambda_0 \in dom(\partial g)$, the trajectory given by $\frac{d^+\lambda(t)}{dt} = -\prod_{\partial g(\lambda(t))}(\mathbf{0})$ with $\lambda(0) = \lambda_0$ converges to a point which minimizes g when $t \to +\infty$.

THEOREM 8.16. [1, Thm. 4, Eq (28), p. 35-36] Let $A \in \mathcal{M}_{m \times n}(\mathbb{R})$ and $U \in \Gamma_0(\mathbb{R}^n), V \in \Gamma_0(\mathbb{R}^m)$. Assume that $0 \in int(A \ dom \ U - dom \ V)$. Then, for all $u \in dom \ U \cap dom \ V(A \cdot)$ we have $\partial (U + V(A \cdot))(u) = \partial U(u) + A^T \partial V(A u)$.

882 PROPOSITION 8.17. [1, Prop. 1, p.163] Let $A \in \mathcal{M}_{m \times n}(\mathbb{R})$ and $U \in \Gamma_0(\mathbb{R}^n), V \in$ 883 $\Gamma_0(\mathbb{R}^m)$. Assume that $\mathbf{0} \in int(A^T \ dom \ V^* + dom \ U^*)$. Then, $\inf_{\mathbf{u} \in \mathbb{R}^n} (U(\mathbf{u}) + V(A\mathbf{u}))$ 884 has a solution.

THEOREM 8.18. [1, Thm. 2, p. 167] Let $A \in \mathcal{M}_{m \times n}(\mathbb{R})$, $U \in \Gamma_0(\mathbb{R}^n)$ and $V \in \Gamma_0(\mathbb{R}^m)$. Assume that assumptions of theorem 8.16 and proposition 8.17 hold. Then, $\inf_{\mathbf{\lambda} \in \mathbb{R}^m} (U^*(-A^T \mathbf{\lambda}) + V^*(\mathbf{\lambda}))$ has a solution.

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