Solving Constrained Variational Inequalities via a First-order Interior Point-based Method

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Abstract

We focus on the open problem to develop a first-order method that can solve constrained variational inequality (cVI) problems when given general constraints. We generalize the *alternating direction method of multipliers* (ADMM) method and combine it with interior-point approaches, yielding a first-order method that we refer to as **A**DMM-based interior-point method for **cVI**s (ACVI). We provide convergence guarantees for ACVI in two general classes of problems: *(i)* when the operator is ξ-monotone, and *(ii)* when it is monotone, some constraints are active and the game is not purely rotational. When the operator is, in addition, L-Lipschitz for the latter case, we match known lower bounds on rates for the gap function of $\mathcal{O}(1/\sqrt{K})$ and $\mathcal{O}(1/K)$ for the last and average iterate, respectively. To our knowledge, this is the first *first*-order interior-point method for the general cVI problem that has a global convergence guarantee. Empirical analyses demonstrate clear advantages of ACVI over common first-order methods. In particular, *(i)* cyclical behavior is notably reduced as our method approaches the solution from the analytic center, and *(ii)* unlike projection-based methods that zigzag when near a constraint, ACVI efficiently handles the constraints.

1. Introduction

We are interested in the *constrained variational inequality* problem [75]:

$$\text{find } \boldsymbol{x}^\star \in \mathcal{X} \quad \text{s.t.} \quad \langle \boldsymbol{x} - \boldsymbol{x}^\star, F(\boldsymbol{x}^\star) \rangle \geq 0, \quad \forall \boldsymbol{x} \in \mathcal{X}\,, \tag{cVI}$$

where $\mathcal{X}$ is a subset of the Euclidean n-dimensional space $\mathbb{R}^n$, and where $F : \mathcal{X} \mapsto \mathbb{R}^n$ is a continuous map. Finding (an element of) the solution set $\mathcal{S}^\star_{\mathcal{X},F}$ of cVI is a key problem in multiple fields such as economics and game theory. More pertinent to machine learning, CVIs generalize standard single-objective optimization, complementarity problems [17], zero-sum games [70, 82] and multi-player games. For example, solving cVI is the optimization problem underlying *reinforcement learning* [e.g., 62]—and *generative adversarial networks* [37]. Moreover, even when training one set of parameters with one loss f, that is $F(\boldsymbol{x}) \equiv \nabla_{\boldsymbol{x}} f(\boldsymbol{x})$, a natural way to improve the model's robustness in some regard is to introduce an adversary to perturb the objective or the input, or to consider the worst sample distribution of the empirical objective. As has been noted in many problem domains, including robust classification [54], adversarial training [76], causal inference [16], and robust objectives [e.g., 71], this leads to a min-max structure, which is an instance of the cVI problem. Solving cVIs is significantly more challenging than solving classical single-objective optimization problems, due to the fact that F is a general vector field, leading to possible "rotational" trajectories in parameter space. In response, the development of efficient algorithms with provable convergence has recently been the focus of interest in machine learning and optimization, particularly in the unconstrained setting, where $\mathcal{X} \equiv \mathbb{R}^n$ [e.g., 3, 8, 13, 19, 36, 38, 56, 57, 78].

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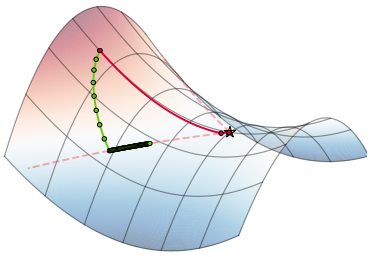


Figure 1: ACVI and EG—shown in *red* and *green*, resp.—on:

$$\min_{x_1 \in \mathbb{R}_+} \max_{x_2 \in \mathbb{R}_+} 0.05 \cdot x_1^2 + x_1 x_2 - 0.05 \cdot x_2^2 \,.$$

The constraints are depicted with dashed lines and the iterates with circles. ACVI gets close to the Nash Equilibrium ($\star$) in a single step, whereas EG zigzags when hitting a constraint. The remaining methods perform similarly to EG, see App. G.

In many applications, however, we have *constraints* on (part of) the decision variable $\boldsymbol{x}$, that is, $\mathcal{X}$ is often a strict subset of $\mathbb{R}^n$. As an example, let us revisit the aforementioned distributionally robust prediction problem: consider a linear setting [cf. Eq. 1 in 71] and class of parametrized distributions $\triangle \equiv \{\boldsymbol{w} \in \mathbb{R}^d | \boldsymbol{w} \geq 0, \boldsymbol{e}^\intercal \boldsymbol{w} = 1\}$, where $\boldsymbol{e} \in \mathbb{R}^d$ is a vector of all ones. Thus, the robust problem is: $\min_{\boldsymbol{x} \in \mathbb{R}^n} \max_{\boldsymbol{w} \in \mathbb{R}^d} \boldsymbol{w}^\intercal(\boldsymbol{y} - \boldsymbol{D}\boldsymbol{x}), \;\; \textit{subject to} \;\; \boldsymbol{w} \geq 0, \;\; \boldsymbol{e}^\intercal \boldsymbol{w} = 1$, where $\boldsymbol{D} \in \mathbb{R}^{d \times n}$ contains d samples of an n-dimensional covariate vector, and $\boldsymbol{y} \in \mathbb{R}^d$ is the vector of target variables (the constraint $\boldsymbol{w} \leq 1$ is implied), leading to an instance of the cVI problem. Additional example applications include (i) applications in business, finance, and economics where often the sum of the decision variables—representing for example *resources*—cannot exceed a specific value, (ii) contract theory (e.g. §2.3.2 in [5]) and (iii) solving optimal control problems numerically, among others. Significantly fewer works address the convergence of first-order optimization methods in the constrained setting; see § A for an overview. Recently, Cai et al. [11] established a convergence rate for the *projected extragradient method* [46], when F is monotone and Lipschitz (see § 2 for definitions). However, (i) the presented proof is computer-assisted making it hard to interpret, and (ii) the considered setting assumes the projection is fast to compute and thus ignores the projection in the rate. The latter assumption only holds in rare cases when the constraints are simple, so that operations such as clipping suffice. However, when given general constraints, each EG update requires *two* projections (see App. B.5). Each projection requires solving a separate *constrained* optimization problem, which if given general constraints requires using a *second-order method* as explained next.

Interior point (IP) methods are the de facto family of iterative algorithms for constrained optimization. These methods enjoy well established guarantees and theoretical understanding in the context of single-objective optimization [see e.g., 9, 55, 84] and have extensions to a wide range of problem settings [e.g., 59, 60, 69, 77, 85]. Several works extend IP methods to the cVI problem, by applying the second-order Newton method to a modified *Karush-Kuhn-Tucker* (KKT) system appropriate for cVI [15, 27, 58, 65, 68]. Many of these approaches, however, rely on strong assumptions—see App. A for more details. Moreover, although these methods enjoy fast convergence in terms of iterations, each iteration involves the computation of the Jacobian of F which quickly becomes prohibitive for large dimension n. Hence first-order methods are preferred in practice.

We are currently missing a *first*-order optimization method for solving cVI with *general* constraints. Accordingly, in this paper, we focus on the following open question:

*Can we derive **first-order** algorithms for the cVI problem that (i) can be applied when **general** constraints are given, and that (ii) have **global** convergence guarantees?*

In this paper, to mitigate the computational burden of second-derivative computation we replace the Newton step of the traditional IP methods with the *alternating direction method of multipliers* (ADMM) method. ADMM was designed with a different purpose: it is applicable only when the objective is separable into two or more different functions—see § B.3—and can be seen as equivalent to *Douglas–Rachford* operator splitting [24] applied in the *dual* space [see e.g. 51]. ADMM owes its

popularity primarily to its computational efficiency [10] for large-scale problems and fast convergence in some settings [e.g., 61]. The core idea of our approach is to reformulate the original cVI problem with equality and inequality constraints via the KKT conditions, so as to apply ADMM in a way that the subproblems of the resulting algorithm have desirable properties (see § C). This framework can also be used to design novel algorithms for solving cVIs; see App. E. Our contributions include:

- Based on the KKT system for cVIs and the ADMM technique, we derive an algorithm that we refer to as the ***ADMM-based IP Method for cVIs*** (ACVI)—see § C and Alg. 1.
- We derive parameter-free global convergence guarantees of ACVI given two sets of assumptions: (i) when F is ξ-monotone, and (ii) when it is monotone, the constraints are active at the solution, and the game is not purely rotational. By further assuming F is a Lipschitz operator, we upper bound the rate of decrease of the gap function and we match the known lower bound for the gap function of $\mathcal{O}(1/\sqrt{K})$ for the last iterate—see § 3.
- Empirically, we document two notable advantages of ACVI over projection-based methods: (i) the ACVI iterates exhibit significantly reduced rotations, as they approach the solution from the analytic center, and (ii) while projection-based methods show extensive zigzagging when hitting a constraint, ACVI avoids this, resulting in more efficient updates—see § 4.

2. Preliminaries

The constraint set is defined as an intersection of finitely many inequalities and linear equalities:

$$\mathcal{C} = \{\boldsymbol{x} \in \mathbb{R}^n | \varphi_i(\boldsymbol{x}) \leq 0, i \in [m],\ \boldsymbol{C}\boldsymbol{x} = \boldsymbol{d}\}\,, \qquad \text{and} \quad \mathcal{C} \subseteq \mathcal{X}\,, \tag{CS}$$

where each $\varphi_i : \mathbb{R}^n \mapsto \mathbb{R}$, $\boldsymbol{C} \in \mathbb{R}^{p\times n}$, $\boldsymbol{d} \in \mathbb{R}^p$, with $rank(\boldsymbol{C}) = p$. With φ we denote the concatenated $\varphi_i(\cdot), i \in [m]$, and each $\varphi_i \in C^1(\mathbb{R}^n), i \in [m]$ and is convex. Let: $\mathcal{C}_{\leq} \triangleq \{\boldsymbol{x} \in \mathbb{R}^n \,|\varphi(\boldsymbol{x}) \leq \boldsymbol{0}\}$, $\mathcal{C}_{<} \triangleq \{\boldsymbol{x} \in \mathbb{R}^n \,|\varphi(\boldsymbol{x}) < \boldsymbol{0}\}$, and $\mathcal{C}_{=} \triangleq \{\boldsymbol{y} \in \mathbb{R}^n | \boldsymbol{C}\boldsymbol{y} = \boldsymbol{d}\}$; thus the *relative* interior of $\mathcal{C}$ is $int\,\mathcal{C} \triangleq \mathcal{C}_{<} \cap \mathcal{C}_{=}$, and we consider $int\,\mathcal{C} \neq \emptyset$ and $\mathcal{C}$ is compact.

Definition 1 ((strong/ξ) monotonicity) *An operator $F : \mathcal{X} \supseteq \mathcal{S} \to \mathbb{R}^n$ is* monotone *on $\mathcal{S}$ if: $\langle \boldsymbol{x} - \boldsymbol{x}', F(\boldsymbol{x}) - F(\boldsymbol{x}')\rangle \geq 0, \forall \boldsymbol{x}, \boldsymbol{x}' \in \mathcal{S}$. F is said to be ξ–*monotone *on $\mathcal{S}$ iff there exist $c > 0$ and $\xi > 1$ such that $\langle \boldsymbol{x} - \boldsymbol{x}', F(\boldsymbol{x}) - F(\boldsymbol{x}')\rangle \geq c\|\boldsymbol{x} - \boldsymbol{x}'\|^{\xi}$, for all $\boldsymbol{x}, \boldsymbol{x}' \in \mathcal{S}$. Finally, F is μ-*strongly monotone *on $\mathcal{S}$ if there exists $\mu > 0$, such that $\langle \boldsymbol{x} - \boldsymbol{x}', F(\boldsymbol{x}) - F(\boldsymbol{x}')\rangle \geq \mu\|\boldsymbol{x} - \boldsymbol{x}'\|^2$, for all $\boldsymbol{x}, \boldsymbol{x}' \in \mathcal{S}$. Moreover, we say that an operator F is* star–monotone, star–ξ-monotone *or* star–strongly-monotone *(on $\mathcal{S}$) if the respective definition holds for $\boldsymbol{x}' \equiv \boldsymbol{x}^\star$, where $\boldsymbol{x}^\star \in \mathcal{S}^\star_{\mathcal{S},F}$.*

Assumption 1 (first-order smoothness) *Let $F : \mathcal{X} \supseteq \mathcal{S} \to \mathbb{R}^n$ be an operator, we say that F satisfies L-first-order smoothness on $\mathcal{S}$, or L-smoothness, if F is an L-Lipschitz map; that is, there exists $L > 0$ such that $\|F(\boldsymbol{x}) - F(\boldsymbol{x}')\| \leq L\,\|\boldsymbol{x} - \boldsymbol{x}'\|$, for all $\boldsymbol{x}, \boldsymbol{x}' \in \mathcal{S}$.*

App. B.3 provides background on path-following IP methods and describes the ADMM method [30, 32, 33, 52]; whereas App. B.2 summarizes the solution existence guarantees for cVI.

3. *ACVI*: first-order ADMM-based IP method for constrained VIs

Algorithm 1 describes the ACVI method, derived in App. C. See App. D.5 for further discussion. In the remaining discussion, where clear from context, we drop the superscript from the iterate $\boldsymbol{x}_k^{(t)}$.

We consider two broad problem classes, see App D.1 for a discussion. For $r, s > 0$, let $\hat{\mathcal{C}}_r \triangleq \{\boldsymbol{x} \in \mathbb{R}^n | \boldsymbol{C}\boldsymbol{x} = \boldsymbol{d}, \varphi(\boldsymbol{x}) \leq r\boldsymbol{e}\}$, and let $\tilde{\mathcal{C}}_s \triangleq \{\boldsymbol{x} \in \mathbb{R}^n |\, \|\boldsymbol{C}\boldsymbol{x} - \boldsymbol{d}\| \leq s, \varphi(\boldsymbol{x}) \leq \boldsymbol{0}\}$.

Algorithm 1 ACVI pseudocode.

Input: operator $F : \mathcal{X} \to \mathbb{R}^n$, constraints $\boldsymbol{C}\boldsymbol{x} = \boldsymbol{d}$ and $\varphi_i(\boldsymbol{x}) \leq 0, i = [m]$, hyperparameters $\mu_{-1}, \beta > 0, \delta \in (0,1)$, number of outer and inner loop iterations T and K, resp.
Initialize: $\boldsymbol{y}_0^{(0)} \in \mathbb{R}^n$, $\boldsymbol{\lambda}_0^{(0)} \in \mathbb{R}^n$
$\boldsymbol{P}_c \triangleq \boldsymbol{I} - \boldsymbol{C}^\intercal(\boldsymbol{C}\boldsymbol{C}^\intercal)^{-1}\boldsymbol{C}$; $\boldsymbol{d}_c \triangleq \boldsymbol{C}^\intercal(\boldsymbol{C}\boldsymbol{C}^\intercal)^{-1}\boldsymbol{d}$ where $\boldsymbol{P}_c \in \mathbb{R}^{n\times n}$, $\boldsymbol{d}_c \in \mathbb{R}^n$
for $t = 0, \ldots, T-1$ **do**
  $\mu_t = \delta\mu_{t-1}$
  for $k = 0, \ldots, K-1$ **do**
    Set $\boldsymbol{x}_{k+1}^{(t)}$ to be the solution of: $\boldsymbol{x} + \frac{1}{\beta}\boldsymbol{P}_c F(\boldsymbol{x}) - \boldsymbol{P}_c\boldsymbol{y}_k^{(t)} + \frac{1}{\beta}\boldsymbol{P}_c\boldsymbol{\lambda}_k^{(t)} - \boldsymbol{d}_c = \boldsymbol{0}$ (w.r.t. $\boldsymbol{x}$)
    $\boldsymbol{y}_{k+1} = \underset{\boldsymbol{y}}{argmin} - \mu\sum_{i=1}^m \log\big(-\varphi_i(\boldsymbol{y})\big) + \frac{\beta}{2}\left\|\boldsymbol{y} - \boldsymbol{x}_{k+1} - \frac{1}{\beta}\boldsymbol{\lambda}_k\right\|^2$
    $\boldsymbol{\lambda}_{k+1} = \boldsymbol{\lambda}_k + \beta(\boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1})$
  end for
  $(\boldsymbol{y}_0^{(t+1)}, \boldsymbol{\lambda}_0^{(t+1)}) \triangleq (\boldsymbol{y}_K^{(t)}, \boldsymbol{\lambda}_K^{(t)})$
end for

Theorem 1 (Last and average iterate convergence for star-ξ-monotone operator) *Given an operator $F : \mathcal{X} \to \mathbb{R}^n$ monotone on $\mathcal{C}_=$ (Def. 1), assume that either F is strictly monotone on $\mathcal{C}$ or one of φ_i is strictly convex. Assume there exists $r > 0$ or $s > 0$ such that F is star-ξ-monotone on either $\hat{\mathcal{C}}_r$ or $\tilde{\mathcal{C}}_s$. Let $\boldsymbol{x}_K^{(t)}$ and $\hat{\boldsymbol{x}}_K^{(t)} \triangleq \frac{1}{K}\sum_{k=1}^K \boldsymbol{x}_k^{(t)}$ denote the last and average iterate of Algorithm 1, respectively, run with sufficiently small μ_{-1}. Then for all $t \in [T]$, we have: $\left\|\boldsymbol{x}_K^{(t)} - \boldsymbol{x}^\star\right\| \leq \mathcal{O}(\frac{1}{K^{1/(2\xi)}})$.*

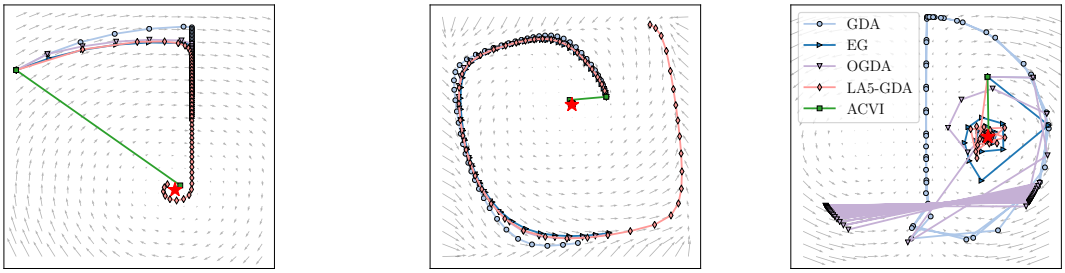
1. *If in addition F is ξ-monotone on $C_=$, we have $\left\|\hat{\boldsymbol{x}}_K^{(t)} - \boldsymbol{x}^\star\right\| \leq \mathcal{O}(\frac{1}{K^{1/\xi}})$.*
2. *Moreover, if F is L-Lipschitz on $\mathcal{C}_=$ we have: $\mathcal{G}(\boldsymbol{x}_K^{(t)}, \mathcal{C}) \leq \mathcal{O}(\frac{L}{K^{1/(2\xi)}})$; $\mathcal{G}(\hat{\boldsymbol{x}}_K^{(t)}, \mathcal{C}) \leq \mathcal{O}(\frac{L}{K^{1/\xi}})$.*

Theorem 2 (Last and average iterate convergence for monotone operator) *Given an operator $F : \mathcal{X} \to \mathbb{R}^n$, assume (i) F is monotone on $\mathcal{C}_=$, and (ii) either F is strictly monotone on $\mathcal{C}$ or one of φ_i is strictly convex, and (iii) $\underset{\boldsymbol{x}\in S\setminus\{\boldsymbol{x}^\star\}}{inf} F(\boldsymbol{x})^\intercal \frac{\boldsymbol{x}-\boldsymbol{x}^\star}{\|\boldsymbol{x}-\boldsymbol{x}^\star\|} > 0$, where $S \equiv \hat{\mathcal{C}}_r$ or $\tilde{\mathcal{C}}_s$. Let $\boldsymbol{x}_K^{(t)}$ and $\hat{\boldsymbol{x}}_K^{(t)} \triangleq \frac{1}{K}\sum_{k=1}^K \boldsymbol{x}_k^{(t)}$ denote the last and average iterate of Algorithm 1, respectively, run with sufficiently small μ_{-1}. Then for all $t \in [T]$, we have that: $\left\|\boldsymbol{x}_K^{(t)} - \boldsymbol{x}^\star\right\| \leq \mathcal{O}(\frac{1}{\sqrt{K}})$.*

1. *If in addition $\underset{\boldsymbol{x}\in S\setminus\{\boldsymbol{x}^\star\}}{inf} F(\boldsymbol{x}^\star)^\intercal \frac{\boldsymbol{x}-\boldsymbol{x}^\star}{\|\boldsymbol{x}-\boldsymbol{x}^\star\|} > 0$ (with $S \equiv \hat{\mathcal{C}}_r$ or $\tilde{\mathcal{C}}_s$), then $\left\|\hat{\boldsymbol{x}}_K^{(t)} - \boldsymbol{x}^\star\right\| \leq \mathcal{O}(\frac{1}{K})$.*
2. *Moreover, if F is L-Lipschitz on $\mathcal{C}_=$ we have: $\mathcal{G}(\boldsymbol{x}_K^{(t)}, \mathcal{C}) \leq \mathcal{O}(\frac{L}{\sqrt{K}})$ and $\mathcal{G}(\hat{\boldsymbol{x}}_K^{(t)}, \mathcal{C}) \leq \mathcal{O}(\frac{L}{K})$.*

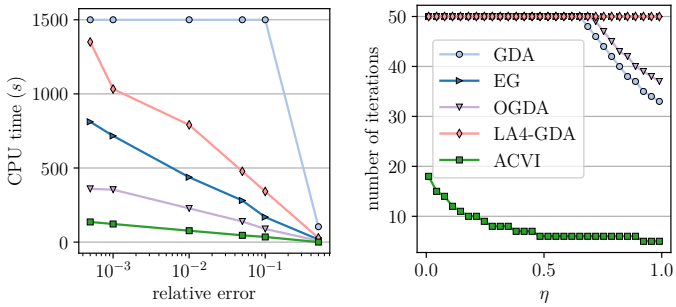
4. Experiments

Problems. We use the following 2D problems: (i) *cBG*: the bilinear game, constrained on $\mathbb{R}_+$, stated in Fig. 1, (ii) *Von Neumann's ratio game* [20, 23, 81], (iii) *Forsaken* game [43]–which exhibits *limit cycles*, as well as (iv) *toy GAN*—used in [1, 19]. We also consider a constrained BG on the probability simplex, with $n = 1000$, given in (HBG). Finally, we use GANs on MNIST [49] with added linear inequalities, see App. F.3. App. G lists additional experiments, including on Fashion-MNIST [86].
Methods. We compare with the projection-based methods (described in App. B.5) (i) **GDA**, (ii) **EG** [46], (iii) **OGDA** [64], and (iv) **LA$\tilde{k}$-GDA** [14, 87], where $\tilde{k}$ is the hyperparameter of LA. For ACVI on MNIST, l denotes the number of steps to solve the subproblems; see Alg. 3.



(*a*) Von Neumann's ratio game (*b*) Forsaken game (*c*) toy GAN

Figure 2: **Convergence of GDA, EG, OGDA, LA-GDA, and ACVI on three different 2d problems**, for a *fixed* number of *total* iterations, where markers denote the *iterates* of the respective method.



(*a*) fixed error (*b*) rotational intensity

Figure 3: **Comparison on** (HBG)**: (a)**–CPU time given fixed error, **(b)**–number of iterations needed to reach ϵ-distance to solution for varying intensity of the rotational component $(1-\eta)$. We set maximum iterations/time to run.

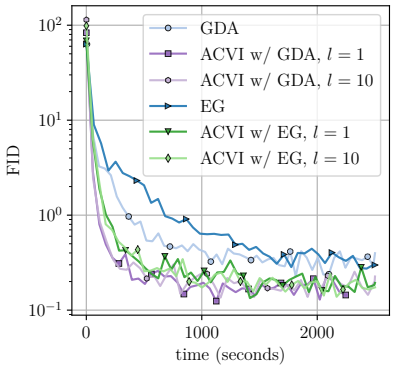


Figure 4: FID (lower is better) on **MNIST** with added constraints, over wall-clock time; averaged over 3 seeds. See also § 4 and App. F.

Results. From Fig. 1 we observe that projection-based algorithms may zigzag when hitting a constraint, due to the rotational nature of F, behavior that ACVI avoids because it incorporates the constraints in its update rule; see Fig. 5 for the remaining baselines. Fig. 2 shows that even with problems that go beyond our theoretical assumptions a single step of ACVI significantly reduces the distance to the solution. Moreover, from Fig. 2(*b*)subfigure we observe that ACVI escapes the limit cycles; see also Fig. 6. Fig. 3 shows results for (HBG), indicating that ACVI is time efficient, and that ACVI performs well relative to projection-based methods for varying rotational intensity $(1-\eta)$. Fig. 4 summarizes the experiments on MNIST with linear inequality constraints; we observe that ACVI converges significantly faster than the corresponding baseline.

5. Conclusion

Motivated by the lack of a *first*-order method to solve cVIs with general constraints, we proposed a framework that combines *(i) interior-point* methods—needed to be able to handle general constraints—with *(ii)* the ADMM method—designed to deal with *separable* objectives. The combination yields *ACVI*—a first-order ADMM-based interior point method for cVIs. We proved convergence for two broad classes of problems, and derived the corresponding convergence rates. Numerical experiments showed that while projection-based methods zigzag when hitting a constraint due to the rotational vector field, ACVI avoids this by incorporating the constraints in the update rule.

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Appendix A. Related Work

Unconstrained VIs: methods and guarantees. Apart from the standard gradient descent ascent (GDA) method, among the most commonly used methods for VI optimization are the extragradient method [EG, 46], optimistic GDA [OGDA, 64], and the lookahead method [LA, 14, 87]. (See App. B for a full description). In contrast to gradient fields (as in a single-objective setting), when F is a general vector field, the last iterate can be far from the solution even though the average iterate converges to it [12, 19]. This is problematic since it implies that the average convergence guarantee is weaker in the sense that it may not extend to more general setups where we can no longer rely on the convexity of $\mathcal{X}$.

Golowich et al. [35, 36] provided a last-iterate lower bound of $\mathcal{O}(\frac{1}{\tilde{p}\sqrt{K}})$ for the broad class of *$\tilde{p}$-stationary canonical linear iterative* ($\tilde{p}$-SCLI) first-order methods [2]. An extensive line of further work has provided guarantees for the last iterate for other problem classes. For the general MVI class, the following $\tilde{p}$-SCLI methods come with guarantees that match the lower bound: (i) Golowich et al. [36] obtained a rate in terms of the gap function relying on first- and second-order smoothness of F, and Gorbunov et al. [38] obtained a rate of $\mathcal{O}(\frac{1}{K})$ in terms of reducing the squared norm of the operator relying on first-order smoothness of F (Assumption 1), using a computer-assisted proof, and (ii) Golowich et al. [36] and Chavdarova et al. [13] provided the best iterate rate for OGDA.

Constrained zero-sum and VI classes of problems. Gidel et al. [31] extended the Frank-Wolfe [28, 45, 48] method—also known as the *conditional gradient* [79]—to solve a subclass of cVI, specifically constrained zero-sum problems. This extension was carried out under a strong convex-concavity assumption and also under the assumption that the constraint set is *strongly* convex; that is, it has sublevel sets that are strongly convex functions [80]. Daskalakis and Panageas [18] provided an asymptotic proof for the last iterate for zero-sum convex-concave constrained problems for the *optimistic multiplicative weights update* (OMWU) method. Wei et al. [83] focused on OGDA and OMWU in the constrained setting and provided convergence rates for bilinear games over the simplex. In her seminal work, Korpelevich [46] proposed the classical (projected) extragradient method (EG)—see App. B—and proved its convergence for monotone (c)VIs with an L-Lipschitz operator, and, as mentioned above, Cai et al. [11] established a rate with respect to the gap function using a computer-aided proof. For the constrained setting, Tseng [78] built on [63] and provided a linear convergence rate for EG in the setting of strongly monotone F, whereas Malitsky [53] focused on the same setting but on the projected reflected gradient method. Diakonikolas [22] obtained parameter-free guarantee for Halpern iteration [39] for cocoercive operators. Goffin et al. [34] described a second-order cutting-plane method for solving pseudomonotone VIs with linear inequalities.

Interior point (IP) methods in single-objective and VI settings. Traditionally IP methods primarily express the inequality constraints by augmenting the objective with a log-barrier penalty (see § B.3), and then use Newton's method to solve the subproblem [9]. The latter involves computing either the inverse of a large matrix or a Cholesky decomposition, and yet it can be highly efficient in low dimensions as it requires only a few iterations to converge. When the dimensionality of the variable is large, however, the computation becomes infeasible. Among other IP variants that address this issue, Lin et al. [50] replaced the Newton step with the ADMM method, which is known to be highly scalable in terms of the dimension [10]. In the context of cVIs, a few works apply IP methods, mostly Newton-based [e.g., 59, Chapter 7]. Monteiro and Pang [58] analyze path-following IP methods for complementarity problems, which are a subclass of cVI, using local homeomorphic maps. Chen

et al. [15] provided a superlinear global convergence rate of the smoothing Newton method when F is semi-smooth for *box constrained* VIs. Similarly, Qi and Sun [65], Qi et al. [66] focused on the smoothing Newton method and provided the rate for the outer loop. Ralph and Wright [68] showed superlinear convergence for MVI problems under inequality constraints, under the following set of assumptions: (i) existence of a *strictly* complementary solution, (ii) full rank of the Jacobian of the active constraints at the solution, and (iii) twice differentiable constraints. They provided a *local* convergence rate. Fan and Yan [27] considered inequality constraints and proposed a second-order Newton-based method that has global convergence guarantees under certain conditions. A rate was not provided.

Appendix B. Background: Additional Details

This section lists additional background such as omitted definitions, and a description of ADMM and the used baseline methods.

B.1. Additional VI definitions & equivalent formulations

The analog for cVI of the function values used as a performance measure for convergence rates in convex optimization is the *gap function* (a.k.a., the *optimality* gap or *primal* gap), defined next.

Definition 2 (gap function) *Given a candidate point $\boldsymbol{x}' \in \mathcal{X}$ and a map $F : \mathcal{X} \supseteq \mathcal{S} \to \mathbb{R}^n$ where $\mathcal{S}$ is compact, the gap function $\mathcal{G} : \mathbb{R}^n \to \mathbb{R}$ is defined as $\mathcal{G}(\boldsymbol{x}', \mathcal{S}) = \max_{\boldsymbol{x} \in \mathcal{S}} \langle F(\boldsymbol{x}'), \boldsymbol{x}' - \boldsymbol{x} \rangle$.*

Note that the gap function requires $\mathcal{S}$ to be compact in order to be defined (as otherwise, it can be infinite).

Seeing an operator $F : \mathcal{X} \to \mathbb{R}^n$ as the graph $GrF = \{(\boldsymbol{x}, \boldsymbol{y}) | \boldsymbol{x} \in \mathcal{X}, \boldsymbol{y} = F(\boldsymbol{x})\}$, its inverse F^{-1} is defined as $GrF^{-1} = \{(\boldsymbol{y}, \boldsymbol{x}) | (\boldsymbol{x}, \boldsymbol{y}) \in GrF\}$. See for example [72] for further discussion. We denote the projection to the set $\mathcal{X}$ with $\Pi_{\mathcal{X}}$.

Definition 3 ($\frac{1}{\mu}$-cocoercive operator) *An operator $F : \mathcal{X} \supseteq \mathcal{S} \to \mathbb{R}^n$ is* $\frac{1}{\mu}$-cocoercive *(or* $\frac{1}{\mu}$-inverse strongly monotone*) on $\mathcal{S}$ if its inverse (graph) F^{-1} is μ-strongly monotone on $\mathcal{S}$, that is,*

$$\exists \mu > 0, \quad s.t. \quad \langle \boldsymbol{x} - \boldsymbol{x}', F(\boldsymbol{x}) - F(\boldsymbol{x}') \rangle \geq \mu \left\| F(\boldsymbol{x}) - F(\boldsymbol{x}') \right\|^2, \forall \boldsymbol{x}, \boldsymbol{x}' \in \mathcal{S} .$$

It is star $\frac{1}{\mu}$-cocoercive *if the above holds when setting $\boldsymbol{x}' \equiv \boldsymbol{x}^\star$ where $\boldsymbol{x}^\star$ denotes a solution, that is:*

$$\exists \mu > 0, \quad s.t. \quad \langle \boldsymbol{x} - \boldsymbol{x}^\star, F(\boldsymbol{x}) - F(\boldsymbol{x}^\star) \rangle \geq \mu \left\| F(\boldsymbol{x}) - F(\boldsymbol{x}^\star) \right\|^2, \forall \boldsymbol{x} \in \mathcal{S}, \boldsymbol{x}^\star \in \mathcal{S}^\star_{\mathcal{X},F} .$$

Note from Def. 3 that cocoercivity is a strict subclass of monotone and L-Lipschitz operators, thus is it is a stronger assumption. See Chapter 4.2 of [6] for further relations of cocoercivity with other properties of operators.

In the following, we will make use of the natural and normal mappings of an operator $F : \mathcal{X} \to \mathbb{R}^n$, where $\mathcal{X} \subset \mathbb{R}^n$. Following the notation of [26], the natural map $F^{NAT}_{\mathcal{X}} : \mathcal{X} \to \mathbb{R}^n$ is:

$$F^{NAT}_{\mathcal{X}} \triangleq \boldsymbol{x} - \Pi_{\mathcal{X}}\big(\boldsymbol{x} - F(\boldsymbol{x})\big), \qquad \forall \boldsymbol{x} \in \mathcal{X} , \tag{F-NAT}$$

whereas the normal map $F^{NOR}_{\mathcal{X}} : \mathbb{R}^n \to \mathbb{R}^n$ is:

$$F^{NOR}_{\mathcal{X}} \triangleq F\big(\Pi_{\mathcal{X}}(\boldsymbol{x})\big) + \boldsymbol{x} - \Pi_{\mathcal{X}}(\boldsymbol{x}), \qquad \forall \boldsymbol{x} \in \mathbb{R}^n . \tag{F-NOR}$$

Moreover, we have the following solution characterizations:

1. $\boldsymbol{x}^\star \in \mathcal{S}^\star_{\mathcal{X},F}$ iff $F^{NAT}_{\mathcal{X}}(\boldsymbol{x}^\star) = \mathbf{0}$, and
2. $\boldsymbol{x}^\star \in \mathcal{S}^\star_{\mathcal{X},F}$ iff $\exists \boldsymbol{x}' \in R^n$ s.t. $\boldsymbol{x}^\star = \Pi_{\mathcal{X}}(\boldsymbol{x}')$ and $F^{NOR}_{\mathcal{X}}(\boldsymbol{x}') = \mathbf{0}$.

B.2. Existence of solution

We provide brief informal summary of some sufficient conditions for solution existence, that $\mathcal{S}^\star_{\cdot,F} \neq \emptyset$. See Chapter 2 of [26] for a full treatment of the topic.

The common underlying tool to establish that a solution to the VI($\mathcal{X}$, F) problem exists is using *topological degree*. The topological degree tool is *designed* so as to satisfy the so-called *homotopy invariance* axiom, which in turn allows for reducing a solution existence question of a complicated map to a simpler one (which is homotopy-invariant to the original one) for which we can more easily show that it has a solution (e.g., the identity map on a closed domain). It can be used (as one way) to prove the celebrated *Brouwer fixed-point theorem*, which states that any continuous map $\Phi : \mathcal{S} \to \mathcal{S}$, where $\mathcal{S}$ is nonempty convex compact set, has a fixed point in $\mathcal{S}$.

We have the following sufficient condition, see [Cor. 2.2.5, 26]

Theorem 3 (sufficient condition for existence of the solution, Cor. 2.2.5, [26]) *If $\mathcal{X} \subseteq \mathbb{R}^n$ is compact and convex, and $F : \mathcal{X} \to \mathbb{R}^n$ is continuous, then the solution set is nonempty and compact.*

It can also be shown that when $\mathcal{X}$ is closed convex (and F continuous), if one can find $\boldsymbol{x}' \in \mathcal{X}$ s.t. $\langle F(\boldsymbol{x}), \boldsymbol{x} - \boldsymbol{x}' \rangle \geq 0, \forall \boldsymbol{x} \in \mathcal{X}$, then $\mathcal{S}^\star_{\mathcal{X},F} \neq \emptyset$. The same conclusion follows if one can show that $\exists \boldsymbol{x}' \in \mathcal{X}$ and the set $\{\boldsymbol{x} \in \mathcal{X} | \langle F(\boldsymbol{x}), \boldsymbol{x} - \boldsymbol{x}' \rangle < 0\}$ is bounded (possibly empty, see Prop. 2.2.3 in [26]).

Sufficient conditions can also be established via the natural and the normal map due to the above solution characterizations. In this case, we require that F is continuous on an open set $\mathcal{S}$ and we are interested if $\mathcal{S}^\star_{\mathcal{X},F} \neq \emptyset$, where $\mathcal{X}$ is assumed closed and convex and subset of $\mathcal{S}$, $\mathcal{X} \subseteq \mathcal{S}$. If one establishes that a solution exists for $F^{NAT}_{\mathcal{X}}$ on a bounded open set $\mathcal{U}$, and if $cl\, \mathcal{U} \subseteq \mathcal{S}$, then it follows that $\mathcal{S}^\star_{\mathcal{X},F} \neq \emptyset$. A similar implication holds when we have such a guarantee for $F^{NOR}_{\mathcal{X}}$. See Theorem 2.2.1 of [26].

In summary, the solution existence guarantee follows from the boundness of some set which includes $\mathcal{X}$, the boundness of the set of potential solutions (if we can construct such set), or the compactness of $\mathcal{X}$ itself.

B.3. Relevant path-following interior-point methods and ADMM

In this section, we overview the interior-point approach to single-objective optimization, focusing on aspects that are most relevant to our proposed method. Consider the following problem:

$$\min_{\boldsymbol{x}} f(\boldsymbol{x}) \qquad s.t. \quad \varphi(\boldsymbol{x}) \leq \boldsymbol{0} \quad and \quad \boldsymbol{C}\boldsymbol{x} = \boldsymbol{d}\,, \tag{cCVX}$$

where $f, \varphi_i : \mathbb{R}^n \to \mathbb{R}$ are convex and continuously differentiable, $\boldsymbol{x} \in \mathbb{R}^n$, $\boldsymbol{C} \in \mathbb{R}^{p \times n}$, and $\boldsymbol{d} \in \mathbb{R}^p$. IP methods solve problem (cCVX) by reducing it to a sequence of linear equality-constrained problems via a logarithmic barrier [see, e.g., 9, Chapter 11]:

$$\min_{\boldsymbol{x}} f(\boldsymbol{x}) - \mu \sum_{i=1}^{m} \log(-\varphi_i(\boldsymbol{x})) \quad s.t. \quad \boldsymbol{C}\boldsymbol{x} = \boldsymbol{d}, \qquad \text{with} \quad \mu > 0\,. \tag{l-cCVX}$$

Assume that (l-cCVX) has a solution for each $\mu > 0$, and let $\boldsymbol{x}^\mu$ denote the solution of (l-cCVX) for a given μ. The *central path* of (l-cCVX) is defined as the set of points $\boldsymbol{x}^\mu, \mu > 0$. Note that $\boldsymbol{x}^\mu \in \mathbb{R}^n$ is a strictly feasible point of (cCVX) as it satisfies $\varphi(\boldsymbol{x}^\mu) < \boldsymbol{0}$ and $\boldsymbol{C}\boldsymbol{x}^\mu = \boldsymbol{d}$.

Alternating direction method of multipliers (ADMM) method. ADMM [30, 32, 33, 52] is a gradient-based algorithm for convex optimization problems that splits the objective into subproblems each of which is easier to solve. Its popularity is due to its computational scalability [10]. Consider a problem of the following form:

$$\min_{\boldsymbol{x},\boldsymbol{y}} f(\boldsymbol{x}) + g(\boldsymbol{y}) \quad s.t. \quad \boldsymbol{A}\boldsymbol{x} + \boldsymbol{B}\boldsymbol{y} = \boldsymbol{b}\,, \tag{ADMM-Pr}$$

where $f, g : \mathbb{R}^n \to \mathbb{R}$ are convex, $\boldsymbol{x}, \boldsymbol{y} \in \mathbb{R}^n$, $\boldsymbol{A}, \boldsymbol{B} \in \mathbb{R}^{n' \times n}$, and $\boldsymbol{b} \in \mathbb{R}^{n'}$. The augmented Lagrangian function, $\mathcal{L}_\beta(\cdot)$, of the (ADMM-Pr) problem is:

$$\mathcal{L}_\beta(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{\lambda}) = f(\boldsymbol{x}) + g(\boldsymbol{y}) + \langle \boldsymbol{x} - \boldsymbol{y}, \boldsymbol{\lambda} \rangle + \frac{\beta}{2}\|\boldsymbol{x} - \boldsymbol{y}\|^2, \tag{AL-CVX}$$

where $\beta > 0$ is referred to as the *penalty parameter*. If the augmented Lagrangian method is used to solve (AL-CVX), at each step k we have: $\boldsymbol{x}_{k+1}, \boldsymbol{y}_{k+1} = \underset{\boldsymbol{x},\boldsymbol{y}}{argmin}\ \mathcal{L}_\beta(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{\lambda}_k)$ and $\boldsymbol{\lambda}_{k+1} = \boldsymbol{\lambda}_k + \beta(\boldsymbol{A}\boldsymbol{x}_{k+1} + \boldsymbol{B}\boldsymbol{y}_{k+1} - \boldsymbol{b})$, where the latter step is gradient ascent on the dual. In contrast, ADMM updates $\boldsymbol{x}$ and $\boldsymbol{y}$ in an alternating way as follows:

$$\begin{aligned} \boldsymbol{x}_{k+1} &= \underset{\boldsymbol{x}}{argmin}\ \mathcal{L}_\beta(\boldsymbol{x}, \boldsymbol{y}_k, \boldsymbol{\lambda}_k)\,, \\ \boldsymbol{y}_{k+1} &= \underset{\boldsymbol{y}}{argmin}\ \mathcal{L}_\beta(\boldsymbol{x}_{k+1}, \boldsymbol{y}_k, \boldsymbol{\lambda}_k)\,, \\ \boldsymbol{\lambda}_{k+1} &= \boldsymbol{\lambda}_k + \beta(\boldsymbol{A}\boldsymbol{x}_{k+1} + \boldsymbol{B}\boldsymbol{y}_{k+1} - \boldsymbol{b})\,. \end{aligned} \tag{ADMM}$$

B.4. Existence of Central path

In this section, we discuss the results that establish guarantees of the existence of the central path. Let $L(\boldsymbol{x}, \boldsymbol{\lambda}, \boldsymbol{\nu}) \triangleq F(\boldsymbol{x}) + \nabla\varphi^\intercal(\boldsymbol{x})\boldsymbol{\lambda} + \boldsymbol{C}^\intercal\boldsymbol{\nu}$, $h(\boldsymbol{x}) = \boldsymbol{C}^\intercal\boldsymbol{x} - \boldsymbol{d}$. For $(\boldsymbol{\lambda}, \boldsymbol{w}, \boldsymbol{x}, \boldsymbol{\nu}) \in \mathbb{R}^{2m+n+p}$, let

$$G(\boldsymbol{\lambda}, \boldsymbol{w}, \boldsymbol{x}, \boldsymbol{\nu}) \triangleq \begin{pmatrix} \boldsymbol{w} \circ \boldsymbol{\lambda} \\ \boldsymbol{w} + \varphi(\boldsymbol{x}) \\ L(\boldsymbol{x}, \boldsymbol{\lambda}, \boldsymbol{\nu}) \\ h(\boldsymbol{x}) \end{pmatrix} \in \mathbb{R}^{2m+n+p},$$

and

$$H(\boldsymbol{\lambda}, \boldsymbol{w}, \boldsymbol{x}, \boldsymbol{\nu}) \triangleq \begin{pmatrix} \boldsymbol{w} + \varphi(\boldsymbol{x}) \\ L(\boldsymbol{x}, \boldsymbol{\lambda}, \boldsymbol{\nu}) \\ h(\boldsymbol{x}) \end{pmatrix} \in \mathbb{R}^{m+n+p}.$$

Let $H_{++} \triangleq H(\mathbb{R}^{2m}_{++} \times \mathbb{R}^n \times \mathbb{R}^p)$. By [Corollary 11.4.24, 26] we have the following proposition.

Proposition 1 (sufficient condition for the existence of the central path.) *If F is monotone, either F is strictly monotone or one of φ_i is strictly convex, and $\mathcal{C}$ is bounded. The following four statements hold for the functions G and H:*

1. *G maps $\mathbb{R}^{2m}_{++} \times \mathbb{R}^{n+p}$ homeomorphically onto $\mathbb{R}^m_{++} \times H_{++}$;*
2. *$\mathbb{R}^m_{++} \times H_{++} \subseteq G(\mathbb{R}^{2m}_{+} \times \mathbb{R}^{n+p})$;*

3. *for every vector* $\boldsymbol{a} \in \mathbb{R}^m_+$*, the system*

$$H(\boldsymbol{\lambda}, \boldsymbol{w}, \boldsymbol{x}, \boldsymbol{\nu}) = \mathbf{0}, \quad \boldsymbol{w} \circ \boldsymbol{\lambda} = \boldsymbol{a}$$

   has a solution $(\boldsymbol{\lambda}, \boldsymbol{w}, \boldsymbol{x}, \boldsymbol{\nu}) \in \mathbb{R}^{2m}_+ \times \mathbb{R}^{n+p}$*; and*

4. *the set* H_{++} *is convex.*

B.5. Saddle-point optimization methods

In this section, we describe in detail the saddle point methods that we compare with in the main paper in § 4. We denote the projection to the set $\mathcal{X}$ with $\Pi_{\mathcal{X}}$, and when the method is applied in the unconstrained setting $\Pi_{\mathcal{X}} \equiv \boldsymbol{I}$.

For an example of the associated vector field and its Jacobian, consider the following constrained zero-sum game:

$$\min_{\boldsymbol{x}_1 \in \mathcal{X}_1} \max_{\boldsymbol{x}_2 \in \mathcal{X}_2} f(\boldsymbol{x}_1, \boldsymbol{x}_2)\,, \tag{ZS-G}$$

where $f : \mathcal{X}_1 \times \mathcal{X}_2 \to \mathbb{R}$ is smooth and convex in $\boldsymbol{x}_1$ and concave in $\boldsymbol{x}_2$. As in the main paper, we write $\boldsymbol{x} \triangleq (\boldsymbol{x}_1, \boldsymbol{x}_2) \in \mathbb{R}^n$. The vector field $F : \mathcal{X} \to \mathbb{R}^n$ and its Jacobian J are defined as:

$$F(\boldsymbol{x}) = \begin{bmatrix} \nabla_{\boldsymbol{x}_1} f(\boldsymbol{x}) \\ -\nabla_{\boldsymbol{x}_2} f(\boldsymbol{x}) \end{bmatrix}, \qquad J(\boldsymbol{x}) = \begin{bmatrix} \nabla^2_{\boldsymbol{x}_1} f(\boldsymbol{x}) & \nabla_{\boldsymbol{x}_2} \nabla_{\boldsymbol{x}_1} f(\boldsymbol{x}) \\ -\nabla_{\boldsymbol{x}_1} \nabla_{\boldsymbol{x}_2} f(\boldsymbol{x}) & -\nabla^2_{\boldsymbol{x}_2} f(\boldsymbol{x}) \end{bmatrix}.$$

In the remaining of this section, we will only refer to the joint variable $\boldsymbol{x}$, and (with abuse of notation) the subscript will denote the step. Let $\gamma \in [0, 1]$ denote the step size.

(Projected) Gradient Descent Ascent (GDA). The extension of gradient descent for the cVI problem is *gradient descent ascent* (GDA). The GDA update at step k is then:

$$\boldsymbol{x}_{k+1} = \Pi_{\mathcal{X}}\big(\boldsymbol{x}_k - \gamma F(\boldsymbol{x}_k)\big)\,. \tag{GDA}$$

(Projected) Extragradient (EG). EG [46] uses a "prediction" step to obtain an extrapolated point $\boldsymbol{x}_{k+\frac{1}{2}}$ using GDA: $\boldsymbol{x}_{k+\frac{1}{2}} = \Pi_{\mathcal{X}}\big(\boldsymbol{x}_k - \gamma F(\boldsymbol{x}_k)\big)$, and the gradients at the *extrapolated* point are then applied to the *current* iterate $\boldsymbol{x}_t$:

$$\boldsymbol{x}_{k+1} = \Pi_{\mathcal{X}}\Big(\boldsymbol{x}_k - \gamma F\Big(\Pi_{\mathcal{X}}\big(\boldsymbol{x}_k - \gamma F(\boldsymbol{x}_k)\big)\Big)\Big)\,. \tag{EG}$$

In the original EG paper, [46] proved that the EG method (with a fixed step size) converges for monotone VIs, as follows.

Theorem 4 (Korpelevich [46]) *Given a map* $F : \mathcal{X} \mapsto \mathbb{R}^n$*, if the following is satisfied:*

1. *the set* $\mathcal{X}$ *is closed and convex,*

2. F *is single-valued, definite, and monotone on* $\mathcal{X}$*–as per Def. 1,*

3. F *is L-Lipschitz–as per Asm. 1.*

then there exists a solution $\boldsymbol{x}^\star \in \mathcal{X}$*, such that the iterates* $\boldsymbol{x}_k$ *produced by the EG update rule with a fixed step size* $\gamma \in (0, \frac{1}{L})$ *converge to it, that is* $\boldsymbol{x}_k \to \boldsymbol{x}^\star$*, as* $k \to \infty$*.*

Facchinei and Pang [26] also show that for any *convex-concave* function f and any closed convex sets $\boldsymbol{x}_1 \in \mathcal{X}_1$ and $\boldsymbol{x}_2 \in \mathcal{X}_2$, the EG method converges [26, Theorem 12.1.11].

(Projected) Optimistic Gradient Descent Ascent (OGDA). The update rule of Optimistic Gradient Descent Ascent OGDA [(OGDA) 64] is:

$$\boldsymbol{x}_{n+1} = \Pi_{\mathcal{X}}\big(\boldsymbol{x}_n - 2\gamma F(\boldsymbol{x}_n) + \gamma F(\boldsymbol{x}_{n-1})\big)\,. \tag{OGDA}$$

(Projected) Lookahead–Minmax (LA). The LA algorithm for min-max optimization [14], originally proposed for minimization by [87], is a general wrapper of a "base" optimizer where at every step t: (i) a copy of the current iterate $\tilde{\boldsymbol{x}}_n$ is made: $\tilde{\boldsymbol{x}}_n \leftarrow \boldsymbol{x}_n$, (ii) $\tilde{\boldsymbol{x}}_n$ is updated $k \geq 1$ times, yielding $\tilde{\boldsymbol{\omega}}_{n+k}$, and finally (iii) the actual update $\boldsymbol{x}_{n+1}$ is obtained as a *point that lies on a line between* the current $\boldsymbol{x}_n$ iterate and the predicted one $\tilde{\boldsymbol{x}}_{n+k}$:

$$\boldsymbol{x}_{n+1} \leftarrow \boldsymbol{x}_n + \alpha(\tilde{\boldsymbol{x}}_{n+k} - \boldsymbol{x}_n), \quad \alpha \in [0, 1]\,. \tag{LA}$$

In this work, we use solely GDA as a base optimizer for LA, and denote it with *LAk-GDA*.

Appendix C. Deriving the ACVI algorithm

In this section, we derive an interior-point method for the cVI problem that we refer to as ACVI (**A**DMM-based **i**nterior problem for constrained **VI**s). We first restate the cVI problem in a form that will allow us to derive an interior-point procedure. By the definition of cVI it follows [see §1.3 in 26] that:

$$\boldsymbol{x} \in \mathcal{S}^{\star}_{\mathcal{C},F} \Leftrightarrow \begin{cases} \boldsymbol{w} = \boldsymbol{x} \\ \boldsymbol{x} = \underset{\boldsymbol{z}}{argmin} F(\boldsymbol{w})^{\intercal}\boldsymbol{z} \\ s.t. \quad \varphi(\boldsymbol{z}) \leq \boldsymbol{0} \\ \qquad \boldsymbol{C}\boldsymbol{z} = \boldsymbol{d} \end{cases} \Leftrightarrow \begin{cases} F(\boldsymbol{x}) + \nabla\varphi^{\intercal}(\boldsymbol{x})\boldsymbol{\lambda} + \boldsymbol{C}^{\intercal}\boldsymbol{\nu} = \boldsymbol{0} \\ \boldsymbol{C}\boldsymbol{x} = \boldsymbol{d} \\ \boldsymbol{0} \leq \boldsymbol{\lambda} \perp \varphi(\boldsymbol{x}) \leq \boldsymbol{0}, \end{cases} \tag{KKT}$$

where $\boldsymbol{\lambda} \in \mathbb{R}^m$ and $\boldsymbol{\nu} \in \mathbb{R}^p$ are dual variables. Recall that we assume that *int* $\mathcal{C} \neq \emptyset$, thus, by the Slater condition (using the fact that $\varphi_i(\boldsymbol{x}), i \in [m]$ are convex) and the KKT conditions, the second equivalence holds, yielding the KKT system of cVI. Note that the above equivalence also guarantees the two solutions coincide; see Facchinei and Pang [26, Prop. 1.3.4 (b)]. Analogous to the method described in § 2, we add a log-barrier term to the objective to remove the inequality constraints and obtain the following modified version of (KKT):

$$\begin{cases} \boldsymbol{w} = \boldsymbol{x} \\ \boldsymbol{x} = \underset{\boldsymbol{z}}{argmin}\ F(\boldsymbol{w})^{\intercal}\boldsymbol{z} - \mu \sum_{i=1}^{m} \log\left(-\varphi_i(\boldsymbol{z})\right) \\ s.t. \quad \boldsymbol{C}\boldsymbol{z} = \boldsymbol{d} \end{cases} \Leftrightarrow \begin{cases} F(\boldsymbol{x}) + \nabla\varphi^{\intercal}(\boldsymbol{x})\boldsymbol{\lambda} + \boldsymbol{C}^{\intercal}\boldsymbol{\nu} = \boldsymbol{0} \\ \boldsymbol{\lambda} \odot \varphi(\boldsymbol{x}) + \mu\boldsymbol{e} = \boldsymbol{0} \\ \boldsymbol{C}\boldsymbol{x} - \boldsymbol{d} = \boldsymbol{0} \\ \varphi(\boldsymbol{x}) < \boldsymbol{0}, \boldsymbol{\lambda} > \boldsymbol{0}, \end{cases} \tag{KKT-2}$$

with $\mu > 0$, $\boldsymbol{e} \triangleq [1, \ldots, 1]^{\intercal} \in \mathbb{R}^m$. Again, the equivalence holds by the KKT and the Slater condition. We derive the update rule at step k via the following subproblem: $\min_{\boldsymbol{x}} F(\boldsymbol{w}_k)^{\intercal}\boldsymbol{x} - \mu \sum_{i=1}^{m} \log\left(-\varphi_i(\boldsymbol{x})\right), s.t.\ \boldsymbol{C}\boldsymbol{x} = \boldsymbol{d}$, where we fix $\boldsymbol{w} = \boldsymbol{w}_k$. Notice that (i) $\boldsymbol{w}_k$ is a constant vector in this subproblem, and (ii) the objective is split, making ADMM a natural choice to solve the subproblem. To apply an ADMM-type method, we introduce a new variable $\boldsymbol{y} \in \mathbb{R}^n$ yielding:

$$\begin{cases} \min_{\boldsymbol{x},\boldsymbol{y}} F(\boldsymbol{w}_k)^{\intercal}\boldsymbol{x} + \mathbb{1}[\boldsymbol{C}\boldsymbol{x} = \boldsymbol{d}] - \mu \sum_{i=1}^{m} \log\left(-\varphi_i(\boldsymbol{y})\right) \\ s.t. \qquad \boldsymbol{x} = \boldsymbol{y} \end{cases}, \quad \mathbb{1}[\boldsymbol{C}\boldsymbol{x} = \boldsymbol{d}] \triangleq \begin{cases} 0, & \text{if } \boldsymbol{C}\boldsymbol{x} = \boldsymbol{d} \\ +\infty, & \text{if } \boldsymbol{C}\boldsymbol{x} \neq \boldsymbol{d}. \end{cases} \tag{1}$$

Note that $\mathbb{1}[\boldsymbol{C}\boldsymbol{x} = \boldsymbol{d}]$ is a generalized real-valued convex function of $\boldsymbol{x}$. We introduce the following:

$$\boldsymbol{P}_c \triangleq \boldsymbol{I} - \boldsymbol{C}^{\intercal}(\boldsymbol{C}\boldsymbol{C}^{\intercal})^{-1}\boldsymbol{C}, \qquad (\boldsymbol{P}_c) \qquad \text{and} \qquad \boldsymbol{d}_c \triangleq \boldsymbol{C}^{\intercal}(\boldsymbol{C}\boldsymbol{C}^{\intercal})^{-1}\boldsymbol{d}, \qquad (d_c\text{-EQ})$$

where $\boldsymbol{P}_c \in \mathbb{R}^{n\times n}$ and $\boldsymbol{d}_c \in \mathbb{R}^n$. The augmented Lagrangian of (1) is thus:

$$\mathcal{L}_\beta(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{\lambda}) = F(\boldsymbol{w}_k)^{\intercal}\boldsymbol{x} + \mathbb{1}(\boldsymbol{C}\boldsymbol{x} = \boldsymbol{d}) - \mu \sum_{i=1}^{m} \log(-\varphi_i(\boldsymbol{y})) + \langle \boldsymbol{\lambda}, \boldsymbol{x} - \boldsymbol{y} \rangle + \frac{\beta}{2} \|\boldsymbol{x} - \boldsymbol{y}\|^2, \tag{AL}$$

where $\beta > 0$ is the penalty parameter. Finally, using ADMM, we have the following update rule for $\boldsymbol{x}$ at step k:

$$\boldsymbol{x}_{k+1} = \arg\min_{\boldsymbol{x}\in\mathcal{C}_=} \mathcal{L}_\beta(\boldsymbol{x}, \boldsymbol{y}_k, \boldsymbol{\lambda}_k) = \arg\min_{\boldsymbol{x}\in\mathcal{C}_=} \frac{\beta}{2} \left\| \boldsymbol{x} - \boldsymbol{y}_k + \frac{1}{\beta}(F(\boldsymbol{w}_k) + \boldsymbol{\lambda}_k) \right\|^2. \tag{2}$$

This yields the following update for $\boldsymbol{x}$:

$$\boldsymbol{x}_{k+1} = \boldsymbol{P}_c\Big(\boldsymbol{y}_k - \frac{1}{\beta}\big(F(\boldsymbol{w}_k) + \boldsymbol{\lambda}_k\big)\Big) + \boldsymbol{d}_c\,. \tag{X-EQ}$$

For $\boldsymbol{y}$ and the dual variable $\boldsymbol{\lambda}$, we have:

$$\boldsymbol{y}_{k+1} = \underset{\boldsymbol{y}}{argmin}\ \mathcal{L}_\beta(\boldsymbol{x}_{k+1}, \boldsymbol{y}, \boldsymbol{\lambda}_k) = \underset{\boldsymbol{y}}{argmin}\left(-\mu\sum_{i=1}^{m}\log\big(-\varphi_i(\boldsymbol{y})\big) + \frac{\beta}{2}\left\|\boldsymbol{y} - \boldsymbol{x}_{k+1} - \frac{1}{\beta}\boldsymbol{\lambda}_k\right\|^2\right), \tag{Y-EQ}$$

$$\boldsymbol{\lambda}_{k+1} = \boldsymbol{\lambda}_k + \beta(\boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1}). \tag{λ-EQ}$$

Next, we derive the update rule for $\boldsymbol{w}$. We set $\boldsymbol{w}_k$ to be the solution of the following equation:

$$\boldsymbol{w} + \frac{1}{\beta}\boldsymbol{P}_c F(\boldsymbol{w}) - \boldsymbol{P}_c\boldsymbol{y}_k + \frac{1}{\beta}\boldsymbol{P}_c\boldsymbol{\lambda}_k - \boldsymbol{d}_c = \boldsymbol{0}. \tag{W-EQ}$$

The following theorem ensures the solution of (W-EQ) exists and is unique, see App. D.2 for proof.

Theorem 5 (W-EQ: solution uniqueness) *If F is monotone on $\mathcal{C}_=$, the following statements hold true for the solution of* (W-EQ)*: (i) it always exists, (ii) it is unique, and (iii) it is contained in $\mathcal{C}_=$.*

Remark 1 *Note that when there are no equality constraints, $\mathcal{C}_=$ becomes the entire space $\mathbb{R}^n$. Further notice that $\boldsymbol{w}_k = \boldsymbol{x}_{k+1}$, thus it is redundant to state it in the algorithm, and we remove $\boldsymbol{w}$.*

Appendix D. Omitted Proofs and Discussions concerning Algorithm 1

This section first discusses the theoretical settings–that is, the assumptions of the main theorems, and then provides the proofs of the theoretical results in § 3.

D.1. Discussion on the assumptions of the main theorems

We consider two broad classes of problems. The first class assumes that F is ξ-monotone on $\mathcal{C}_=$—a stronger assumption than monotonicity, yet weaker than strong monotonicity. The second setup requires that (i) F is monotone, (ii) the constraints are active at the solution, and (iii) F is not purely rotational. Note that (iii) is weaker than requiring that the active constraints at the solution form an acute angle with the operator; in other words, given the latter, the former holds due to monotonicity of F. (See App. D). Note that (iii) is not strong, as purely rotational games occur "almost never" in a Baire category sense [4, 43, 47, 74]. The proofs of the main theorems use the following lemma.

Lemma 1 (Upper bound for $\mathcal{G}(\cdot)$) *When F is L-Lipschitz on $\mathcal{C}_=$—as per Assumption 1—we have that any iterate $\boldsymbol{x}_k$ produced by Algorithm 1 satisfies $\mathcal{G}(\boldsymbol{x}_k, \mathcal{C}) \leq M_0 \|\boldsymbol{x}_k - \boldsymbol{x}^\star\|$, where $M_0 > 0$ depends linearly on L, and $\boldsymbol{x}^\star \in \mathcal{S}^\star_{\mathcal{C},F}$.*

Assumption (iii) in Theorem 2 requires the angle of $F(\boldsymbol{x})$ and $\boldsymbol{x} - \boldsymbol{x}^\star$ to be acute on $\mathcal{S} \setminus \{\boldsymbol{x}^\star\}$, where $\mathcal{S} = \hat{\mathcal{C}}_r$ or $\tilde{\mathcal{C}}_s$. For example, when there are no equality constraints, Assumption (iii) becomes $\underset{\boldsymbol{x} \in \mathcal{C} \setminus \{\boldsymbol{x}^\star\}}{inf} F(\boldsymbol{x})^\intercal \frac{\boldsymbol{x} - \boldsymbol{x}^\star}{\|\boldsymbol{x} - \boldsymbol{x}^\star\|} > 0$. From (cVI) and by the monotonicity of F, we can see that for any point $x \in \mathcal{C} \setminus \{\boldsymbol{x}^\star\}$, the angle between $F(\boldsymbol{x})$ and $\boldsymbol{x} - \boldsymbol{x}^\star$ is always less than or equal to $\pi/2$. And assumption (iii) requires that $F(\boldsymbol{x}^\star) \neq \mathbf{0}$, which means some constraints are active at $\boldsymbol{x}^\star$, and $\exists \theta \in (0, \pi/2)$ s.t. for any $\boldsymbol{x} \in \mathcal{C} \setminus \{\boldsymbol{x}^\star\}$, the angle between $F(\boldsymbol{x})$ and $\boldsymbol{x} - \boldsymbol{x}^\star$ is upper bounded by θ.

Remark 2 *Note that the convergence guarantee does not rely on Assumption 1, and it is solely used to relate the rate to the gap function. Also, note that μ_{-1} does not impact the convergence rate. Moreover, for simplicity we state the result with sufficiently small μ_{-1}, however, the proof extends to any $\mu_{-1} > 0$. That is, the above result can be made parameter-free; see App. D.*

Remark 3 *Our proofs rely on the existence of the central path—see Appendix B. Note that since $\mathcal{C}$ is compact, it suffices that either: (i) F is strictly monotone on $\mathcal{C}$, or that (ii) one of the inequality constraints φ_i is strictly convex for the central path to exist [26, Corollary 11.4.24]. Thus, if F is ξ-monotone on $\mathcal{C}$, then the central path exists. However, to relax the former assumption, notice that—by the compactness of $\mathcal{C}$—there exists a sufficiently large M such that for any $\boldsymbol{x} \in \mathcal{C}$, $\boldsymbol{x}^\intercal \boldsymbol{x} \leq M$. Thus, one can add a strictly convex inequality constraint $\varphi_{m+1}(\boldsymbol{x})$—e.g., $\boldsymbol{x}^\intercal \boldsymbol{x} - M \leq 0$—and the solution set remains intact. That is, as μ tends to 0 the original problem is recovered. This ensures the existence of the central path without changing the original problem.*

D.2. Proof of Theorem 5: uniqueness of the solution of Eq. W-EQ

Recall first that Eq. W-EQ is as follows:

$$\boldsymbol{x} + \frac{1}{\beta} \boldsymbol{P}_c F(\boldsymbol{x}) - \boldsymbol{P}_c \boldsymbol{y}_k + \frac{1}{\beta} \boldsymbol{P}_c \boldsymbol{\lambda}_k - \boldsymbol{d}_c = \mathbf{0}\,,$$

since $\boldsymbol{w} = \boldsymbol{x}$.

Proof [Proof of Theorem 5: uniqueness of the solution of (W-EQ)] Let $G(\boldsymbol{x})$ denote the LHS of (W-EQ), that is:

$$G(\boldsymbol{x}) \triangleq \boldsymbol{x} + \frac{1}{\beta}\boldsymbol{P}_c F(\boldsymbol{x}) - \boldsymbol{P}_c \boldsymbol{y}_k + \frac{1}{\beta}\boldsymbol{P}_c \boldsymbol{\lambda}_k - \boldsymbol{d}_c \tag{3}$$

We claim that $G(\boldsymbol{x})$ is strongly monotone on $\mathcal{C}_=$. In fact, $\forall \boldsymbol{x}, \boldsymbol{y} \in \mathcal{C}_=$, $\boldsymbol{P}_c(\boldsymbol{x} - \boldsymbol{y}) = \boldsymbol{x} - \boldsymbol{y}$. Note that $\boldsymbol{P}_c$ is symmetric, thus we have:

$$\begin{aligned}
\langle G(\boldsymbol{x}) - G(\boldsymbol{y}), \boldsymbol{x} - \boldsymbol{y}\rangle &= \|\boldsymbol{x} - \boldsymbol{y}\|^2 + \frac{1}{\beta}\langle \boldsymbol{P}_c F(\boldsymbol{x}) - \boldsymbol{P}_c F(\boldsymbol{y}), \boldsymbol{x} - \boldsymbol{y}\rangle \\
&= \|\boldsymbol{x} - \boldsymbol{y}\|^2 + \frac{1}{\beta}\langle \boldsymbol{x} - \boldsymbol{y}, F(\boldsymbol{x}) - F(\boldsymbol{y})\rangle \\
&\geq \|\boldsymbol{x} - \boldsymbol{y}\|^2.
\end{aligned}$$

Therefore, according to Theorem 2.3.3 (b) in [26], $\mathcal{S}^{\star}_{\mathcal{C}_=, G}$ has a unique solution $\tilde{\boldsymbol{x}} \in \mathcal{C}_=$. Thus, we have:

$$G(\tilde{\boldsymbol{x}})^{\intercal}(\boldsymbol{x} - \tilde{\boldsymbol{x}}) = 0, \ \forall \boldsymbol{x} \in \mathcal{C}_= .$$

From the above, we deduce that $G(\boldsymbol{x}) \in Span\{\boldsymbol{c}_1, \cdots, \boldsymbol{c}_p\}$, where $\boldsymbol{c}_i$ is the row vectors of $\boldsymbol{C}$, $i \in [p]$.

Suppose that $G(\tilde{\boldsymbol{x}}) = \sum_{i=1}^{p} \alpha_i \boldsymbol{c}_i$. Notice that $\boldsymbol{C}G(\boldsymbol{x}) = \boldsymbol{0}$, $\forall \boldsymbol{x} \in \mathcal{C}_=$. Thus, we have that:

$$\boldsymbol{c}_j^{\intercal} G(\tilde{\boldsymbol{x}}) = \boldsymbol{c}_j^{\intercal} \sum_{i=1}^{p} \alpha_i \boldsymbol{c}_i = 0, \ \forall j \in [p] .$$

Hence,

$$\langle G(\tilde{\boldsymbol{x}}), G(\tilde{\boldsymbol{x}})\rangle = \langle \sum_{i=1}^{p} \alpha_i \boldsymbol{c}_i, \sum_{i=1}^{p} \alpha_i \boldsymbol{c}_i\rangle = 0 ,$$

which indicates that $G(\tilde{\boldsymbol{x}}) = \boldsymbol{0}$. Hence, $\tilde{\boldsymbol{x}}$ is a solution of (W-EQ) in $\mathcal{C}_=$.

On the other hand, $\forall x \in \mathbb{R}^n$, if $\boldsymbol{x}$ is a solution of (W-EQ), i.e. $G(\boldsymbol{x}) = \boldsymbol{0}$, then $\boldsymbol{x} \in \mathcal{C}_=$. By the uniqueness of $\tilde{\boldsymbol{x}}$ in $\mathcal{C}_=$ we have that $\boldsymbol{x} = \tilde{\boldsymbol{x}}$, which means $\tilde{\boldsymbol{x}}$ is unique in $\mathbb{R}^n$. ■

D.3. Proof of Lemma 1: Upper bound on the gap function

Proof [Proof of Lemma 1: Upper bound on the gap function] Let $\boldsymbol{x}_k$ denote an iterate produced by Algorithm 1, and let $\boldsymbol{x} \in \mathcal{C}$. Note that we always have $\boldsymbol{x}_k \in \mathcal{C}_=$. We have that:

$$\begin{aligned}
\langle F(\boldsymbol{x}_k), \boldsymbol{x}_k - \boldsymbol{x}\rangle &= \langle F(\boldsymbol{x}^\star), \boldsymbol{x}^\star - \boldsymbol{x}\rangle + \langle F(\boldsymbol{x}_k), \boldsymbol{x}_k - \boldsymbol{x}\rangle - \langle F(\boldsymbol{x}^\star), \boldsymbol{x}^\star - \boldsymbol{x}\rangle \\
&= \langle F(\boldsymbol{x}^\star), \boldsymbol{x}^\star - \boldsymbol{x}\rangle + \langle F(\boldsymbol{x}_k), \boldsymbol{x}_k - \boldsymbol{x}^\star + \boldsymbol{x}^\star - \boldsymbol{x}\rangle - \langle F(\boldsymbol{x}^\star), \boldsymbol{x}^\star - \boldsymbol{x}\rangle .
\end{aligned}$$

From the proof of Theorem 6 we know that $\boldsymbol{x}_k$ is bounded, which gives: $\langle F(\boldsymbol{x}_k), \boldsymbol{x}_k - \boldsymbol{x}^\star\rangle \leq M\|\boldsymbol{x}_k - \boldsymbol{x}^\star\|$, with $M > 0$, as well as that $\langle F(\boldsymbol{x}^\star), \boldsymbol{x}^\star - \boldsymbol{x}\rangle \leq 0$. Thus, for the above we get:

$$\begin{aligned}
\langle F(\boldsymbol{x}_k), \boldsymbol{x}_k - \boldsymbol{x}\rangle &\leq \langle F(\boldsymbol{x}_k) - F(\boldsymbol{x}^\star), \boldsymbol{x}^\star - \boldsymbol{x}\rangle + M\|\boldsymbol{x}_k - \boldsymbol{x}^\star\| \\
&\leq D\|F(\boldsymbol{x}_k) - F(\boldsymbol{x}^\star)\| + M\|\boldsymbol{x}_k - \boldsymbol{x}^\star\| \\
&\leq (DL + M)\|\boldsymbol{x}_k - \boldsymbol{x}^\star\| ,
\end{aligned}$$

where for the second row we used $\boldsymbol{x}^\star - \boldsymbol{x} \leq D$ where $D \triangleq \max_{\boldsymbol{x}' \in \mathcal{C}} \|\boldsymbol{x}^\star - \boldsymbol{x}'\|$ is the largest distance between any point in $\mathcal{C}$ and $\boldsymbol{x}^\star$. For the last row we used that F is L-Lipschitz—Assumption 1—which concludes the proof. ■

The proof is analogous for the $\boldsymbol{y}_k \in \mathcal{C}_\leq$ iterates produced by Algorithm 1.

D.4. Proofs of the convergence rate: Theorems 1 and 2

Let

$$f^{(t)}(\boldsymbol{x}) \triangleq F(\boldsymbol{x}^{\mu_t})^\intercal \boldsymbol{x} + \mathbb{1}(\boldsymbol{C}\boldsymbol{x} = \boldsymbol{d}),$$
$$f_k^{(t)}(\boldsymbol{x}) \triangleq F(\boldsymbol{x}_k^{(t)})^\intercal \boldsymbol{x} + \mathbb{1}(\boldsymbol{C}\boldsymbol{x} = \boldsymbol{d}),$$
$$g^{(t)}(\boldsymbol{y}) \triangleq -\mu_t \sum_{i=1}^{m} \log\big(-\varphi_i(\boldsymbol{y})\big),$$

where $\boldsymbol{x}^{\mu_t}$ is a solution of (KKT-2) when $\mu = \mu_t$. Note that the existence of $\boldsymbol{x}^{\mu_t}$ is guaranteed by the existence of the central path-see App. B, and that $f^{(t)}$, $f_k^{(t)}$ and $g^{(t)}$ are all convex.

In the following proofs, unless causing confusion, we drop the subscript t to simplify notations.

Let $\boldsymbol{y}^\mu = \boldsymbol{x}^\mu$, from (KKT-2) we can see that $(\boldsymbol{x}^\mu, \boldsymbol{y}^\mu)$ is an optimal solution of

$$\begin{cases} \min f(\boldsymbol{x}) + g(\boldsymbol{y}) \\ \quad s.t. \quad \boldsymbol{x} = \boldsymbol{y} \end{cases}. \tag{4}$$

There exists $\boldsymbol{\lambda}^\mu \in \mathbb{R}^n$ such that $(\boldsymbol{x}^\mu, \boldsymbol{y}^\mu, \boldsymbol{\lambda}^\mu)$ is a KKT point of (4). We give the following proposition which we will repeatedly use in the proofs:

Proposition 2 *If F is monotone, then $\forall k \in \mathbb{N}$,*

$$f_k(\boldsymbol{x}_{k+1}) - f_k(\boldsymbol{x}^\mu) \geq f(\boldsymbol{x}_{k+1}) - f(\boldsymbol{x}^\mu).$$

Furthermore, if F is ξ-monotone, as per Def. 1

$$f_k(\boldsymbol{x}_{k+1}) - f_k(\boldsymbol{x}^\mu) \geq f(\boldsymbol{x}_{k+1}) - f(\boldsymbol{x}^\mu) + c\|\boldsymbol{x}_{k+1} - \boldsymbol{x}^\mu\|_2^\xi.$$

Proof [Proof of Proposition 2] It suffices to note that:

$$f_k(\boldsymbol{x}_{k+1}) - f_k(\boldsymbol{x}^\mu) - (f(\boldsymbol{x}_{k+1}) - f(\boldsymbol{x}^\mu)) = (F(\boldsymbol{x}_{k+1}) - F(\boldsymbol{x}^\mu))^\intercal (\boldsymbol{x}_{k+1} - \boldsymbol{x}^\mu).$$

■

Some of our proofs that follow are inspired by some convergence proofs in ADMM [21, 25, 29, 40, 41, 51]. However, although Algorithm 1 adopts the high level idea of ADMM, we can not directly refer to the convergence proofs of ADMM, but need to substantially modify these.

We will use the following lemma.

Lemma 2 $f(\boldsymbol{x}) + g(\boldsymbol{y}) - f(\boldsymbol{x}^\mu) - g(\boldsymbol{y}^\mu) + \langle \boldsymbol{\lambda}^\mu, \boldsymbol{x} - \boldsymbol{y} \rangle \geq 0, \forall \boldsymbol{x}, \boldsymbol{y}.$

Proof The Lagrange function of (4) is

$$L(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{\lambda}) = f(\boldsymbol{x}) + g(\boldsymbol{y}) + \boldsymbol{\lambda}^{\mathsf{T}}(\boldsymbol{x} - \boldsymbol{y}).$$

And by the property of KKT point we have

$$L(\boldsymbol{x}^{\mu}, \boldsymbol{y}^{\mu}, \boldsymbol{\lambda}) \leq L(\boldsymbol{x}^{\mu}, \boldsymbol{y}^{\mu}, \boldsymbol{\lambda}^{\mu}) \leq L(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{\lambda}^{\mu}),\ \forall(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{\lambda}),$$

from which the conclusion follows. ■

The following lemma is straightforward to verify:

Lemma 3 *If*

$$\begin{aligned} f(\boldsymbol{x}) + g(\boldsymbol{y}) - f(\boldsymbol{x}^{\mu}) - g(\boldsymbol{y}^{\mu}) + \langle \boldsymbol{\lambda}^{\mu}, \boldsymbol{x} - \boldsymbol{y} \rangle &\leq \alpha_1, \\ \|\boldsymbol{x} - \boldsymbol{y}\| &\leq \alpha_2 \end{aligned}$$

then we have

$$-\|\boldsymbol{\lambda}^{\mu}\|\alpha_2 \leq f(\boldsymbol{x}) + g(\boldsymbol{y}) - f(\boldsymbol{x}^{\mu}) - g(\boldsymbol{y}^{\mu}) \leq \|\boldsymbol{\lambda}^{\mu}\|\alpha_2 + \alpha_1.$$

The following lemma lists some simple but useful facts that we will use in the following proofs.

Lemma 4 *For* (4) *and Algorithm 1, we have*

$$\mathbf{0} \in \partial f_k(\boldsymbol{x}_{k+1}) + \boldsymbol{\lambda}_k + \beta(\boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1}) \tag{5}$$

$$\mathbf{0} \in \partial g(\boldsymbol{y}_{k+1}) - \boldsymbol{\lambda}_k - \beta(\boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1}), \tag{6}$$

$$\boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}_k = \beta(\boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1}), \tag{7}$$

$$\mathbf{0} \in \partial f(\boldsymbol{x}^{\mu}) + \boldsymbol{\lambda}^{\mu}, \tag{8}$$

$$\mathbf{0} \in \partial g(\boldsymbol{y}^{\mu}) - \boldsymbol{\lambda}^{\mu}, \tag{9}$$

$$\boldsymbol{x}^{\mu} = \boldsymbol{y}^{\mu}. \tag{10}$$

We define:

$$\hat{\nabla} f_k(\boldsymbol{x}_{k+1}) \triangleq -\boldsymbol{\lambda}_k - \beta(\boldsymbol{x}_{k+1} - \boldsymbol{y}_k), \tag{11}$$

$$\hat{\nabla} g(\boldsymbol{y}_{k+1}) \triangleq \boldsymbol{\lambda}_k + \beta(\boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1}). \tag{12}$$

Then from (5) and (6) we can see that

$$\hat{\nabla} f_k(\boldsymbol{x}_{k+1}) \in \partial f_k(\boldsymbol{x}_{k+1}) \text{ and } \hat{\nabla} g(\boldsymbol{y}_{k+1}) \in \partial g(\boldsymbol{y}_{k+1}). \tag{13}$$

Lemma 5 *For Algorithm 1, we have*

$$\langle \hat{\nabla} g(\boldsymbol{y}_{k+1}), \boldsymbol{y}_{k+1} - \boldsymbol{y} \rangle = -\langle \boldsymbol{\lambda}_{k+1}, \boldsymbol{y} - \boldsymbol{y}_{k+1} \rangle, \tag{14}$$

and

$$\begin{aligned} &\langle \hat{\nabla} f_k(\boldsymbol{x}_{k+1}), \boldsymbol{x}_{k+1} - \boldsymbol{x} \rangle + \langle \hat{\nabla} g(\boldsymbol{y}_{k+1}), \boldsymbol{y}_{k+1} - \boldsymbol{y} \rangle \\ =& - \langle \boldsymbol{\lambda}_{k+1}, \boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1} - \boldsymbol{x} + \boldsymbol{y} \rangle + \beta \langle -\boldsymbol{y}_{k+1} + \boldsymbol{y}_k, \boldsymbol{x}_{k+1} - \boldsymbol{x} \rangle. \end{aligned} \tag{15}$$

Proof [Proof of Lemma 5] From (7), (11) and (12) we have:

$$\begin{aligned}&\langle \hat{\nabla} f_k(\boldsymbol{x}_{k+1}), \boldsymbol{x}_{k+1} - \boldsymbol{x}\rangle \\ =& -\langle \boldsymbol{\lambda}_k + \beta(\boldsymbol{x}_{k+1} - \boldsymbol{y}_k), \boldsymbol{x}_{k+1} - \boldsymbol{x}\rangle \\ =& -\langle \boldsymbol{\lambda}_{k+1}, \boldsymbol{x}_{k+1} - \boldsymbol{x}\rangle + \beta\langle -\boldsymbol{y}_{k+1} + \boldsymbol{y}_k, \boldsymbol{x}_{k+1} - \boldsymbol{x}\rangle,\end{aligned}$$

and

$$\langle \hat{\nabla} g(\boldsymbol{y}_{k+1}), \boldsymbol{y}_{k+1} - \boldsymbol{y}\rangle = -\langle \boldsymbol{\lambda}_{k+1}, \boldsymbol{y} - \boldsymbol{y}_{k+1}\rangle.$$

Adding these together yields (15). ■

Lemma 6 *For Algorithm 1, we have*

$$\begin{aligned}&\langle \hat{\nabla} f_k(\boldsymbol{x}_{k+1}), \boldsymbol{x}_{k+1} - \boldsymbol{x}^\mu\rangle + \langle \hat{\nabla} g(\boldsymbol{y}_{k+1}), \boldsymbol{y}_{k+1} - \boldsymbol{y}^\mu\rangle + \langle \boldsymbol{\lambda}^\mu, \boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1}\rangle \\ \le& \frac{1}{2\beta}\|\boldsymbol{\lambda}_k - \boldsymbol{\lambda}^\mu\|^2 - \frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}^\mu\|^2 \\ &+ \frac{\beta}{2}\|\boldsymbol{y}^\mu - \boldsymbol{y}_k\|^2 - \frac{\beta}{2}\|\boldsymbol{y}^\mu - \boldsymbol{y}_{k+1}\|^2 \\ &- \frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}_k\|^2 - \frac{\beta}{2}\|\boldsymbol{y}_k - \boldsymbol{y}_{k+1}\|^2\end{aligned}$$

Proof [Proof of Lemma 6] Letting $(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{\lambda}) = (\boldsymbol{x}^\mu, \boldsymbol{y}^\mu, \boldsymbol{\lambda}^\mu)$ in(15), adding $\langle \boldsymbol{\lambda}^\mu, \boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1}\rangle$ to both sides, and using (7) and (10), we have:

$$\begin{aligned}&\langle \hat{\nabla} f_k(\boldsymbol{x}_{k+1}), \boldsymbol{x}_{k+1} - \boldsymbol{x}^\mu\rangle + \langle \hat{\nabla} g(\boldsymbol{y}_{k+1}), \boldsymbol{y}_{k+1} - \boldsymbol{y}^\mu\rangle + \langle \boldsymbol{\lambda}^\mu, \boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1}\rangle \\ =& -\langle \boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}^\mu, \boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1}\rangle + \beta\langle -\boldsymbol{y}_{k+1} + \boldsymbol{y}_k, \boldsymbol{x}_{k+1} - \boldsymbol{x}^\mu\rangle \\ =& -\frac{1}{\beta}\langle \boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}^\mu, \boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}_k\rangle + \langle -\boldsymbol{y}_{k+1} + \boldsymbol{y}_k, \boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}_k\rangle \\ &- \beta\langle -\boldsymbol{y}_{k+1} + \boldsymbol{y}_k, -\boldsymbol{y}_{k+1} + \boldsymbol{y}^\mu\rangle \qquad (16) \\ =& \frac{1}{2\beta}\|\boldsymbol{\lambda}_k - \boldsymbol{\lambda}^\mu\|^2 - \frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}^\mu\|^2 - \frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}_k\|^2 \\ &+ \frac{\beta}{2}\|-\boldsymbol{y}_k + \boldsymbol{y}^\star\|^2 - \frac{\beta}{2}\|-\boldsymbol{y}_{k+1} + \boldsymbol{y}^\star\|^2 - \frac{\beta}{2}\|-\boldsymbol{y}_{k+1} + \boldsymbol{y}_k\|^2 \\ &+ \langle -\boldsymbol{y}_{k+1} + \boldsymbol{y}_k, \boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}_k\rangle\,. \qquad (17)\end{aligned}$$

On the other hand, (14) gives

$$\langle \hat{\nabla} g(\boldsymbol{y}_k), \boldsymbol{y}_k - \boldsymbol{y}\rangle + \langle \boldsymbol{\lambda}_k, -\boldsymbol{y}_k + \boldsymbol{y}\rangle = 0\,. \qquad (18)$$

Letting $\boldsymbol{y} = \boldsymbol{y}_k$ in (14) and $\boldsymbol{y} = \boldsymbol{y}_{k+1}$ in (18), and adding them together, we have:

$$\langle \hat{\nabla} g(\boldsymbol{y}_{k+1}) - \hat{\nabla} g(\boldsymbol{y}_k), \boldsymbol{y}_{k+1} - \boldsymbol{y}_k\rangle + \langle \boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}_k, -\boldsymbol{y}_{k+1} + \boldsymbol{y}_k\rangle = 0\,.$$

By the monotonicity of ∂g we know that the first term of the above equality is non-negative. Thus, we have:

$$\langle \boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}_k, -\boldsymbol{y}_{k+1} + \boldsymbol{y}_k\rangle \le 0\,. \qquad (19)$$

Plugging it into (17), we have the conclusion. ■

Lemma 7 *For Algorithm 1, we have*

$$\begin{aligned}
&f(\boldsymbol{x}_{k+1}) + g(\boldsymbol{y}_{k+1}) - f(\boldsymbol{x}^\mu) - g(\boldsymbol{y}^\mu) + \langle \boldsymbol{\lambda}^\mu, \boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1} \rangle \\
\leq &\frac{1}{2\beta}\|\boldsymbol{\lambda}_k - \boldsymbol{\lambda}^\mu\|^2 - \frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}^\mu\|^2 \\
+ &\frac{\beta}{2}\|-\boldsymbol{y}_k + \boldsymbol{y}^\mu\|^2 - \frac{\beta}{2}\|-\boldsymbol{y}_{k+1} + \boldsymbol{y}^\mu\|^2 \\
- &\frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}_k\|^2 - \frac{\beta}{2}\|-\boldsymbol{y}_{k+1} + \boldsymbol{y}_k\|^2
\end{aligned} \tag{20}$$

Furthermore, if F is ξ-monotone on $\mathcal{C}_=$, we have

$$\begin{aligned}
&c\|\boldsymbol{x}_{k+1} - \boldsymbol{x}^\mu\|_2^\xi + f(\boldsymbol{x}_{k+1}) + g(\boldsymbol{y}_{k+1}) - f(\boldsymbol{x}^\mu) - g(\boldsymbol{y}^\mu) + \langle \boldsymbol{\lambda}^\mu, \boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1} \rangle \\
\leq &\frac{1}{2\beta}\|\boldsymbol{\lambda}_k - \boldsymbol{\lambda}^\mu\|^2 - \frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}^\mu\|^2 \\
+ &\frac{\beta}{2}\|-\boldsymbol{y}_k + \boldsymbol{y}^\mu\|^2 - \frac{\beta}{2}\|-\boldsymbol{y}_{k+1} + \boldsymbol{y}^\mu\|^2 \\
- &\frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}_k\|^2 - \frac{\beta}{2}\|-\boldsymbol{y}_{k+1} + \boldsymbol{y}_k\|^2.
\end{aligned} \tag{21}$$

Proof From the convexity of $f_k(\boldsymbol{x})$ and $g(\boldsymbol{y})$ and using Proposition 2 and (13), we have

$$\begin{aligned}
&f(\boldsymbol{x}_{k+1}) + g(\boldsymbol{y}_{k+1}) - f(\boldsymbol{x}^\mu) - g(\boldsymbol{y}^\mu) + \langle \boldsymbol{\lambda}^\mu, \boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1} \rangle \\
\leq &f_k(\boldsymbol{x}_{k+1}) + g(\boldsymbol{y}_{k+1}) - f_k(\boldsymbol{x}^\mu) - g(\boldsymbol{y}^\mu) + \langle \boldsymbol{\lambda}^\mu, \boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1} \rangle \\
\leq &\langle \hat{\nabla} f_k(\boldsymbol{x}_{k+1}), \boldsymbol{x}_{k+1} - \boldsymbol{x}^\mu \rangle + \langle \hat{\nabla} g(\boldsymbol{y}_{k+1}), \boldsymbol{y}_{k+1} - \boldsymbol{y}^\mu \rangle \\
+ &\langle \boldsymbol{\lambda}^\mu, \boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1} \rangle \\
\leq &\frac{1}{2\beta}\|\boldsymbol{\lambda}_k - \boldsymbol{\lambda}^\mu\|^2 - \frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}^\mu\|^2 \\
+ &\frac{\beta}{2}\|-\boldsymbol{y}_k + \boldsymbol{y}^\mu\|^2 - \frac{\beta}{2}\|-\boldsymbol{y}_{k+1} + \boldsymbol{y}^\mu\|^2 \\
- &\frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}_k\|^2 - \frac{\beta}{2}\|-\boldsymbol{y}_{k+1} + \boldsymbol{y}_k\|^2.
\end{aligned} \tag{22}$$

If F is ξ-monotone on $\mathcal{C}_=$, again by Proposition 2, we can add the term $c\|\boldsymbol{x}_{k+1} - \boldsymbol{x}^\mu\|_2^\xi$ to the first line and the inequality still holds. ■

Theorem 6 *For Algorithm 1, we have*

$$\begin{aligned}
f(\boldsymbol{x}_{k+1}) - f(\boldsymbol{x}^\mu) + g(\boldsymbol{y}_{k+1}) - g(\boldsymbol{y}^\mu) &\to 0, \\
f_k(\boldsymbol{x}_{k+1}) - f_k(\boldsymbol{x}^\mu) + g(\boldsymbol{y}_{k+1}) - g(\boldsymbol{y}^\mu) &\to 0, \\
\boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1} &\to \boldsymbol{0},
\end{aligned}$$

as $k \to \infty$. Furthermore, if F is ξ-monotone on $\mathcal{C}_=$, we have

$$\boldsymbol{x}_{k+1} \to \boldsymbol{x}^\mu, \ k \to \infty$$

Proof [Proof of Theorem 6] Proof From Lemma 2 and (20), we have

$$\begin{aligned}&\frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k\|^2+\frac{\beta}{2}\|-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k\|^2\\ \leq&\frac{1}{2\beta}\|\boldsymbol{\lambda}_k-\boldsymbol{\lambda}^\mu\|^2-\frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}^\mu\|^2\\ +&\frac{\beta}{2}\|-\boldsymbol{y}_k+\boldsymbol{y}^\mu\|^2-\frac{\beta}{2}\|-\boldsymbol{y}_{k+1}+\boldsymbol{y}^\mu\|^2.\end{aligned} \tag{23}$$

Summing over $k=0,\cdots,\infty$, we have

$$\begin{aligned}&\sum_{k=0}^{\infty}(\frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k\|^2+\frac{\beta}{2}\|-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k\|^2)\\ \leq&\frac{1}{2\beta}\|\boldsymbol{\lambda}_0-\boldsymbol{\lambda}^\mu\|^2+\frac{\beta}{2}\|-\boldsymbol{y}_0+\boldsymbol{y}^\mu\|^2.\end{aligned}$$

from which we deduce that $\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k\to\mathbf{0}$ and $-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k\to\mathbf{0}$. Moreover, $\|\boldsymbol{\lambda}_k-\boldsymbol{\lambda}^\mu\|^2$ and $\|-\boldsymbol{y}_k+\boldsymbol{y}^\mu\|^2$ are bounded for all k, as well as $\|\boldsymbol{\lambda}_k\|$. Since

$$\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k=\beta(\boldsymbol{x}_{k+1}-\boldsymbol{y}_{k+1})=\beta(\boldsymbol{x}_{k+1}-\boldsymbol{x}^\mu)+\beta(-\boldsymbol{y}_{k+1}+\boldsymbol{y}^\mu)$$

we deduce that $\boldsymbol{x}_{k+1}-\boldsymbol{y}_{k+1}\to\mathbf{0}$ and $\boldsymbol{x}_{k+1}-\boldsymbol{x}^\mu$ is also bounded.

From (15) and the convexity of f and g, and using Proposition 2, we have:

$$\begin{aligned}&f(\boldsymbol{x}_{k+1})-f(\boldsymbol{x}^\mu)+g(\boldsymbol{y}_{k+1})-g(\boldsymbol{y}^\mu)\\ \leq&f_k(\boldsymbol{x}_{k+1})-f_k(\boldsymbol{x}^\mu)+g(\boldsymbol{y}_{k+1})-g(\boldsymbol{y}^\mu)\\ \leq&-\langle\boldsymbol{\lambda}_{k+1},\boldsymbol{x}_{k+1}-\boldsymbol{y}_{k+1}\rangle+\beta\langle-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k,\boldsymbol{x}_{k+1}-\boldsymbol{x}^\mu\rangle\to 0.\end{aligned}$$

On the other hand, from (8), (9), and (10), we have:

$$\begin{aligned}&f_k\left(\boldsymbol{x}_{k+1}\right)-f_k\left(\boldsymbol{x}^\mu\right)+g\left(\boldsymbol{y}_{k+1}\right)-g\left(\boldsymbol{y}^\mu\right)\\ \geq&f\left(\boldsymbol{x}_{k+1}\right)-f\left(\boldsymbol{x}^\mu\right)+g\left(\boldsymbol{y}_{k+1}\right)-g\left(\boldsymbol{y}^\mu\right)\\ \geq&\langle-\boldsymbol{\lambda}^\mu,\boldsymbol{x}_{k+1}-\boldsymbol{x}^\mu\rangle+\langle\boldsymbol{\lambda}^\mu,\boldsymbol{y}_{k+1}-\boldsymbol{y}^\mu\rangle\\ =&-\langle\boldsymbol{\lambda}^\mu,\boldsymbol{x}_{k+1}-\boldsymbol{y}_{k+1}\rangle\to 0.\end{aligned}$$

Thus, we have $f\left(\boldsymbol{x}_{k+1}\right)-f\left(\boldsymbol{x}^\mu\right)+g\left(\boldsymbol{y}_{k+1}\right)-g\left(\boldsymbol{y}^\mu\right)\to 0$ and $f_k\left(\boldsymbol{x}_{k+1}\right)-f_k\left(\boldsymbol{x}^\mu\right)+g\left(\boldsymbol{y}_{k+1}\right)-g\left(\boldsymbol{y}^\mu\right)\to 0, k\to\infty$.

If F is ξ-monotone on $\mathcal{C}_=$, from Lemma 2 and (21) we have

$$\begin{aligned}&c\|\boldsymbol{x}_{k+1}-\boldsymbol{x}^\mu\|_2^\xi+\frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k\|^2+\frac{\beta}{2}\|-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k\|^2\\ \leq&\frac{1}{2\beta}\|\boldsymbol{\lambda}_k-\boldsymbol{\lambda}^\mu\|^2-\frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}^\mu\|^2\\ +&\frac{\beta}{2}\|-\boldsymbol{y}_k+\boldsymbol{y}^\mu\|^2-\frac{\beta}{2}\|-\boldsymbol{y}_{k+1}+\boldsymbol{y}^\mu\|^2\end{aligned} \tag{24}$$

From which we deduce that:

$$
\begin{aligned}
&c\|\boldsymbol{x}_{k+1}-\boldsymbol{x}^{\mu}\|_2^{\xi}+\sum_{k=0}^{\infty}\left(\frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k\|^2+\frac{\beta}{2}\|-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k\|^2\right)\\
\leq&\frac{1}{2\beta}\|\boldsymbol{\lambda}^0-\boldsymbol{\lambda}^{\mu}\|^2+\frac{\beta}{2}\|-\boldsymbol{y}^0+\boldsymbol{y}^{\mu}\|^2.
\end{aligned}
$$

Therefore, $\|\boldsymbol{x}_{k+1}-\boldsymbol{x}^{\mu}\|_2\to 0, k\to\infty$. ■

Lemma 8 *For Algorithm 1, we have*

$$
\begin{aligned}
&\frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k\|^2+\frac{\beta}{2}\|-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k\|^2\\
\leq&\frac{1}{2\beta}\|\boldsymbol{\lambda}_k-\boldsymbol{\lambda}_{k-1}\|^2+\frac{\beta}{2}\|-\boldsymbol{y}_k+\boldsymbol{y}_{k-1}\|^2.
\end{aligned}
\tag{25}
$$

Furthermore, if F is ξ-monotone on $\mathcal{C}_=$, we have

$$
\begin{aligned}
&c\|\boldsymbol{x}_{k+1}-\boldsymbol{x}_k\|^2+\frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k\|^2+\frac{\beta}{2}\|-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k\|^2\\
\leq&\frac{1}{2\beta}\|\boldsymbol{\lambda}_k-\boldsymbol{\lambda}_{k-1}\|^2+\frac{\beta}{2}\|-\boldsymbol{y}_k+\boldsymbol{y}_{k-1}\|^2.
\end{aligned}
\tag{26}
$$

Proof [Proof of Lemma 8] (15) gives:

$$
\begin{aligned}
&\langle\hat{\nabla}f_{k-1}\left(\boldsymbol{x}_k\right),\boldsymbol{x}_k-\boldsymbol{x}\rangle+\langle\hat{\nabla}g\left(\boldsymbol{y}_k\right),\boldsymbol{y}_k-\boldsymbol{y}\rangle\\
=&-\langle\boldsymbol{\lambda}_k,\boldsymbol{x}_k-\boldsymbol{y}_k-\boldsymbol{x}+\boldsymbol{y}\rangle+\beta\langle-\boldsymbol{y}_k+\boldsymbol{y}_{k-1},\boldsymbol{x}_k-\boldsymbol{x}\rangle.
\end{aligned}
\tag{27}
$$

Letting $(\boldsymbol{x},\boldsymbol{y},\boldsymbol{\lambda})=(\boldsymbol{x}_k,\boldsymbol{y}_k,\boldsymbol{\lambda}_k)$ in (15) and $(\boldsymbol{x},\boldsymbol{y},\boldsymbol{\lambda})=(\boldsymbol{x}_{k+1},\boldsymbol{y}_{k+1},\boldsymbol{\lambda}_{k+1})$ in (27), and adding them together, and using (7), we have

$$
\begin{aligned}
&\langle\hat{\nabla}f_k\left(\boldsymbol{x}_{k+1}\right)-\hat{\nabla}f_{k-1}\left(\boldsymbol{x}_k\right),\boldsymbol{x}_{k+1}-\boldsymbol{x}_k\rangle+\langle\hat{\nabla}g\left(\boldsymbol{y}_{k+1}\right)-\hat{\nabla}g\left(\boldsymbol{y}_k\right),\boldsymbol{y}_{k+1}-\boldsymbol{y}_k\rangle\\
=&-\langle\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k,\boldsymbol{x}_{k+1}-\boldsymbol{y}_{k+1}-\boldsymbol{x}_k+\boldsymbol{y}_k\rangle+\beta\langle-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k-\left(-\boldsymbol{y}_k+\boldsymbol{y}_{k-1}\right),\boldsymbol{x}_{k+1}-\boldsymbol{x}_k\rangle\\
=&-\frac{1}{\beta}\langle\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k,\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k-\left(\boldsymbol{\lambda}_k-\boldsymbol{\lambda}_{k-1}\right)\rangle\\
&+\langle-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k+\left(\boldsymbol{y}_k-\boldsymbol{y}_{k-1}\right),\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k+\beta\boldsymbol{y}_{k+1}-\left(\boldsymbol{\lambda}_k-\boldsymbol{\lambda}_{k-1}+\beta\boldsymbol{y}_k\right)\rangle\\
=&\frac{1}{2\beta}\left[\|\boldsymbol{\lambda}_k-\boldsymbol{\lambda}_{k-1}\|^2-\|\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k\|^2-\|\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k-\left(\boldsymbol{\lambda}_k-\boldsymbol{\lambda}_{k-1}\right)\|^2\right]\\
&+\frac{\beta}{2}\left[\|-\boldsymbol{y}_k+\boldsymbol{y}_{k-1}\|^2-\|-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k\|^2-\|-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k-\left(-\boldsymbol{y}_k+\boldsymbol{y}_{k-1}\right)\|^2\right]\\
&+\langle-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k-\left(-\boldsymbol{y}_k+\boldsymbol{y}_{k-1}\right),\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k-\left(\boldsymbol{\lambda}_k-\boldsymbol{\lambda}_{k-1}\right)\rangle\\
=&\frac{1}{2\beta}\left(\|\boldsymbol{\lambda}_k-\boldsymbol{\lambda}_{k-1}\|^2-\|\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k\|^2\right)+\frac{\beta}{2}\left(\|-\boldsymbol{y}_k+\boldsymbol{y}_{k-1}\|^2-\|-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k\|^2\right)\\
&-\frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k-\left(\boldsymbol{\lambda}_k-\boldsymbol{\lambda}_{k-1}\right)\|^2-\frac{\beta}{2}\|-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k-\left(-\boldsymbol{y}_k+\boldsymbol{y}_{k-1}\right)\|^2\\
&+\langle-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k-\left(-\boldsymbol{y}_k+\boldsymbol{y}_{k-1}\right),\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k-\left(\boldsymbol{\lambda}_k-\boldsymbol{\lambda}_{k-1}\right)\rangle\\
\leq&\frac{1}{2\beta}\left(\|\boldsymbol{\lambda}_k-\boldsymbol{\lambda}_{k-1}\|^2-\|\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k\|^2\right)+\frac{\beta}{2}\left(\|-\boldsymbol{y}_k+\boldsymbol{y}_{k-1}\|^2-\|-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k\|^2\right)
\end{aligned}
$$

Note that

$$\begin{aligned}
&\langle \hat{\nabla} f_k\left(\boldsymbol{x}_{k+1}\right)-\hat{\nabla} f_{k-1}\left(\boldsymbol{x}_k\right), \boldsymbol{x}_{k+1}-\boldsymbol{x}_k\rangle \\
=&\langle \hat{\nabla} f_k\left(\boldsymbol{x}_{k+1}\right)-\hat{\nabla} f_k\left(\boldsymbol{x}_k\right), \boldsymbol{x}_{k+1}-\boldsymbol{x}_k\rangle+\langle \hat{\nabla} f_k\left(\boldsymbol{x}_k\right)-\hat{\nabla} f_{k-1}\left(\boldsymbol{x}_k\right), \boldsymbol{x}_{k+1}-\boldsymbol{x}_k\rangle \\
=&\langle \hat{\nabla} f_k\left(\boldsymbol{x}_{k+1}\right)-\hat{\nabla} f_k\left(\boldsymbol{x}_k\right), \boldsymbol{x}_{k+1}-\boldsymbol{x}_k\rangle+\langle F\left(\boldsymbol{x}_{k+1}\right)-F\left(\boldsymbol{x}_k\right), \boldsymbol{x}_{k+1}-\boldsymbol{x}_k\rangle
\end{aligned}$$

Using the monotonicity of $\partial f_k, \partial g$ and F, we draw the conclusion. ■

Theorem 7 *If F is monotone on $\mathcal{C}_=$, then for Algorithm 1, we have*

$$\begin{aligned}
-\|\boldsymbol{\lambda}^\mu\|\sqrt{\frac{\Delta}{\beta(K+1)}} \leq & f\left(\boldsymbol{x}_{K+1}\right)+g\left(\boldsymbol{y}_{K+1}\right)-f\left(\boldsymbol{x}^\mu\right)-g\left(\boldsymbol{y}^\mu\right) \\
\leq & f_K\left(\boldsymbol{x}_{K+1}\right)+g\left(\boldsymbol{y}_{K+1}\right)-f_K\left(\boldsymbol{x}^\mu\right)-g\left(\boldsymbol{y}^\mu\right) \\
\leq & \frac{\Delta}{K+1}+\frac{2 \Delta}{\sqrt{K+1}}+\|\boldsymbol{\lambda}^\mu\|\sqrt{\frac{\Delta}{\beta(K+1)}},
\end{aligned} \tag{28}$$

$$\|\boldsymbol{x}_{K+1}-\boldsymbol{y}_{K+1}\| \leq \sqrt{\frac{\Delta}{\beta(K+1)}}, \tag{29}$$

where $\Delta \triangleq \frac{1}{\beta}\|\boldsymbol{\lambda}_0-\boldsymbol{\lambda}^\mu\|^2+\beta\|\boldsymbol{y}_0-\boldsymbol{y}^\mu\|^2$.

Furthermore, if F is ξ-monotone on $\mathcal{C}_=$, we have

$$c\|\hat{\boldsymbol{x}}_{K+1}-\boldsymbol{x}^\mu\|^\xi \leq \frac{\Delta}{K+1}, \tag{30}$$

$$c\|\boldsymbol{x}_{K+1}-\boldsymbol{x}^\mu\|_2^\xi \leq \frac{\Delta}{K+1}+\frac{2 \Delta}{\sqrt{K+1}}. \tag{31}$$

where $\hat{\boldsymbol{x}}_{K+1}=\frac{1}{K+1} \sum_{k=1}^{K+1} \boldsymbol{x}_K$.

Proof [Proof of Theorem 7.] Summing (23) over $k=0,1,\ldots,K$ and using the monotonicity of $\frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}_k\|^2+\frac{\beta}{2}\|-\boldsymbol{y}_{k+1}+\boldsymbol{y}_k\|^2$ from Lemma 8, we have:

$$\frac{1}{\beta}\|\boldsymbol{\lambda}_{K+1}-\boldsymbol{\lambda}_K\|^2+\beta\|-\boldsymbol{y}_{K+1}+\boldsymbol{y}_K\|^2 \leq \frac{1}{K+1}\left(\frac{1}{\beta}\|\boldsymbol{\lambda}_0-\boldsymbol{\lambda}^\mu\|^2+\beta\|-\boldsymbol{y}_0+\boldsymbol{y}^\mu\|^2\right) \tag{32}$$

From which we deduce that

$$\beta\|\boldsymbol{x}_{K+1}-\boldsymbol{y}_{K+1}\|=\|\boldsymbol{\lambda}_{K+1}-\boldsymbol{\lambda}_K\| \leq \sqrt{\frac{\beta \Delta}{(K+1)}},$$

$$\|-\boldsymbol{y}_{K+1}+\boldsymbol{y}_K\| \leq \sqrt{\frac{\Delta}{\beta(K+1)}}.$$

On the other hand, (23) gives:

$$\begin{aligned}&\frac{1}{2\beta}\|\boldsymbol{\lambda}_{k+1}-\boldsymbol{\lambda}^{\mu}\|^2+\frac{\beta}{2}\|-\boldsymbol{y}_{k+1}+\boldsymbol{y}^{\mu}\|^2\\ \leq&\frac{1}{2\beta}\|\boldsymbol{\lambda}_{k}-\boldsymbol{\lambda}^{\mu}\|^2+\frac{\beta}{2}\|-\boldsymbol{y}_{k}+\boldsymbol{y}^{\mu}\|^2\\ \leq&\frac{1}{2\beta}\|\boldsymbol{\lambda}_{0}-\boldsymbol{\lambda}^{\mu}\|^2+\frac{\beta}{2}\|-\boldsymbol{y}_{0}+\boldsymbol{y}^{\mu}\|^2=\frac{1}{2}\Delta\,.\end{aligned}$$

Hence, we have:

$$\begin{aligned}\|\boldsymbol{\lambda}_{K+1}-\boldsymbol{\lambda}^{\mu}\|&\leq\sqrt{\beta\Delta},\\ \|-\boldsymbol{y}_{K+1}+\boldsymbol{y}^{\mu}\|&\leq\sqrt{\frac{\Delta}{\beta}}\,.\end{aligned}\tag{33}$$

Then from (16) and the convexity of f and g, we have:

$$\begin{aligned}&f\left(\boldsymbol{x}_{K+1}\right)-f\left(\boldsymbol{x}^{\mu}\right)+g\left(\boldsymbol{y}_{K+1}\right)-g\left(\boldsymbol{y}^{\mu}\right)+\langle\boldsymbol{\lambda}^{\mu},\boldsymbol{x}_{K+1}-\boldsymbol{y}_{K+1}\rangle\\ \leq&f_K\left(\boldsymbol{x}_{K+1}\right)-f_K\left(\boldsymbol{x}^{\mu}\right)+g\left(\boldsymbol{y}_{K+1}\right)-g\left(\boldsymbol{y}^{\mu}\right)+\langle\boldsymbol{\lambda}^{\mu},\boldsymbol{x}_{K+1}-\boldsymbol{y}_{K+1}\rangle\\ \leq&\frac{1}{\beta}\|\boldsymbol{\lambda}_{K+1}-\boldsymbol{\lambda}^{\mu}\|\|\boldsymbol{\lambda}_{K+1}-\boldsymbol{\lambda}_{K}\|+\|-\boldsymbol{y}_{K+1}+\boldsymbol{y}_{K}\|\|\boldsymbol{\lambda}_{K+1}-\boldsymbol{\lambda}_{K}\|\\ &+\beta\|-\boldsymbol{y}_{K+1}+\boldsymbol{y}_{K}\|\|-\boldsymbol{y}_{K+1}+\boldsymbol{y}^{\mu}\|\\ \leq&\frac{\Delta}{K+1}+\frac{2\Delta}{\sqrt{K+1}}\,.\end{aligned}\tag{34}$$

From Lemma 3, we have (28).

If in addition F is ξ-monotone on $\mathcal{C}_{=}$, using (24), we can obtain the following inequality similar to (32):

$$\begin{aligned}&c\frac{\sum_{k=0}^{K}\|\boldsymbol{x}_{k+1}-\boldsymbol{x}^{\mu}\|_2^{\xi}}{K+1}+\frac{1}{\beta}\|\boldsymbol{\lambda}_{K+1}-\boldsymbol{\lambda}_{K}\|^2+\beta\|-\boldsymbol{y}_{K+1}+\boldsymbol{y}_{K}\|^2\\ \leq&\frac{1}{K+1}\left(\frac{1}{\beta}\|\boldsymbol{\lambda}_{0}-\boldsymbol{\lambda}^{\mu}\|^2+\beta\|-\boldsymbol{y}_{0}+\boldsymbol{y}^{\mu}\|^2\right)\end{aligned}$$

By the convexity of $\|\cdot\|_2^{\xi}$ we have $c\|\hat{\boldsymbol{x}}_{K+1}-\boldsymbol{x}^{\mu}\|_2^{\xi}\leq\frac{\Delta}{K+1}$. And from (34) we can see that $c\|\boldsymbol{x}_{K+1}-\boldsymbol{x}^{\mu}\|_2^{\xi}\leq\frac{\Delta}{K+1}+\frac{2\Delta}{\sqrt{K+1}}$. ■

Theorem 8 *If F is monotone on $\mathcal{C}_{=}$, then for Algorithm 1, we have*

$$\left|f\left(\hat{\boldsymbol{x}}_{K+1}\right)+g\left(\hat{\boldsymbol{y}}_{K+1}\right)-f\left(\boldsymbol{x}^{\mu}\right)-g\left(\boldsymbol{y}^{\mu}\right)\right|\leq\frac{\Delta}{2\left(K+1\right)}+\frac{2\sqrt{\beta\Delta}\|\boldsymbol{\lambda}^{\mu}\|}{\beta\left(K+1\right)},$$

$$\|\hat{\boldsymbol{x}}_{K+1}-\hat{\boldsymbol{y}}_{K+1}\|\leq\frac{2\sqrt{\beta\Delta}}{\beta\left(K+1\right)}$$

where $\hat{\boldsymbol{x}}_{K+1}=\frac{1}{K+1}\sum_{k=1}^{K+1}\boldsymbol{x}_k$, $\hat{\boldsymbol{y}}_{K+1}=\frac{1}{K+1}\sum_{k=1}^{K+1}\boldsymbol{y}_k$, and $\Delta\triangleq\frac{1}{\beta}\|\boldsymbol{\lambda}_0-\boldsymbol{\lambda}^{\mu}\|^2+\beta\|\boldsymbol{y}_0-\boldsymbol{y}^{\mu}\|^2$.

Proof [Proof of Theorem 8] Summing (20) over $k = 0, 1, \ldots, K$, dividing both sides with $K+1$, and using the definitions of $\hat{\boldsymbol{x}}_{K+1}$ and $\hat{\boldsymbol{y}}_{K+1}$ and the convexity of f and g, we have

$$f\left(\hat{\boldsymbol{x}}_{K+1}\right) + g\left(\hat{\boldsymbol{y}}_{K+1}\right) - f\left(\boldsymbol{x}^{\mu}\right) - g\left(\boldsymbol{y}^{\mu}\right) + \langle \boldsymbol{\lambda}^{\mu}, \hat{\boldsymbol{x}}_{K+1} - \hat{\boldsymbol{y}}_{K+1} \rangle \leq \frac{\Delta}{2\left(K+1\right)}.$$

From (7) and (33), we have:

$$\begin{aligned}
\|\hat{\boldsymbol{x}}_{K+1} - \hat{\boldsymbol{y}}_{K+1}\| =& \frac{1}{\beta\left(K+1\right)} \|\sum_{k=0}^{K} (\boldsymbol{\lambda}_{k+1} - \boldsymbol{\lambda}_k)\| \\
=& \frac{1}{\beta\left(K+1\right)} \|\boldsymbol{\lambda}_{K+1} - \boldsymbol{\lambda}_0\| \\
\leq& \frac{1}{\beta\left(K+1\right)} \left(\|\boldsymbol{\lambda}_0 - \boldsymbol{\lambda}^{\mu}\| + \|\boldsymbol{\lambda}_{K+1} - \boldsymbol{\lambda}^{\mu}\|\right) \\
\leq& \frac{2\sqrt{\beta\Delta}}{\beta\left(K+1\right)}
\end{aligned}$$

Finally, from Lemma 3, the conclusion follows. ■

From Proposition 1 and the fact that $\mathcal{C}$ is compact we can see that $\lim_{\mu \to 0} \boldsymbol{x}^{\mu}$ and $\lim_{\mu \to 0} \boldsymbol{\lambda}^{\mu}$ exist and are finite. Let $\boldsymbol{x}^{\star} = \lim_{\mu \to 0} \boldsymbol{x}^{\mu}$ and $\boldsymbol{\lambda}^{\star} = \lim_{\mu \to 0} \boldsymbol{\lambda}^{\mu}$, then $\boldsymbol{x}^{\star} \in \mathcal{S}^{\star}_{\mathcal{C},F}$.

Theorem 9 *$\exists \tilde{\mu} > 0$, s.t. if $\mu_t < \tilde{\mu}$, then* $\left|F(\boldsymbol{x}^{(t)}_{K+1})^{\intercal}(\boldsymbol{x}^{(t)}_{K+1} - \boldsymbol{x}^{\star})\right| \leq 2(\frac{\Delta}{K+1} + \frac{2\Delta}{\sqrt{K+1}} + \|\boldsymbol{\lambda}^{\star}\|\sqrt{\frac{\Delta}{\beta(K+1)}})$, $\left|F(\boldsymbol{x}^{\star})^{\intercal}(\boldsymbol{x}^{(t)}_{K+1} - \boldsymbol{x}^{\star})\right| \leq 2(\frac{\Delta}{K+1} + \frac{2\Delta}{\sqrt{K+1}} + \|\boldsymbol{\lambda}^{\star}\|\sqrt{\frac{\Delta}{\beta(K+1)}})$ *and* $\left|F(\boldsymbol{x}^{\star})^{\intercal}(\hat{\boldsymbol{x}}^{(t)}_{K+1} - \boldsymbol{x}^{\star})\right| \leq 2(\frac{\Delta}{2(K+1)} + \frac{2\sqrt{\beta\Delta}\|\boldsymbol{\lambda}^{\star}\|}{\beta(K+1)})$, $\forall K \geq 0$.

Proof [Proof of Theorem 9.]

For simplicity we let $B(\boldsymbol{x})$ denote the log-barrier term $-\sum_{i=1}^{m} log(-\varphi_i(\boldsymbol{x}))$. From Theorem 7 and 8 we have

$$
\begin{aligned}
&\left|F(\boldsymbol{x}^{\mu})^{\intercal}(\boldsymbol{x}_{K+1}^{(t)}-\boldsymbol{x}^{\mu})+\mu(B(\boldsymbol{y}_{K+1}^{(t)})-B(\boldsymbol{y}^{\mu}))\right| \\
=&\left|f(\boldsymbol{x}_{K+1}^{(t)})-f(\boldsymbol{x}^{\mu})+g(\boldsymbol{y}_{K+1}^{(t)})-g(\boldsymbol{y}^{\mu})\right| \\
\leq&\frac{\Delta}{K+1}+\frac{2\Delta}{\sqrt{K+1}}+\|\boldsymbol{\lambda}^{\mu}\|\sqrt{\frac{\Delta}{\beta\,(K+1)}} \\
&\left|F(\boldsymbol{x}^{K+1})^{\intercal}(\boldsymbol{x}_{K+1}^{(t)}-\boldsymbol{x}^{\mu})+\mu(B(\boldsymbol{y}_{K+1}^{(t)})-B(\boldsymbol{y}^{\mu}))\right| \\
=&\left|f_K(\boldsymbol{x}_{K+1}^{(t)})-f_K(\boldsymbol{x}^{\mu})+g(\boldsymbol{y}_{K+1}^{(t)})-g(\boldsymbol{y}^{\mu})\right| \\
\leq&\frac{\Delta}{K+1}+\frac{2\Delta}{\sqrt{K+1}}+\|\boldsymbol{\lambda}^{\mu}\|\sqrt{\frac{\Delta}{\beta\,(K+1)}} \\
&\left|F(\boldsymbol{x}^{\mu})^{\intercal}(\hat{\boldsymbol{x}}_{K+1}^{(t)}-\boldsymbol{x}^{\mu})+\mu(B(\hat{\boldsymbol{y}}_{K+1}^{(t)})-B(\boldsymbol{y}^{\mu}))\right| \\
=&\left|f(\hat{\boldsymbol{x}}_{K+1}^{(t)})-f(\boldsymbol{x}^{\mu})+g(\hat{\boldsymbol{y}}_{K+1}^{(t)})-g(\boldsymbol{y}^{\mu})\right| \\
\leq&\frac{\Delta}{2\,(K+1)}+\frac{2\sqrt{\beta\Delta}\|\boldsymbol{\lambda}^{\mu}\|}{\beta\,(K+1)}.
\end{aligned}
$$

From Proposition 1 in App. B we know that when $\mu\to 0$, $\boldsymbol{x}^{\mu}\to\boldsymbol{x}^{\star}$, $\boldsymbol{\lambda}^{\mu}\to\boldsymbol{\lambda}^{\star}$ and $\mu(B(\boldsymbol{y}_{K+1}^{(t)})-B(\boldsymbol{y}^{\mu}))\to 0$, so $\exists\tilde{\mu}>0$, s.t. if $\mu_t<\tilde{\mu}$, then we have

$$
\begin{aligned}
&\left|F(\boldsymbol{x}^{\star})^{\intercal}(\boldsymbol{x}_{K+1}^{(t)}-\boldsymbol{x}^{\star})\right|\leq 2\left(\frac{\Delta}{K+1}+\frac{2\Delta}{\sqrt{K+1}}+\|\boldsymbol{\lambda}^{\star}\|\sqrt{\frac{\Delta}{\beta\,(K+1)}}\right) \\
&\left|F(\boldsymbol{x}_{K+1}^{(t)})^{\intercal}(\boldsymbol{x}_{K+1}^{(t)}-\boldsymbol{x}^{\star})\right|\leq 2\left(\frac{\Delta}{K+1}+\frac{2\Delta}{\sqrt{K+1}}+\|\boldsymbol{\lambda}^{\star}\|\sqrt{\frac{\Delta}{\beta\,(K+1)}}\right) \\
&\left|F(\boldsymbol{x}^{\star})^{\intercal}(\hat{\boldsymbol{x}}_{K+1}^{(t)}-\boldsymbol{x}^{\star})\right|\leq 2\left(\frac{\Delta}{2\,(K+1)}+\frac{2\sqrt{\beta\Delta}\|\boldsymbol{\lambda}^{\star}\|}{\beta\,(K+1)}\right).
\end{aligned}
$$

■

To make the dependencies of the constants clear, here we restate Theorem 1 and Theorem 2 and provide their proofs.

Theorem 10 (restatement of Theorem 1) *Given an operator $F:\mathcal{X}\to\mathbb{R}^n$ monotone on $\mathcal{C}_=$ (Def. 1), and either F is strictly monotone on $\mathcal{C}$ or one of φ_i is strictly convex. Assume there exists $r>0$ or $s>0$ such that F is star-ξ-monotone—as per Def. 1—on either $\hat{\mathcal{C}}_r$ or $\tilde{\mathcal{C}}_s$, resp.*

Let $\boldsymbol{x}_K^{(t)}$ and $\hat{\boldsymbol{x}}_K^{(t)}\triangleq\frac{1}{K}\sum_{k=1}^{K}\boldsymbol{x}_k^{(t)}$ denote the last and average iterate of Algorithm 1, respectively, run with sufficiently small μ_{-1}. Then there exists $K_0\in\mathbb{N}$, K_0 depends on r or s, $s.t.\forall K>K_0$, for all $t\in[T]$, we have that:

1. $\|\boldsymbol{x}_K^{(t)}-\boldsymbol{x}^{\star}\|\leq(\frac{4}{c}(\frac{\Delta}{K}+\frac{2\Delta}{\sqrt{K}}+\|\boldsymbol{\lambda}^{\star}\|\sqrt{\frac{\Delta}{\beta K}}))^{1/\xi}$.

2. *If in addition F is ξ-monotone on $C_=$, we have $\|\hat{\boldsymbol{x}}_K^{(t)} - \boldsymbol{x}^\star\| \leq (\frac{2\Delta}{cK})^{1/\xi}$.*

3. *Moreover, if F is L-Lipschitz on $\mathcal{C}_=$—as per Assumption 1— then $\mathcal{G}(\boldsymbol{x}_K^{(t)}, \mathcal{C}) \leq M_0(\frac{4}{c}(\frac{\Delta}{K} + \frac{2\Delta}{\sqrt{K}} + \|\boldsymbol{\lambda}^\star\|\sqrt{\frac{\Delta}{\beta K}}))^{1/\xi}$ and $\mathcal{G}(\hat{\boldsymbol{x}}_K^{(t)}, \mathcal{C}) \leq M_0(\frac{4}{c}(\frac{\Delta}{2K} + \frac{2\sqrt{\beta\Delta}\|\boldsymbol{\lambda}^\star\|}{\beta K}))^{1/\xi}$,*

Proof [Proof of Theorem 10.] Without loss of generality, we suppose that F is star-ξ-monotone on $\hat{\mathcal{C}}_r$. Since you $\phi(\boldsymbol{y}_k) \leq \mathbf{0}$, from (29) we can see that $\exists K_0 \in \mathbb{N}$, K_0 depends on r, $s.t. \forall K \geq K_0,\ \boldsymbol{x}_k \in \hat{\mathcal{C}}_r$, so if $\mu_t < \tilde{\mu}$ ($\tilde{\mu}$ is defined in Theorem 9), by Theorem 9 we have

$$
\begin{aligned}
c\|\boldsymbol{x}_{K+1}^{(t)} - \boldsymbol{x}^\star\|^\xi \leq & \left|(F(\boldsymbol{x}_{K+1}^{(t)})^\intercal - F(\boldsymbol{x}^\star)^\intercal)(\boldsymbol{x}_{K+1}^{(t)} - \boldsymbol{x}^\star)\right| \\
\leq & \left|F(\boldsymbol{x}^\star)^\intercal(\boldsymbol{x}_{K+1}^{(t)} - \boldsymbol{x}^\star)\right| + \left|F(\boldsymbol{x}_{K+1}^{(t)})^\intercal(\boldsymbol{x}_{K+1}^{(t)} - \boldsymbol{x}^\star)\right| \\
\leq & 4\left(\frac{\Delta}{K+1} + \frac{2\Delta}{\sqrt{K+1}} + \|\boldsymbol{\lambda}^\star\|\sqrt{\frac{\Delta}{\beta\,(K+1)}}\right).
\end{aligned}
$$

So we have

$$
\|\boldsymbol{x}_{K+1}^{(t)} - \boldsymbol{x}^\star\| \leq \left(\frac{4}{c}\left(\frac{\Delta}{K+1} + \frac{2\Delta}{\sqrt{K+1}} + \|\boldsymbol{\lambda}^\star\|\sqrt{\frac{\Delta}{\beta\,(K+1)}}\right)\right)^{1/\xi}.
$$

If in addition F is ξ-monotone on $C_=$, then from (30) we know that when $\hat{\mu}$ is small enough, we have

$$
\|\hat{\boldsymbol{x}}_{K+1}^{(t)} - \boldsymbol{x}^\star\| \leq \left(\frac{2\Delta}{c(K+1)}\right)^{1/\xi}.
$$

If F is L-Lipschitz on $\mathcal{C}_=$, then from Lemma 1 we can see that

$$
\begin{aligned}
\mathcal{G}(\boldsymbol{x}_K^{(t)}, \mathcal{C}) \leq M_0\left(\frac{4}{c}\left(\frac{\Delta}{K+1} + \frac{2\Delta}{\sqrt{K+1}} + \|\boldsymbol{\lambda}^\star\|\sqrt{\frac{\Delta}{\beta\,(K+1)}}\right)\right)^{1/\xi}, \\
\mathcal{G}(\hat{\boldsymbol{x}}_K^{(t)}, \mathcal{C}) \leq M_0\left(\frac{4}{c}\left(\frac{\Delta}{2\,(K+1)} + \frac{2\sqrt{\beta\Delta}\|\boldsymbol{\lambda}^\star\|}{\beta\,(K+1)}\right)\right)^{1/\xi},
\end{aligned}
$$

where M_0 is a linear function of L, see the proof of Lemma 1 in App.B.2. ■

Theorem 11 (restatement of Theorem 2) *Given an operator $F : \mathcal{X} \to \mathbb{R}^n$, assume: (i) F is monotone on $\mathcal{C}_=$, as per Def. 1; (ii) either F is strictly monotone on $\mathcal{C}$ or one of φ_i is strictly convex; and (iii) $\inf_{\boldsymbol{x} \in S \setminus \{\boldsymbol{x}^\star\}} F(\boldsymbol{x})^\intercal \frac{\boldsymbol{x} - \boldsymbol{x}^\star}{\|\boldsymbol{x} - \boldsymbol{x}^\star\|} = a > 0$, where $S \equiv \hat{\mathcal{C}}_r$ or $\tilde{\mathcal{C}}_s$. Let $\boldsymbol{x}_K^{(t)}$ and $\hat{\boldsymbol{x}}_K^{(t)} \triangleq \frac{1}{K}\sum_{k=1}^K \boldsymbol{x}_k^{(t)}$ denote the last and average iterate of Algorithm 1, respectively, run with sufficiently small μ_{-1}. Then there exists $K_0 \in \mathbb{N}$, K_0 depends on r or s, $s.t. \forall t \in [T]$, $\forall K > K_0$, we have that:*

1. $\left\|\boldsymbol{x}_K^{(t)} - \boldsymbol{x}^\star\right\| \leq \frac{2}{a}\left(\frac{\Delta}{K} + \frac{2\Delta}{\sqrt{K}} + \|\boldsymbol{\lambda}^\star\|\sqrt{\frac{\Delta}{\beta K}}\right)$.

2. *If in addition* $\inf_{\boldsymbol{x}\in S\setminus\{\boldsymbol{x}^\star\}} F(\boldsymbol{x}^\star)^\intercal \frac{\boldsymbol{x}-\boldsymbol{x}^\star}{\|\boldsymbol{x}-\boldsymbol{x}^\star\|} = b > 0$ *(with* $S \equiv \hat{\mathcal{C}}_r$ *or* $\tilde{\mathcal{C}}_s$*), then* $\left\|\hat{\boldsymbol{x}}_K^{(t)} - \boldsymbol{x}^\star\right\| \leq \frac{2}{b}\left(\frac{\Delta}{2K} + \frac{2\sqrt{\beta\Delta}\|\boldsymbol{\lambda}^\star\|}{\beta K}\right)$.

3. *Moreover, if* F *is* L*-Lipschitz on* $\mathcal{C}_=$*—as per Assumption 1—then* $\mathcal{G}(\boldsymbol{x}_K^{(t)}, \mathcal{C}) \leq \frac{2M_0}{a}\left(\frac{\Delta}{K} + \frac{2\Delta}{\sqrt{K}} + \|\boldsymbol{\lambda}^\star\|\sqrt{\frac{\Delta}{\beta K}}\right)$ *and* $\mathcal{G}(\hat{\boldsymbol{x}}_K^{(t)}, \mathcal{C}) \leq \frac{2M_0}{b}\left(\frac{\Delta}{2K} + \frac{2\sqrt{\beta\Delta}\|\boldsymbol{\lambda}^\star\|}{\beta K}\right)$.

Proof [Proof of Theorem 11] Without loss of generality, we suppose $\inf_{\boldsymbol{x}\in\hat{\mathcal{C}}_r\setminus\{\boldsymbol{x}^\star\}} F(\boldsymbol{x})^\intercal \frac{\boldsymbol{x}-\boldsymbol{x}^\star}{\|\boldsymbol{x}-\boldsymbol{x}^\star\|} = a > 0$. When $K \geq K_0$ (K_0 and $\tilde{\mu}$ are defined in the proof of Theorem 10 and in Theorem 9, resp.), by Theorem 9 we have that

$$\begin{aligned}\|\boldsymbol{x}_{K+1}^{(t)} - \boldsymbol{x}^\star\| \leq& \frac{1}{a}\left|F(\boldsymbol{x}_{K+1}^{(t)})^\intercal(\boldsymbol{x}_{K+1}^{(t)} - \boldsymbol{x}^\star)\right| \\ \leq& \frac{2}{a}\left(\frac{\Delta}{K+1} + \frac{2\Delta}{\sqrt{K+1}} + \|\boldsymbol{\lambda}^\star\|\sqrt{\frac{\Delta}{\beta(K+1)}}\right).\end{aligned}$$

Similarly, if $\inf_{\boldsymbol{x}\in\hat{\mathcal{C}}_r\setminus\{\boldsymbol{x}^\star\}} F(\boldsymbol{x}^\star)^\intercal \frac{\boldsymbol{x}-\boldsymbol{x}^\star}{\|\boldsymbol{x}-\boldsymbol{x}^\star\|} = b > 0$, we have that

$$\|\hat{\boldsymbol{x}}_{K+1}^{(t)} - \boldsymbol{x}^\star\| \leq \frac{2}{b}\left(\frac{\Delta}{2(K+1)} + \frac{2\sqrt{\beta\Delta}\|\boldsymbol{\lambda}^\star\|}{\beta(K+1)}\right).$$

If F is L-Lipschitz on $\mathcal{C}_=$, then from Lemma 1 we can see that

$$\begin{aligned}\mathcal{G}(\boldsymbol{x}_K^{(t)}, \mathcal{C}) \leq& \frac{2M_0}{a}\left(\frac{\Delta}{K+1} + \frac{2\Delta}{\sqrt{K+1}} + \|\boldsymbol{\lambda}^\star\|\sqrt{\frac{\Delta}{\beta(K+1)}}\right), \\ \mathcal{G}(\hat{\boldsymbol{x}}_K^{(t)}, \mathcal{C}) \leq& \frac{2M_0}{b}\left(\frac{\Delta}{2(K+1)} + \frac{2\sqrt{\beta\Delta}\|\boldsymbol{\lambda}^\star\|}{\beta(K+1)}\right).\end{aligned}$$

■

From the above proofs we can see that by setting μ_{-1} small enough, Theorem 1 and 2 hold true. But since we do not know exactly how small μ_{-1} should be, in practice, we can set a small μ_{-1}, and μ_t could be smaller than any fixed positive number after a very small number of outer loops, as μ_t decays exponentially. Thus, Algorithm 1 is actually parameter-free.

Note that the assumption (iii) in Theorem 2 is the weakening of the assumption $\inf_{\boldsymbol{x}\in S\setminus\{\boldsymbol{x}^\star\}} F(\boldsymbol{x}^\star)^\intercal \frac{\boldsymbol{x}-\boldsymbol{x}^\star}{\|\boldsymbol{x}-\boldsymbol{x}^\star\|} > 0$, where $S = \hat{\mathcal{C}}_r$ or $\tilde{\mathcal{C}}_s$. In fact, by the monotonicity of F, we have $F(\boldsymbol{x})^\intercal(\boldsymbol{x}-\boldsymbol{x}^\star) \geq F(\boldsymbol{x}^\star)^\intercal(\boldsymbol{x}-\boldsymbol{x}^\star)$, so if $\inf_{\boldsymbol{x}\in S\setminus\{\boldsymbol{x}^\star\}} F(\boldsymbol{x}^\star)^\intercal \frac{\boldsymbol{x}-\boldsymbol{x}^\star}{\|\boldsymbol{x}-\boldsymbol{x}^\star\|} > 0$, there must be $\inf_{\boldsymbol{x}\in S\setminus\{\boldsymbol{x}^\star\}} F(\boldsymbol{x})^\intercal \frac{\boldsymbol{x}-\boldsymbol{x}^\star}{\|\boldsymbol{x}-\boldsymbol{x}^\star\|} > 0$.

D.5. Discussion on Algorithm 1

Advantages and disadvantages of Algorithm 1. In Algorithm 1, the update of $\boldsymbol{x}$ (step 7) requires solving (W-EQ). Our method is especially suitable for problems where (W-EQ) is easy to solve

analytically. This includes the class of affine variational inequalities, low-dimensional problems, and when optimization variables represent probabilities, for example.

For problems where (W-EQ) is hard to solve—for example, min-max optimization problems in GANs—one could use other *unconstrained* methods like GDA or EG methods (*without* projection) so as to solve $VI(\mathbb{R}^n, G)$, where G is defined in (3). Algorithm 3 describes ACVI when using an unconstrained solver for the inner problems. We observe that Algorithm 3 outperforms the projection-based baseline algorithms, see for example Fig. 4.

Note that when there are constraints, the projection-based methods such as GDA, EG, OGDA etc. require solving a quadratic programming problem in each iteration (or twice per iteration for EG). This problem is often nontrivial and solving it may require using an interior point method or variants (such as the Frank-Wolfe algorithm). Because of this, Algorithm 1 can be seen as an orthogonal approach to projection-based methods, or in other words, as a complementary tool to solve (cVI) and particularly relevant when the constraints are non-trivial.

ACVI with only equality or inequality constraints. If there are no equality constraints, then $\mathcal{C}_=$ becomes $\mathbb{R}^n$. In this case, we have that $\boldsymbol{x}_{k+1}^{(t)}$ is the solution of:

$$\boldsymbol{x} + \frac{1}{\beta}(F(\boldsymbol{x}) + \boldsymbol{\lambda}_k^{(t)}) - \boldsymbol{y}_k^{(t)} = \mathbf{0}\,.$$

When there are no inequality constraints, we let $\boldsymbol{y}_k = \boldsymbol{x}_k$ and $\boldsymbol{\lambda}_k = \mathbf{0}$ for every k, and we can remove the outer loop. Thus, Algorithm 1 can be simplified to update only one variable $\boldsymbol{x}$ each iteration with the following updating rule:

$$\boldsymbol{x}_{k+1} \quad \textit{is the unique solution of} \quad \boldsymbol{x} + \frac{1}{\beta}\boldsymbol{P}_c F(\boldsymbol{x}) - \boldsymbol{P}_c \boldsymbol{x}_k - \boldsymbol{d}_c = \mathbf{0}\,.$$

Appendix E. Variant of the ACVI Algorithm (v-ACVI)

The presented approach of combining interior point methods with ADMM can be used as a framework to derive additional algorithms that may be more suitable for some specific problems. More precisely, observe from Eq. 1 that we could also consider a different splitting than that in § 3. Following this approach, we present a variant of Algorithm 1 and discuss its advantages and disadvantages relative to Algorithm 1.

E.1. Introduction of the variant ACVI

Deriving the v-ACVI algorithm. By considering a different splitting in (1) we get:

$$\begin{cases} \min\limits_{\boldsymbol{x},\boldsymbol{y}} F(\boldsymbol{w})^\intercal \boldsymbol{x} - \mu \sum\limits_{i=1}^{m} \log\big(-\varphi_i(\boldsymbol{x})\big) + \mathbb{1}[\boldsymbol{C}\boldsymbol{y} = \boldsymbol{d}] \\ s.t. \qquad \boldsymbol{x} = \boldsymbol{y} \end{cases},$$
$$\text{where:} \quad \mathbb{1}[\boldsymbol{C}\boldsymbol{y} = \boldsymbol{d}] = \begin{cases} 0, & \text{if } \boldsymbol{C}\boldsymbol{y} = \boldsymbol{d} \\ +\infty, & \text{if } \boldsymbol{C}\boldsymbol{y} \neq \boldsymbol{d} \end{cases}. \tag{35}$$

The augmented Lagrangian of (35) is thus:

$$\mathcal{L}_\beta(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{\lambda}) = F(\boldsymbol{w})^\intercal \boldsymbol{x} - \mu \sum_{i=1}^{m} \log(-\varphi_i(\boldsymbol{x})) + \mathbb{1}(\boldsymbol{C}\boldsymbol{y} = \boldsymbol{d}) + \langle \boldsymbol{\lambda}, \boldsymbol{x} - \boldsymbol{y} \rangle + \frac{\beta}{2} \|\boldsymbol{x} - \boldsymbol{y}\|^2 \tag{AL-CVI}$$

where $\beta > 0$ is the penalty parameter. We have that for $\boldsymbol{x}$ at step $k+1$:

$$\begin{aligned} \boldsymbol{x}_{k+1} &= \arg\min_{\boldsymbol{x}} \mathcal{L}_\beta(\boldsymbol{x}, \boldsymbol{y}_k, \boldsymbol{\lambda}_k) \\ &= \arg\min_{\boldsymbol{x}} \frac{1}{2} \left\| \boldsymbol{x} - \boldsymbol{y}_k + \frac{1}{\beta}(F(\boldsymbol{w}) + \boldsymbol{\lambda}_k) \right\|^2 - \frac{\mu}{\beta} \sum_{i=1}^{m} \log(-\varphi_i(\boldsymbol{x})) \end{aligned} \tag{36}$$

The following proposition ensures the existence and uniqueness of x_{k+1} in $\mathcal{C}_<$. i.e. We show that $\boldsymbol{x}_{k+1}$ is the unique solution in $\mathcal{C}_<$ of the following closed form equation (see App. E.2 for its proof):

$$\boldsymbol{x} - \boldsymbol{y}_k + \frac{1}{\beta}(F(\boldsymbol{w}) + \boldsymbol{\lambda}_k) - \frac{\mu}{\beta} \sum_{i=1}^{m} \frac{\nabla \varphi_i(\boldsymbol{x})}{\varphi_i(\boldsymbol{x})} = \boldsymbol{0}\,. \tag{X-CF}$$

Proposition 3 (unique solution) *The problem* (X-CF) *has a solution in $\mathcal{C}_<$ and the solution is unique.*

For $\boldsymbol{y}$, the updating rule is

$$\begin{aligned} \boldsymbol{y}_{k+1} &= arg\min_{\boldsymbol{y}} \mathcal{L}_\beta(\boldsymbol{x}_{k+1}, \boldsymbol{y}, \boldsymbol{\lambda}_k) \\ &= arg\min_{\boldsymbol{y} \in \mathcal{C}_=} -\frac{1}{\beta}(\boldsymbol{\lambda}_k)^\intercal \boldsymbol{y} + \frac{1}{2}\|\boldsymbol{y} - \boldsymbol{x}_{k+1}\|_2^2 \\ &= arg\min_{\boldsymbol{y} \in \mathcal{C}_=} \frac{1}{2}\|\boldsymbol{y} - \boldsymbol{x}_{k+1} - \frac{1}{\beta}\boldsymbol{\lambda}_k\|_2^2 \\ &= \boldsymbol{P}_c(\boldsymbol{x}_{k+1} + \frac{\boldsymbol{\lambda}_k}{\beta}) + \boldsymbol{d}_c \end{aligned} \tag{$\boldsymbol{y}$}$$

And the updating rule for dual variable $\boldsymbol{\lambda}$ is

$$\boldsymbol{\lambda}_{k+1} = \boldsymbol{\lambda}_k + \beta(\boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1}) \tag{$\boldsymbol{\lambda}$}$$

As in § 3, we would like to choose $\boldsymbol{w}_k$ so that $\boldsymbol{w}_k = \boldsymbol{x}_{k+1}$. To this end, we need the following assumption:

Assumption 2 *$\forall \boldsymbol{b} \in \mathbb{R}^n$ and $\mu > 0$, $\boldsymbol{x} + \frac{1}{\beta}F(\boldsymbol{x}) - \frac{\mu}{\beta}\sum_{i=1}^m \frac{\nabla\varphi_i(\boldsymbol{x})}{\varphi(\boldsymbol{x})} + \boldsymbol{b} = \boldsymbol{0}$ has a solution in $\mathcal{C}_<$.*

If Assumption 2 holds true, we can let $\boldsymbol{w}_k$ be a solution of

$$\boldsymbol{w} - \boldsymbol{y}_{k+1} + \frac{1}{\beta}(F(\boldsymbol{w}) + \boldsymbol{\lambda}_{k+1}) - \frac{\mu}{\beta}\sum_{i=1}^m \frac{\nabla\varphi_i(\boldsymbol{w})}{\varphi_i(\boldsymbol{w})} = \boldsymbol{0} \tag{37}$$

in $\mathcal{C}_<$. And by Proposition 3, we can let $\boldsymbol{x}_{k+1}$ be the unique solution of

$$\boldsymbol{x} - \boldsymbol{y}_k + \frac{1}{\beta}\big(F(\boldsymbol{w}_k) + \boldsymbol{\lambda}_k\big) - \frac{\mu}{\beta}\sum_{i=1}^m \frac{\nabla\varphi_i(\boldsymbol{x})}{\varphi_i(\boldsymbol{x})} = \boldsymbol{0} \tag{$\boldsymbol{x}$}$$

in $\mathcal{C}_<$.

Note that $\boldsymbol{w}_k$ is also a solution of ($\boldsymbol{x}$). By the uniqueness of the solution of ($\boldsymbol{x}$) shown in Prop. 3, we know that $\boldsymbol{w}_k = \boldsymbol{x}_{k+1}$.

So in the $(k+1)$-th step, we can compute $\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{\lambda}$ and $\boldsymbol{w}$ by ($\boldsymbol{x}$), ($\boldsymbol{y}$), ($\boldsymbol{\lambda}$) and (37), respectively. Since $\boldsymbol{w}_k = \boldsymbol{x}_{k+1}$, we can simplify our algorithm by removing variable $\boldsymbol{w}$ and only keep $\boldsymbol{x}, \boldsymbol{y}$ and $\boldsymbol{\lambda}$, see Algorithm 2.

Algorithm 2 v-ACVI pseudocode.

Input: operator $F : \mathcal{X} \to \mathbb{R}^n$, equality $\boldsymbol{C}\boldsymbol{x} = \boldsymbol{d}$ and inequality constraints $\varphi_i(\boldsymbol{x}) \le 0, i = [m]$, hyperparameters $\mu_{-1}, \beta > 0, \delta \in (0,1)$, number of outer and inner loop iterations T and K, resp.
Initialize: $\boldsymbol{y}_0^{(0)} \in \mathbb{R}^n$, $\boldsymbol{\lambda}_0^{(0)} \in \mathbb{R}^n$
$\boldsymbol{P}_c \triangleq \boldsymbol{I} - \boldsymbol{C}^\intercal(\boldsymbol{C}\boldsymbol{C}^\intercal)^{-1}\boldsymbol{C}$ *where $\boldsymbol{P}_c \in \mathbb{R}^{n\times n}$*
$\boldsymbol{d}_c \triangleq \boldsymbol{C}^\intercal(\boldsymbol{C}\boldsymbol{C}^\intercal)^{-1}\boldsymbol{d}$ *where $\boldsymbol{d}_c \in \mathbb{R}^n$*
for $t = 0, \dots, T-1$ **do**
  $\mu_t = \delta\mu_t$
  Denote $\hat{\varphi}(\boldsymbol{\lambda}, \boldsymbol{y})$ as the solution of $\frac{1}{\beta}\big(\mu_t \sum_{i=1}^m \frac{\nabla\varphi_i(\boldsymbol{x})}{\varphi_i(\boldsymbol{x})} - F(\boldsymbol{x}) - \boldsymbol{\lambda}\big) + \boldsymbol{y} - \boldsymbol{x} = 0$ with respect to $\boldsymbol{x}$, where $\boldsymbol{y}, \boldsymbol{\lambda}$ are variables
  for $k = 0, \dots, K-1$ **do**
    $\boldsymbol{x}_{k+1}^{(t)} = \hat{\varphi}(\boldsymbol{\lambda}_k^{(t)}, \boldsymbol{y}_k^{(t)})$ *Ensures $\boldsymbol{x}_{k+1} \in \mathcal{X}_\le$*
    $\boldsymbol{y}_{k+1}^{(t)} = \boldsymbol{P}_c(\boldsymbol{x}_{k+1}^{(t)} + \frac{\boldsymbol{\lambda}_k^{(t)}}{\beta}) + \boldsymbol{d}_c$
    $\boldsymbol{\lambda}_{k+1}^{(t)} = \boldsymbol{\lambda}_k^{(t)} + \beta(\boldsymbol{x}_{k+1}^{(t)} - \boldsymbol{y}_{k+1}^{(t)})$
  end for
  $(\boldsymbol{y}_0^{(t+1)}, \boldsymbol{\lambda}_0^{(t+1)}) \triangleq (\boldsymbol{y}_K^{(t)}, \boldsymbol{\lambda}_K^{(t)})$
end for

Discussion: ACVI Vs. v-ACVI. Relative to Algorithm 1, the subproblem for solving $\boldsymbol{x}$ in line 7 in Algorithm 2 becomes more complex, whereas the subproblem for $\boldsymbol{y}$ becomes simpler. Hence, in cases when the inequality constraints are simpler, or there are no inequality constraints Alg. 2 may be more suitable, as that simplifies the x subproblem. However, Algorithm 1 balances better the complexities of the subproblems, hence it may be simpler to use for general problems.
Convergence of v-ACVI. By analogous proofs to those in App. D, we can get similar convergence results as for Algorithm 1, that is Theorems 1 and 2. Specifically, Theorem 1 and 2 hold for Algorithm 2, provided that we replace the assumption "F is monotone on $\mathcal{C}_=$" with "F is monotone on $\mathcal{C}_\leq$" in Theorems 1 and 2.

E.2. Proof of Proposition 3

To prove proposition 3, we will use the following lemma.

Lemma 9 $\forall \boldsymbol{b} \in \mathbb{R}^n, \forall \boldsymbol{x} \in \mathcal{C}_<, \frac{1}{2}\|\boldsymbol{x}-\boldsymbol{b}\|_2^2 - \frac{\mu}{\beta}\sum_{i=1}^m \log(-\varphi_i(\boldsymbol{x})) \to +\infty, \|\boldsymbol{x}\|_2 \to +\infty$

Proof [Proof of Lemma 9.] We denote $\phi : \mathcal{C}_< \to \mathbb{R}$ by

$$\phi(\boldsymbol{x}) = \frac{1}{2}\|\boldsymbol{x}-\boldsymbol{b}\|_2^2 - \frac{\mu}{\beta}\sum_{i=1}^m \log(-\varphi_i(\boldsymbol{x}))$$

let $B(\boldsymbol{x}) = -\frac{\mu}{\beta}\sum_{i=1}^m \log(-\varphi_i(\boldsymbol{x}))$. We choose an arbitrary $\boldsymbol{x}_0 \in \mathcal{C}_<$. Then by the convexity of $B(\boldsymbol{x})$ we deduce that

$$\forall \boldsymbol{x} \in \mathcal{C}_<, \phi(\boldsymbol{x}) \geqslant \frac{1}{2}\|\boldsymbol{x}-\boldsymbol{b}\|_2^2 + B(\boldsymbol{x}_0) + \nabla B(\boldsymbol{x}_0)^\intercal(\boldsymbol{x}-\boldsymbol{x}_0) \to +\infty, \|\boldsymbol{x}\|_2 \to +\infty$$

■

In the remaining, we prove Proposition 3 which guarantees that (X-CF) has a unique solution.
Proof [Proof of Proposition 3: uniqueness of the solution of (X-CF).]
Let $\phi : \mathcal{C}_< \to \mathbb{R}$ denote:

$$\phi(\boldsymbol{x}) = \frac{1}{2}\|\boldsymbol{x}-\boldsymbol{y}_k + \frac{1}{\beta}(F(\boldsymbol{w})+\lambda_k)\|_2^2 - \frac{\mu}{\beta}\sum_{i=1}^m \log(-\varphi_i(\boldsymbol{x}))$$

We choose $\boldsymbol{x}_0 \in \mathcal{C}_<$. By Lemma 9, $\forall \boldsymbol{x} \in \mathcal{C}_<, \phi(\boldsymbol{x}) \to +\infty, \|\boldsymbol{x}\|_2 \to +\infty$.
So there exists $M > 0$ such that $\boldsymbol{x}_0 \in B(0, M)$ and $\forall \boldsymbol{x} \in S, \phi(\boldsymbol{x}) \leq \phi(\boldsymbol{x}_0), \boldsymbol{x}$ must belong to $B(0, M)$,where

$$B(0, M) = \{\boldsymbol{x} \in \mathbb{R}^n \,|\, \|\boldsymbol{x}\| \leq M\}$$

It's clear that there exists $t > 0$ such that for every $\boldsymbol{x} \in \mathcal{C}_<$ that satisfies $\phi(\boldsymbol{x}) \leq \phi(\boldsymbol{x}_0)$, $\boldsymbol{x}$ must belong to $\mathcal{C}_t$,where

$$\mathcal{C}_t = \{\boldsymbol{x} \in B(0, M) \,|\, \varphi_i(\boldsymbol{x}) \leq -t\} \subset \mathcal{C}_< \tag{38}$$

And we can make t small enough so that $\boldsymbol{x}_0 \in \mathcal{C}_t$. $\mathcal{C}_t$ is a nonempty compact set and ϕ is continuous, so there exists $\boldsymbol{x}^\star \in \mathcal{C}_t$ such that $\phi(\boldsymbol{x}^\star) \leq \phi(\boldsymbol{x}), \forall \boldsymbol{x} \in \mathcal{C}_<$.

Note that $\forall \boldsymbol{x} \in \mathcal{C}_< \backslash \mathcal{C}_t,\ \phi(\boldsymbol{x}) \geq \phi(\boldsymbol{x}_0) \geq \phi(\boldsymbol{x}^\star)$. Therefore, $\boldsymbol{x}^\star$ is a global minimizer of ϕ. ϕ is strongly-convex thus $\boldsymbol{x}_{k+1} = \boldsymbol{x}^\star$ is its unique minimizer. So $\boldsymbol{x} = \boldsymbol{x}_{k+1}$ if and only if $\nabla\phi(\boldsymbol{x}) = \boldsymbol{0}$. Therefore, $\boldsymbol{x}_{k+1}$ is the unique solution of $\boldsymbol{x} - \boldsymbol{y}_k + \frac{1}{\beta}(F(\boldsymbol{w}) + \boldsymbol{\lambda}_k) - \frac{\mu}{\beta}\sum_{i=1}^m \frac{\nabla\varphi_i(\boldsymbol{x})}{\varphi_i(\boldsymbol{x})} = \boldsymbol{0}$. ■

Appendix F. Details on the implementation

In this section, we provide the details on the implementation of the experiments shown in the main part in 2D and higher dimension bilinear game, see § F.1 and § F.2, respectively. We also provide here in § F.3 the details of the MNIST experiments presented later in App. G. In addition, we provide the source code through the following link: `https://anonymous.4open.science/r/ACVI`.

F.1. Experiments in 2D

We first state the considered problem fully, then describe the setup what includes the hyperparameters.

Problems. We consider the following constrained bilinear game (for the same experiment shown in Fig. 1 and 5):

$$\min_{x_1\in\mathbb{R}_+}\max_{x_2\in\mathbb{R}_+} 0.05x_1^2 + x_1x_2 - 0.05x_2^2\,. \tag{cBG}$$

The *Von Neumann's ratio game* [81] (results in Fig. 2(*a*)subfigure) is as follows:

$$\min_{\boldsymbol{x}\in\Delta^2}\max_{\boldsymbol{y}\in\Delta^2}\frac{\langle \boldsymbol{x}, \boldsymbol{R}\boldsymbol{y}\rangle}{\langle \boldsymbol{x}, \boldsymbol{S}\boldsymbol{y}\rangle}\,, \tag{RG}$$

where $\Delta^2 = \left\{\boldsymbol{z}\in\mathbb{R}^2 | \boldsymbol{z}\geq 0, \boldsymbol{e}^\intercal\boldsymbol{z} = 1\right\}$, $\boldsymbol{R} = \begin{pmatrix} -0.6 & -0.3 \\ 0.6 & -0.3\end{pmatrix}$ and $\boldsymbol{S} = \begin{pmatrix} 0.9 & 0.5 \\ 0.8 & 0.4\end{pmatrix}$.

The so called *Forsaken* [43] game—used in Fig. 2(*b*)subfigure and in App. G—is as follows:

$$\min_{x_1\in\mathcal{X}_1}\max_{x_2\in\mathcal{X}_2} x_1(x_2 - 0.45) + h(x_1) - h(x_2)\,, \tag{Forsaken}$$

where $h(z) = \frac{1}{4}z^2 - \frac{1}{2}z^4 + \frac{1}{6}z^6$. The original version is unconstrained $\mathcal{X}\equiv\mathbb{R}^2$. In Fig. 2(*b*)subfigure we use the constraint $x_1^2 + x_2^2 \leq 4$, and in App. G we use two other constraints: $x_1 \geq 0.08$ and $x_2 \geq 0.4$.

For the toy GAN experiments shown in Fig. 2(*c*)subfigure, the problem is as follows:

$$\begin{aligned}\min_{\theta}\max_{\varphi} \quad & \mathop{\mathbb{E}}_{x\sim\mathcal{N}(0,2)}(\varphi x^2) - \mathop{\mathbb{E}}_{z\sim\mathcal{N}(0,1)}(\varphi\theta^2 z^2) \\ \textit{s.t.} \quad & \varphi^2 + \theta^2 \leq 4\end{aligned} \tag{toy-GAN}$$

In the experiment, we look at a finite sum (sample average) approximation, which we then solve deterministically in a full batch fashion.

Setup. For all the 2D problems, we set the step size of GDA, EG and OGDA to 0.1, we use $k = 5$ and $\alpha = 0.5$ for LA-GDA, we set $\beta = 0.08$, $\mu_{-1} = 10^{-5}$, $\delta = 0.5$ and $\boldsymbol{\lambda}_0 = \boldsymbol{0}$ for ACVI; and run for 50 iterations. For ACVI, we set the number of outer loop $T = 20$. In the first 19 outer loop iterations, we only run one inner loop iteration, and in the last outer loop iteration, we run 30 inner loop iterations (for total of 50 updates). In Fig. 1 and Fig. 2(*c*)subfigure, we set the starting point for all algorithms. In Fig. 2(*a*)subfigure and 2(*b*)subfigure, we set the starting point to be $(0.5, 0.5)^\intercal$ for all algorithms.

F.2. High-dimension Bilinear Game

We also consider the following higher-dimension bilinear game on the probability simplex, with $\eta \in (0,1)$, $n = 1000$:

$$\min_{\boldsymbol{x}_1 \in \triangle} \max_{\boldsymbol{x}_2 \in \triangle} \eta \boldsymbol{x}_1^\intercal \boldsymbol{x}_1 + (1-\eta)\, \boldsymbol{x}_1^\intercal \boldsymbol{x}_2 - \eta \boldsymbol{x}_2^\intercal \boldsymbol{x}_2\,; \quad \triangle = \{\boldsymbol{x}_i \in \mathbb{R}^{500} | \boldsymbol{x}_i \geq \mathbf{0}, \text{ and }, \boldsymbol{e}^\intercal \boldsymbol{x}_i = 1\}. \tag{HBG}$$

We set the step size of GDA, EG and OGDA to 0.1, use $k = 4$ and $\alpha = 0.5$ for LA-GDA. For ACVI, we set $\beta = 0.5$, $\mu_{-1} = 10^{-6}$, $\delta = 0.5$ and $\boldsymbol{\lambda}_0 = \mathbf{0}$ for ACVI.

The solution of (HBG) is $\boldsymbol{x}^\star = \frac{1}{500}\boldsymbol{e}$, where $\boldsymbol{e} \in \mathbb{R}^{1000}$. As a metric of the experiments on this problem we use the relative error: $\varepsilon_r(\boldsymbol{x}_k) = \frac{\|\boldsymbol{x}_k - \boldsymbol{x}^\star\|}{\|\boldsymbol{x}^\star\|}$. In Fig.3(*b*)subfigure, we set $\varepsilon = 0.02$ and compute number of iterations of ACVI needed to reach the relative error given different rotation "strength"$1-\eta$, $\eta \in (0,1)$. Here we set the maximum numbers of iterations to be 50 for all algorithms. In Fig. 3(*a*)subfigure, we set $\eta = 0.05$ and compute CPU times needed for ACVI, EG, OGDA and LA4-GDA to reach different relative errors. Here we set the maximum run time to be 1500 seconds for all algorithms. In Fig. 7 in App. G on the other hand, we fix $\eta = 0.05$, and for varying *CPU time* limits, we compute the relative error of the last iterates of ACVI, GDA, EG, OGDA and LA4-GDA.

More general HD-BG game (g-HBG). Since (HBG) has perfect conditioning (that is, the interactive term $\boldsymbol{x}_1^\intercal \boldsymbol{B} \boldsymbol{x}_2$, is with $\boldsymbol{B} \equiv \boldsymbol{I}$), in App. G.2 we present results on the following more general high dimensional bilinear game:

$$\begin{aligned} \min_{\boldsymbol{x}_1 \in \triangle} \max_{\boldsymbol{x}_2 \in \triangle} \quad & \frac{\eta}{2} \cdot \mathbf{x}_1^\intercal \mathbf{A} \boldsymbol{x}_1 + (1-\eta)\, \boldsymbol{x}_1^\intercal \mathbf{B} \boldsymbol{x}_2 - \frac{\eta}{2} \boldsymbol{x}_2^\intercal \mathbf{C} \boldsymbol{x}_2, \\ & \triangle = \{\boldsymbol{x}_i \in \mathbb{R}^{500} | \boldsymbol{x}_i \geq -\boldsymbol{e}, \text{ and }, \boldsymbol{e}^\intercal \boldsymbol{x}_i = 0\}. \end{aligned} \tag{g-HBG}$$

Where $\mathbf{A}$, $\mathbf{B}$ and $\mathbf{C}$ are randomly generated 500×500 matrices, and $\mathbf{A}$, $\mathbf{C}$ are positive semi-definite.

The solution of (g-HBG) is $\boldsymbol{x}^\star = \mathbf{0}$, where $\mathbf{0} \in \mathbb{R}^{1000}$. As a metric of the experiments on this problem we use the error $\varepsilon(\boldsymbol{x}_k) = \|\boldsymbol{x}_k\|$. The remaining settings are identical to those of (HBG), explained above.

F.3. MNIST and Fashion-MNIST experiments

For the experiments on the MNIST dataset [1], we use the source code of Chavdarova et al. [14] for the baselines and we build on it to implement ACVI. For completeness, we provide an overview of the implementation.

Models. We used the DCGAN architectures [67], listed in Table 1, and the parameters of the models are initialized using PyTorch default initialization. For experiments on this dataset, we used the *non-saturating* GAN loss as proposed in [37]:

$$\mathcal{L}_D = \mathop{\mathbb{E}}_{\tilde{\boldsymbol{x}}_d \sim p_d} \log\big(D(\tilde{\boldsymbol{x}}_d)\big) + \mathop{\mathbb{E}}_{\tilde{\boldsymbol{z}} \sim p_z} \log\Big(1 - D\big(G(\tilde{\boldsymbol{z}})\big)\Big) \tag{L-D}$$

$$\mathcal{L}_G = \mathop{\mathbb{E}}_{\tilde{\boldsymbol{z}} \sim p_z} \log\Big(D\big(G(\tilde{\boldsymbol{z}})\big)\Big), \tag{L-G}$$

where $G(\cdot), D(\cdot)$ denote the generator and discriminator, resp., and p_d and p_z denote the data and the latent distributions (the latter predefined as normal distribution).

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Generator	**Discriminator**
Input: $\boldsymbol{z} \in \mathbb{R}^{128} \sim \mathcal{N}(0, I)$	*Input:* $\boldsymbol{x} \in \mathbb{R}^{1\times 28\times 28}$
transposed conv. (ker: 3×3, $128 \rightarrow 512$; stride: 1)	conv. (ker: 4×4, $1 \rightarrow 64$; stride: 2; pad:1)
Batch Normalization	LeakyReLU (negative slope: 0.2)
ReLU	conv. (ker: 4×4, $64 \rightarrow 128$; stride: 2; pad:1)
transposed conv. (ker: 4×4, $512 \rightarrow 256$, stride: 2)	Batch Normalization
Batch Normalization	LeakyReLU (negative slope: 0.2)
ReLU	conv. (ker: 4×4, $128 \rightarrow 256$; stride: 2; pad:1)
transposed conv. (ker: 4×4, $256 \rightarrow 128$, stride: 2)	Batch Normalization
Batch Normalization	LeakyReLU (negative slope: 0.2)
ReLU	conv. (ker: 3×3, $256 \rightarrow 1$; stride: 1)
transposed conv. (ker: 4×4, $128 \rightarrow 1$, stride: 2, pad: 1)	$Sigmoid(\cdot)$
$Tanh(\cdot)$	

Table 1: DCGAN architectures [67] used for experiments on **MNIST**. With "conv." we denote a convolutional layer and "transposed conv" a transposed convolution layer [67]. We use *ker* and *pad* to denote *kernel* and *padding* for the (transposed) convolution layers, respectively. With $h\times w$ we denote the kernel size. With $c_{in} \rightarrow y_{out}$ we denote the number of channels of the input and output, for (transposed) convolution layers. The models use Batch Normalization [44] layers.

Details on the ACVI implementation. When implementing ACVI on MNIST, we "remove" the outer loop of Algorithm 1 (that is we set $T = 1$), and fix μ to be a small number, in particular, we select $\mu = 10^{-9}$. We randomly initialize $\boldsymbol{x}$ and $\boldsymbol{y}$ and initialize $\boldsymbol{\lambda}$ to zero. For line 8 and 9 of Algorithm 1, we run l steps of (unconstrained) EG and GDA, respectively. For the update of $\boldsymbol{x}$ (using EG), we use step-size $\eta_x = 0.001$, whereas for $\boldsymbol{y}$ (using GDA), we use step-size $\eta_y = 0.2$. We present results when $l \in \{1, 10\}$. At every iteration, we update $\boldsymbol{\lambda}$ using the expression in line 11 of Algorithm 1, with $\beta = 0.5$.

Because the problem in step 7 of Algorithm 1 does not change a lot over the iterations (as well as when computing $\boldsymbol{y}$), when we implement Algorithm 1 we do *not* reinitialize the variable $\boldsymbol{x}$. We instead use the one from the previous iteration as an initialization and update it l times. The full details of the training when using an inner optimizer for step 7 of Algorithm 1 are provided in Algorithm 3, where we recall that $G(\boldsymbol{x})$ is defined as:

$$G(\boldsymbol{x}) \triangleq \boldsymbol{x} + \frac{1}{\beta}\boldsymbol{P}_c F(\boldsymbol{x}) - \boldsymbol{P}_c \boldsymbol{y}_k + \frac{1}{\beta}\boldsymbol{P}_c \boldsymbol{\lambda}_k - \boldsymbol{d}_c \tag{G-EQ}$$

Note that in the case of MNIST we consider only inequality constraints (and there are no equality constraints), therefore the matrices $\boldsymbol{P}_c$ and $\boldsymbol{d}_c$ are identity and zero, respectively. Therefore, there is no need for lines 3 and 4 in Algorithm 3.

Algorithm 3 Pseudocode for ACVI when using an inner optimizer (MNIST experiments).

Input: operator $F : \mathcal{X} \to \mathbb{R}^n$, equality $\boldsymbol{C}\boldsymbol{x} = \boldsymbol{d}$ and inequality constraints $\varphi_i(\boldsymbol{x}) \leq 0, i = [m]$, hyperparameters $\mu, \beta > 0, \delta \in (0, 1)$, inner optimizer $\mathcal{A}$ (e.g. EG, GDA, OGDA), l number of steps for the inner-optimizer, number of iterations K.
Initialize: $\boldsymbol{x}_0 \in \mathbb{R}^n$, $\boldsymbol{y}_0 \in \mathbb{R}^n$, $\boldsymbol{\lambda}_0 \in \mathbb{R}^n$
$\boldsymbol{P}_c \triangleq \boldsymbol{I} - \boldsymbol{C}^\intercal(\boldsymbol{C}\boldsymbol{C}^\intercal)^{-1}\boldsymbol{C}$ where $\boldsymbol{P}_c \in \mathbb{R}^{n\times n}$
$\boldsymbol{d}_c \triangleq \boldsymbol{C}^\intercal(\boldsymbol{C}\boldsymbol{C}^\intercal)^{-1}\boldsymbol{d}$ where $\boldsymbol{d}_c \in \mathbb{R}^n$
for $k = 0, \ldots, K-1$ **do**
  To obtain $\boldsymbol{x}_{k+1}$: run l steps of $\mathcal{A}$ solving VI$(\mathbb{R}^n, G)$, where G is defined in (G-EQ)
  To obtain $\boldsymbol{y}_{k+1}$: run l steps of GD to find $\boldsymbol{y}_{k+1}$ that minimizes:

$$-\mu \sum_{i=1}^{m} \log\big(-\varphi_i(\boldsymbol{y})\big) + \frac{\beta}{2}\left\|\boldsymbol{y} - \boldsymbol{x}_{k+1} - \frac{1}{\beta}\boldsymbol{\lambda}_k\right\|^2$$

  $\boldsymbol{\lambda}_{k+1} = \boldsymbol{\lambda}_k + \beta(\boldsymbol{x}_{k+1} - \boldsymbol{y}_{k+1})$
end for
Return: $\boldsymbol{x}_K$

The implementation details for the Fashion-MNIST experiment are identical to those on the MNIST experiment.

Setup 1: MNIST. The MNIST experiment in Fig. 4 in the main part (and in Fig. 9, 10) has 100 randomly generated linear inequality constraints for the Generator and 100 for the Discriminator.

Setup 1: projection details. Suppose the linear inequality constraints for the Generator are $\boldsymbol{A}\boldsymbol{\theta} \leq \boldsymbol{b}$, where $\boldsymbol{\theta} \in \mathbb{R}^n$ is the vector of all parameters of the Generator, $\boldsymbol{A} = (\boldsymbol{a}_1^\intercal, \cdots, \boldsymbol{a}_{100}^\intercal)^\intercal \in \mathbb{R}^{100\times n}$, $\boldsymbol{b} = (b_1, \ldots, b_{100}) \in \mathbb{R}^{100}$. We use the *greedy projection algorithm* described in [7]. A greedy projection algorithm is essentially a projected gradient method, it is easy to implement in high-dimension problems, and it has a convergence rate of $O(1/\sqrt{K})$. See Chapter 8.2.3 in [7] for more details. Since the dimension n is very large, at each step of the projection, one could only project $\boldsymbol{\theta}$ to one hyperplane $\boldsymbol{a}_i^\intercal\boldsymbol{x} = b_i$ for some $i \in \mathcal{I}(\boldsymbol{\theta})$, where

$$\mathcal{I}(\boldsymbol{\theta}) \triangleq \{j | \boldsymbol{a}_j^\intercal\boldsymbol{\theta} > b_j\}.$$

For every $j \in \{1, 2, \ldots, 100\}$, let

$$\mathcal{S}_j \triangleq \{\boldsymbol{x} | \boldsymbol{a}_j^\intercal\boldsymbol{x} \leq b_j\}.$$

The greedy projection method chooses i so that $i \in \arg\max\{dist(\boldsymbol{\theta}, \mathcal{S}_i)\}$. Note that as long as $\boldsymbol{\theta}$ is not in the constraint set $C_\leq = \{\boldsymbol{x} | \boldsymbol{A}\boldsymbol{x} \leq \boldsymbol{b}\}$, i would be in $\mathcal{I}(\boldsymbol{\theta})$. Algorithm 4 gives the details of the greedy projection method we use for the baseline, written for the Generator only for simplicity; the same projection method is used for the Discriminator as well.

Setup 2: MNIST & Fashion-MNIST. We add two constraints for the MNIST experiment: we set the squared sum of all parameters of the Generator and that of the Discriminator (separately) to be less than or equal to a hyperparameter M. We select some large number for M, in particular, we set $M = 50$.

Metrics. We describe the metrics for the MNIST experiments shown later in App. G. We use the two standard GAN metrics, Inception Score (IS, 73) and Fréchet Inception Distance (FID, 42). Both FID and IS rely on a pre-trained classifier, and take finite set of $\tilde{m}$ samples from the generator to compute these. Since **MNIST** has greyscale images, we used a classifier trained on this dataset and used $\tilde{m} = 5000$.

Algorithm 4 Greedy projection method for the baseline.

Input: $\boldsymbol{\theta} \in \mathbb{R}^n$, $\boldsymbol{A} = (\boldsymbol{a}_1^\intercal, \ldots, \boldsymbol{a}_{100}^\intercal)^\intercal \in \mathbb{R}^{100\times n}$, $\boldsymbol{b} = (b_1, \ldots, b_{100}) \in \mathbb{R}^{100}$, $\varepsilon > 0$
while True **do**
  $\mathcal{I}(\boldsymbol{\theta}) \triangleq \{j | \boldsymbol{a}_j^\intercal \boldsymbol{\theta} > b_j\}$
  if $\mathcal{I}(\boldsymbol{\theta}) = \emptyset$ or $\max\limits_{j\in\mathcal{I}(\boldsymbol{\theta})} \frac{|\boldsymbol{a}_j^\intercal \boldsymbol{\theta} - b_j|}{\|\boldsymbol{a}_j\|} < \varepsilon$ **then**
    break
  end if
  choose $i \in \arg\max\limits_{j\in\mathcal{I}(\boldsymbol{\theta})} \frac{|\boldsymbol{a}_j^\intercal \boldsymbol{\theta} - b_j|}{\|\boldsymbol{a}_j\|}$
  $\boldsymbol{\theta} \leftarrow \boldsymbol{\theta} - \frac{|\boldsymbol{a}_i^\intercal \boldsymbol{\theta} - b_i|}{\|\boldsymbol{a}_i\|^2} \boldsymbol{a}_i$
end while
Return: $\boldsymbol{\theta}$

Metrics: IS. Given a sample from the generator $\tilde{\boldsymbol{x}}_g \sim p_g$—where p_g denotes the data distribution of the generator—IS uses the softmax output of the pre-trained network $p(\tilde{\boldsymbol{y}}|\tilde{\boldsymbol{x}}_g)$ which represents the probability that $\tilde{\boldsymbol{x}}_g$ is of class $c_i, i \in 1 \ldots C$, i.e., $p(\tilde{\boldsymbol{y}}|\tilde{\boldsymbol{x}}_g) \in [0,1]^C$. It then computes the marginal class distribution $p(\tilde{\boldsymbol{y}}) = \int_{\tilde{\boldsymbol{x}}} p(\tilde{\boldsymbol{y}}|\tilde{\boldsymbol{x}}_g) p_g(\tilde{\boldsymbol{x}}_g)$. IS measures the Kullback–Leibler divergence $\mathbb{D}_{KL}$ between the predicted conditional label distribution $p(\tilde{\boldsymbol{y}}|\tilde{\boldsymbol{x}}_g)$ and the marginal class distribution $p(\tilde{\boldsymbol{y}})$. More precisely, it is computed as follows:

$$IS(G) = \exp\Big(\mathop{\mathbb{E}}_{\tilde{\boldsymbol{x}}_g\sim p_g}\Big[\mathbb{D}_{KL}\big(p(\tilde{\boldsymbol{y}}|\tilde{\boldsymbol{x}}_g)||p(\tilde{\boldsymbol{y}})\big)\Big]\Big) = \exp\Big(\frac{1}{\tilde{m}}\sum_{i=1}^{\tilde{m}}\sum_{c=1}^{C} p(y_c|\tilde{\boldsymbol{x}}_i)\log\frac{p(y_c|\tilde{\boldsymbol{x}}_i)}{p(y_c)}\Big). \quad \text{(IS)}$$

It aims at estimating (i) if the samples look realistic i.e., $p(\tilde{\boldsymbol{y}}|\tilde{\boldsymbol{x}}_g)$ should have low entropy, and (ii) if the samples are diverse (from different ImageNet classes) i.e., $p(\tilde{\boldsymbol{y}})$ should have high entropy. As these are combined using the Kullback–Leibler divergence, the higher the score is, the better the performance.

Metrics: FID. Contrary to IS, FID compares the synthetic samples $\tilde{\boldsymbol{x}}_g \sim p_g$ with those of the training dataset $\tilde{\boldsymbol{x}}_d \sim p_d$ in a feature space. The samples are embedded using the first several layers of a pretrained classifier. It assumes p_g and p_d are multivariate normal distributions, and estimates the means $\boldsymbol{m}_g$ and $\boldsymbol{m}_d$ and covariances $\boldsymbol{C}_g$ and $\boldsymbol{C}_d$, respectively for p_g and p_d in that feature space. Finally, FID is computed as:

$$\mathbb{D}_{\text{FID}}(p_d, p_g) \approx \mathscr{D}_2\big((\boldsymbol{m}_d, \boldsymbol{C}_d), (\boldsymbol{m}_g, \boldsymbol{C}_g)\big) = \|\boldsymbol{m}_d - \boldsymbol{m}_g\|_2^2 + Tr\big(\boldsymbol{C}_d + \boldsymbol{C}_g - 2(\boldsymbol{C}_d\boldsymbol{C}_g)^{\frac{1}{2}}\big), \quad \text{(FID)}$$

where $\mathscr{D}_2$ denotes the Fréchet Distance. Note that as this metric is a distance, the lower it is, the better the performance.

Hardware. We used the Colab platform (https://colab.research.google.com/) and *Tesla P100* GPUs. The running times are reported in App. G.

Appendix G. Additional Empirical Analysis

In this section, we provide some omitted plots/analyses of the results in the main paper as well as additional experiments. In particular, (i) App. G.1 lists results in 2D, (ii) App. G.2 on (HBG) and (g-HBG), whereas (iii) App. G.3 provides more detailed plots of those experiments summarized in Fig. 4, and presents additional results on other constraints on MNIST where we compare computationally-wise with *unconstained* baselines.

G.1. Additional experiments in 2D, on HBG and on g-HBG

Depicting the omitted baselines of Fig. 1. While Fig. 1 lists solely EG and ACVI for clarity, Fig. 5 depicts all the considered baselines in this paper on the cBG problem for completeness.

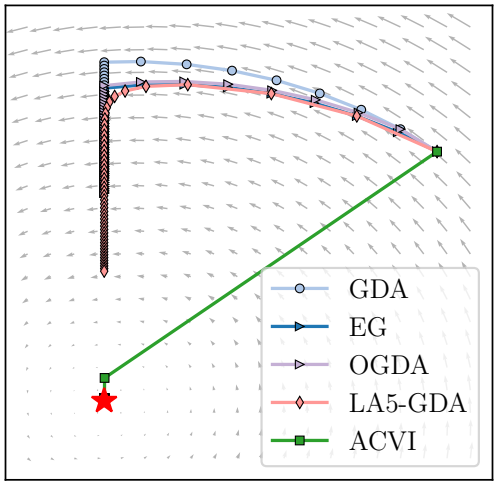
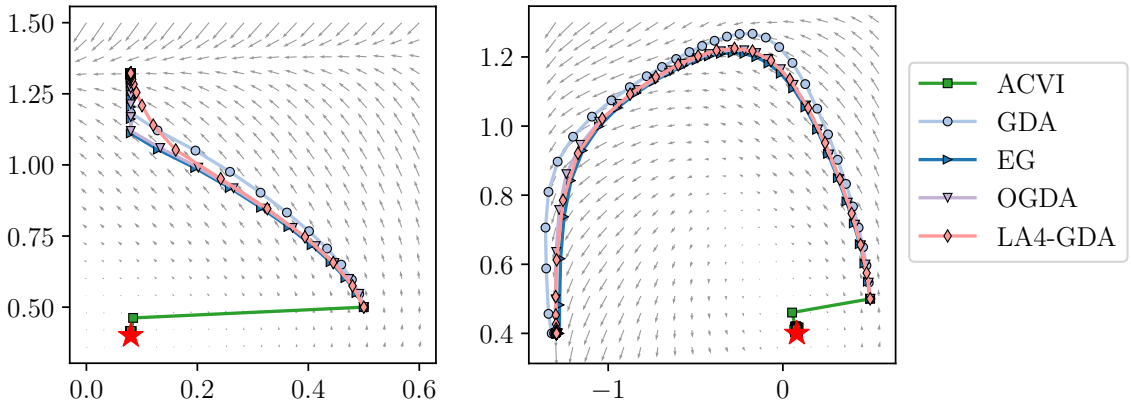


Figure 5: In addition to Fig. 1, here we depict all the considered baselines on the cBG problem. The solution at $(0, 0)$ is denoted with red star ($\star$), and the vector field of this problem with gray arrows. See App. F.1 for details on the implementation.

Additional experiments: varying constraints on the *Forsaken* problem. The *Forsaken* game was first pointed out in [43] and is particularly relevant because it has *limit cycles*, despite that it is in 2D. Since we are missing a tool to detect if we are in a limit cycle when in higher dimensions, this example is a popular benchmark in many recent works. Interestingly, in Fig. 2(*b*)subfigure we observe that ACVI is the only method that escapes the limit cycle. However, since in those simulations, given the initial point the constraints are not active throughout the training, in this section we run experiments with additional constraints. Fig. 6 depicts the baseline methods and ACVI on the Forsaken problem with two different constraints than that considered in Fig. 2 (that $x_1^2 + x_2^2 \leq 4$). Since this game is non-monotone, we observe that for some constraints the baseline methods—GDA, EG, OGDA, LA4-GDA—stay near the constraint (and do not converge). This may indicate that ACVI may have better chances of converging for *broader* problem classes than monotone VIs, relative to baseline methods whose convergence may depend on the constraints, and when hitting a constraint may be significantly slower (as Fig. 1 illustrates).



(*a*) $x_1 \geq 0.08$ (*b*) $x_2 \geq 0.4$

Figure 6: *Forsaken game with different constraints*: we consider two additional (to that in Fig. 2) constraints: 6(*a*)subfigure that $x_1 \geq 0.08$, and 6(*b*)subfigure that $x_2 \geq 0.4$. See App. F.1 for details on the implementation, and App. G.1 for a discussion.

G.2. Additional experiments HBG and on g-HBG

Complementary analysis to those in Fig. 3. Similar to Fig. 3, in Fig. 7 we run experiments on the HBG problem. However, here for a given fixed CPU time, we depict the relative error of the considered baselines and ACVI.

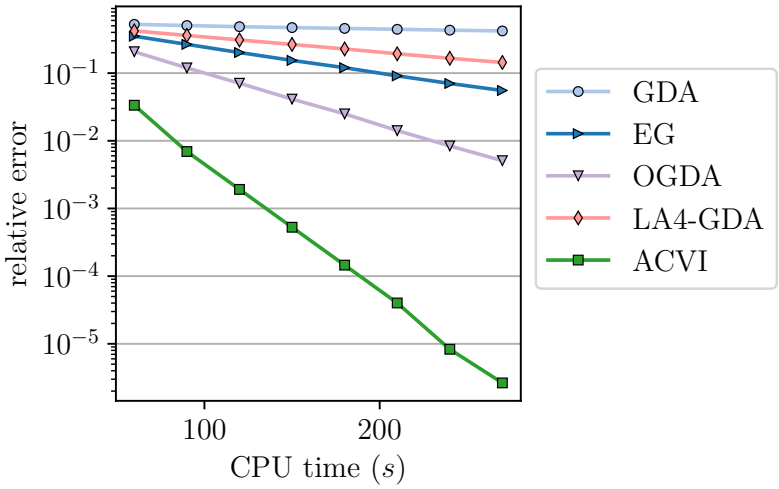
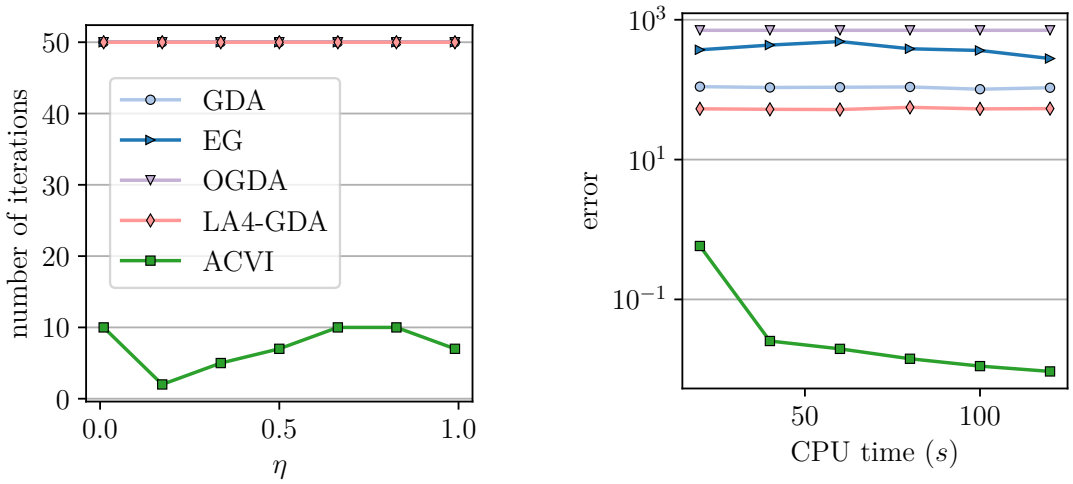


Figure 7: Given varying CPU time (in seconds), depicting the relative error (see App. F.2) of GDA, EG, OGDA, LA4-GDA, and ACVI (Algorithm 1) on the HBG problem where η is fixed to $\eta = .05$ (hence the vector field is highly rotational). This experiment complements those in Fig. 3 in the main paper. For the details on the implementation, see App. F.2.

Additional experiments on (g-HBG). In Fig. 8 we run experiments on the generalized HBG problem (g-HBG). In figure 8(*a*)subfigure, we compute number of iterations needed to reach ε-

distance to solution for varying intensity of the rotational component $(1-\eta)$; in figure 8(b)subfigure, we compute the error of the last iterate given fixed CPU time. We observe that despite the highly rotational monotone vector field, ACVI converges significantly faster in terms of wall clock time in higher dimensions as well.



(a) Varying rotational intensity $(1-\eta)$

(b) Achieved error given varying fixed CPU time

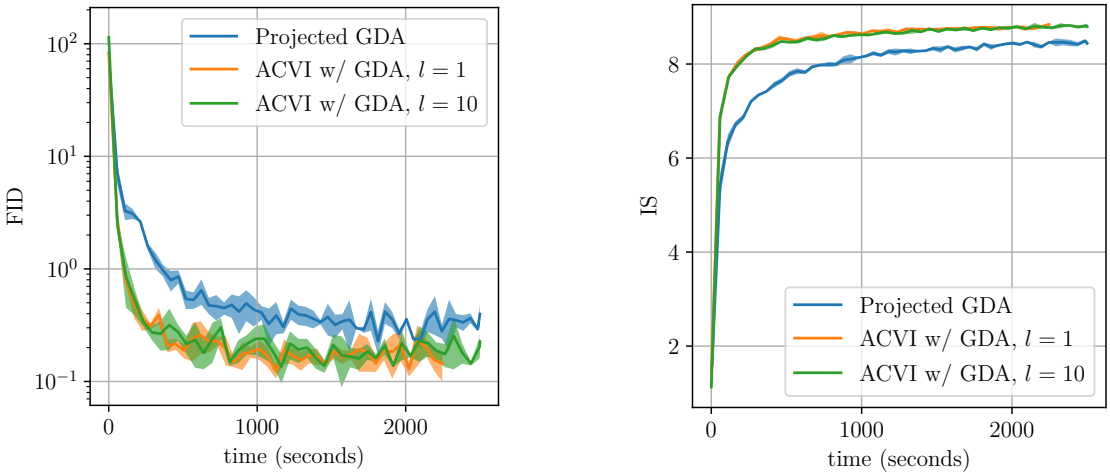
Figure 8: *General high-dimensional bilinear game* (g-HBG): comparison of ACVI with the GDA, EG, OGDA, and LA4-GDA baselines (described in App. B.5). **Left:** number of iterations (y-axis) needed to reach an ϵ-distance to the solution, for varying intensity of the rotational component $1-\eta$ (η is the x-axis) of the vector field (the smaller the η the higher the rotational component). We fix a threshold of the maximum number of iterations, and we stop the experiment. **Right:** distance to the solution (see App. F.2) of the last iterate (y-axis) for a varying wall-clock CPU time allowed to run each experiment (x-axis); in this experiment η is fixed to $\eta = 0.05$. See App. G.2 and F.2 for discussion and details on the implementation, respectively.

G.3. Experiments on MNIST and Fashion-MNIST

G.3.1. Setup 1: experiments on MNIST with linear inequalities

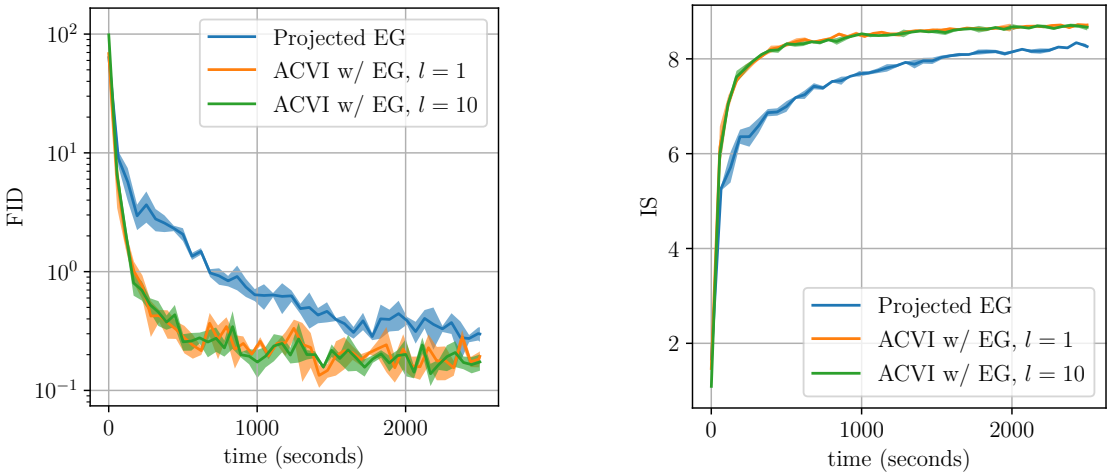
In this section, we present more detailed results of the summarizing plot in Fig. 4 of the main paper. For this experiment, we used linear inequalities as described in § F.3. Unlike in subsection G.3.2, here all the baselines are projected methods (that is, the same problem setting applies to ACVI and the baselines).

Fig. 9 and 10 depict the comparisons with projected GDA and projected EG, respectively. We observe that ACVI converges fast relative to the corresponding baseline. When choosing a larger number of steps for the inner problem $l = 10$ (see Algorithm 3) the wall-clock time per iteration increases, and interestingly the ACVI steps compensate for that and overall converge as fast as when $l = 1$.



(*a*) FID (lower is better) (*b*) IS (higher is better)

Figure 9: *Setup* 1*: Comparison between ACVI and GDA, and the* ***projected GDA*** on MNIST with linear inequalities (described in § F.3). l denotes the number of steps for the inner problems, see Algorithm 3. The depicted results are over multiple seeds. The FID and IS metrics as well as the implementation details are described in App. F.3. See App. G.3.1 for discussion.



(*a*) FID (lower is better) (*b*) IS (higher is better)

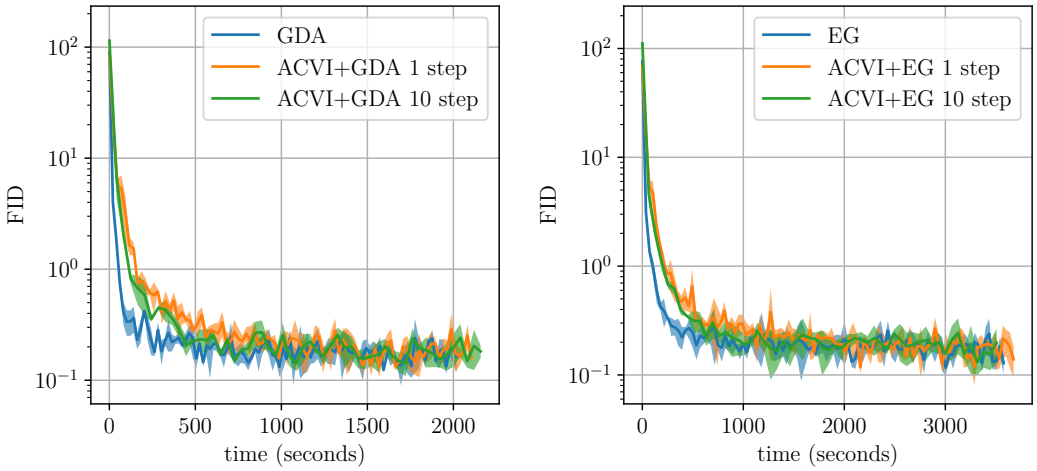
Figure 10: *Setup* 1*: Comparison between ACVI and EG, and the* ***projected EG*** on MNIST with linear inequalities (described in § F.3). l denotes the number of steps for the inner problems, see Algorithm 3. The depicted results are over multiple seeds. The FID and IS metrics as well as the implementation details are described in App. F.3. See App. G.3.1 for discussion.

G.3.2. SETUP 2: EXPERIMENTS ON (FASHION-)MNIST WITH QUADRATIC INEQUALITIES

In this section, we consider the MNIST and Fashion-MNIST datasets, which are unconstrained problems so as to make use of the well-established performance metrics (which are otherwise unclear

in the non-monotone settings, where we do not know the optimal solution apriori). We augment the problem with a mild constraint which requires that the norm of the per-player parameters does not exceed a certain value (see App. F.3). We compare ACVI with *unconstrained* baselines, what sets ACVI in a disadvantage as the projection requires additional computation. However, the primary purpose of these experiments is to observe if Algorithm 1 is competitive computationally-wise when lines 8 and 9 are non-trivial and require an (unconstrained) solver. However note that since MNIST is a relatively easy problem, it may not answer the natural question if ACVI has advantages on problems augmented with constraints over standard unconstrained methods. We leave such analyses for future work. The implementation and the used metrics are described in App. F.3.

Fig. 11 summarizes the experiments in terms of the obtained FID score over time. We observe that ACVI (although it uses two solvers at each iteration) is yet performing *competitively* to *unconstrained* GDA and EG. Figures 12–17 provide in addition samples of the Generator and IS scores, separately for each method. Figures 18 and 19 depict samples generated by the different methods, when trained on the Fashion-MNIST dataset. We believe that further exploring the type of constraints to be added, or the implementation options (e.g., l, step-size) may be proven fruitful even for problems that are originally unconstrained–as such an approach may reduce the rotational component of the original vector field, what in turn causes faster convergence or may help in escaping limit cycles for problems beyond monotone ones.



(*a*) GDA, and ACVI with GDA (*b*) EG, and ACVI with EG

Figure 11: Summary of the experiments on MNIST, using FID (lower is better). 11(*a*)subfigure: GDA and ACVI with GDA, and 11(*b*)subfigure: EG and ACVI with EG, using $l = \{1, 10\}$ for ACVI. Using step size of 0.001. The depicted results are over multiple seeds. See App. F.3 and G.3 for details on the implementation and discussion, resp.

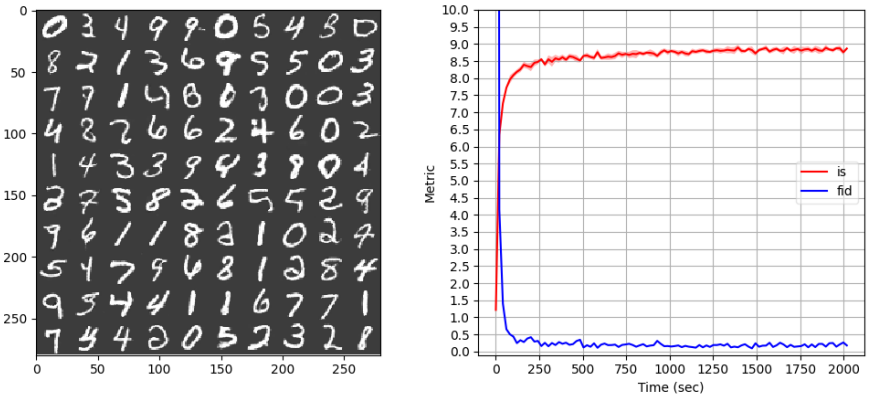


Figure 12: GDA on MNIST. Left: samples $\tilde{\boldsymbol{x}}_g \sim p_g$ of the last iterate of the Generator. Right: FID and IS of GDA, depicted in blue and red, respectively. Using step size of 0.001.

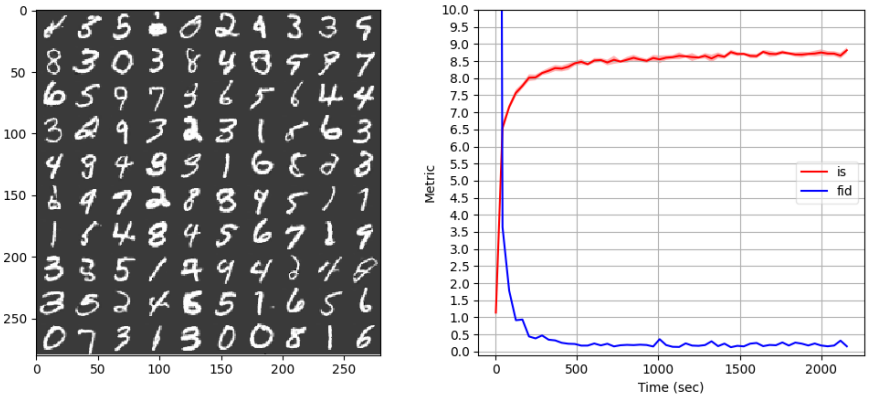


Figure 13: ACVI (Algorithm 1) with 10 GDA steps on MNIST. Left: samples $\tilde{\boldsymbol{x}}_g \sim p_g$ of the last iterate of the Generator. Right: FID and IS of GDA, depicted in blue and red, respectively. Using step size of 0.001 for $\boldsymbol{x}$ and 0.2 for $\boldsymbol{y}$, and $l = 10$ both for $\boldsymbol{x}$ and $\boldsymbol{y}$.

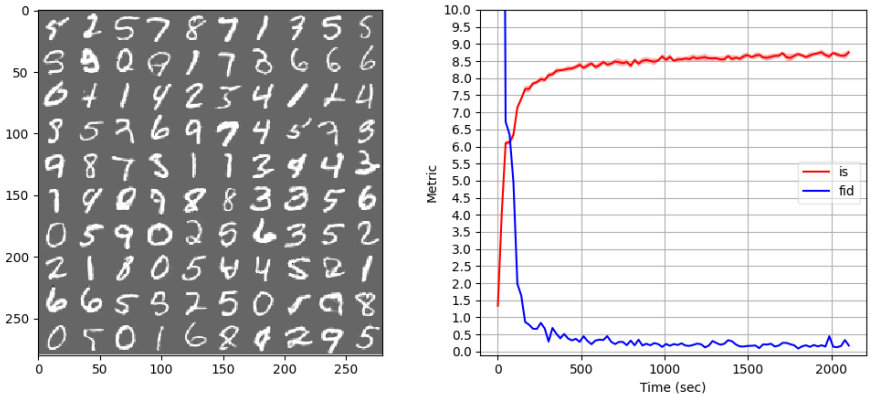


Figure 14: ACVI (Algorithm 1) with 1 GDA step on MNIST. Left: samples $\tilde{\boldsymbol{x}}_g \sim p_g$ of the last iterate of the Generator. Right: FID and IS of GDA, depicted in blue and red, respectively. Using step size of 0.001 for $\boldsymbol{x}$ and 0.2 for $\boldsymbol{y}$, and $l = 1$ both for $\boldsymbol{x}$ and $\boldsymbol{y}$.

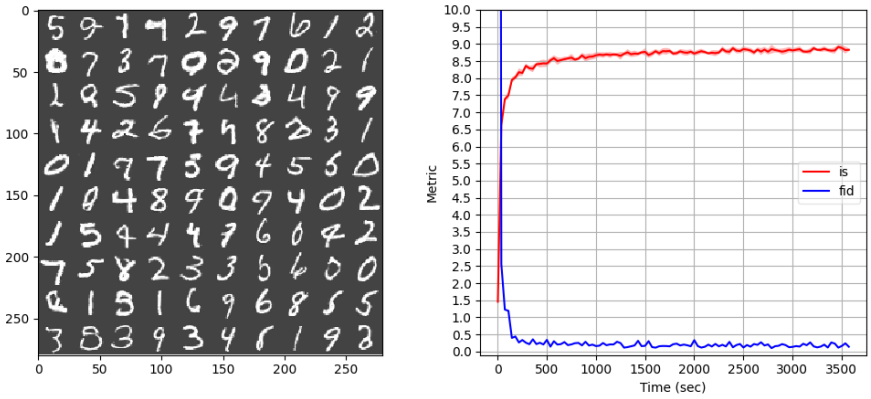


Figure 15: EG on MNIST. Left: samples $\tilde{\boldsymbol{x}}_g \sim p_g$ of the last iterate of the Generator. Right: FID and IS of EG, depicted in blue and red, respectively. Using step size of 0.001.

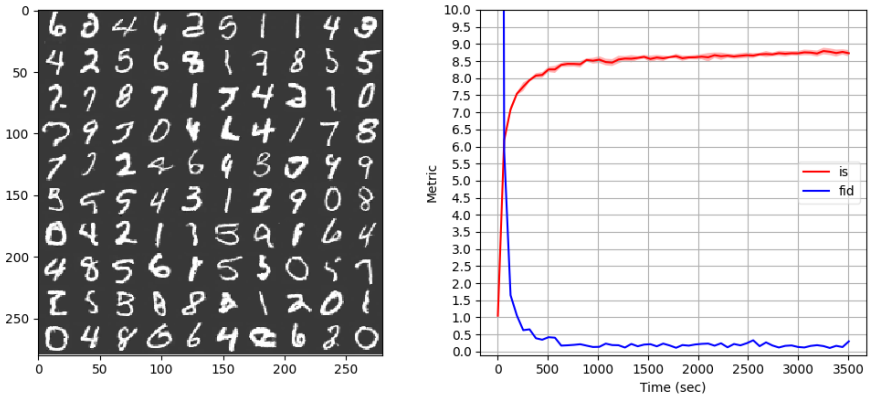


Figure 16: ACVI (Algorithm 1) with 10 EG steps on MNIST. Left: samples $\tilde{\boldsymbol{x}}_g \sim p_g$ of the last iterate of the Generator. Right: FID and IS of GDA, depicted in blue and red, respectively. Using step size of 0.001 for $\boldsymbol{x}$ and 0.2 for $\boldsymbol{y}$, and $l = 10$ both for $\boldsymbol{x}$ and $\boldsymbol{y}$.

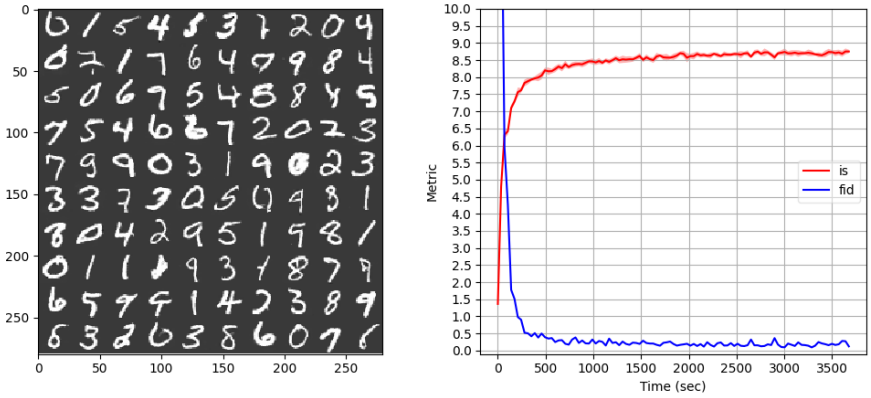
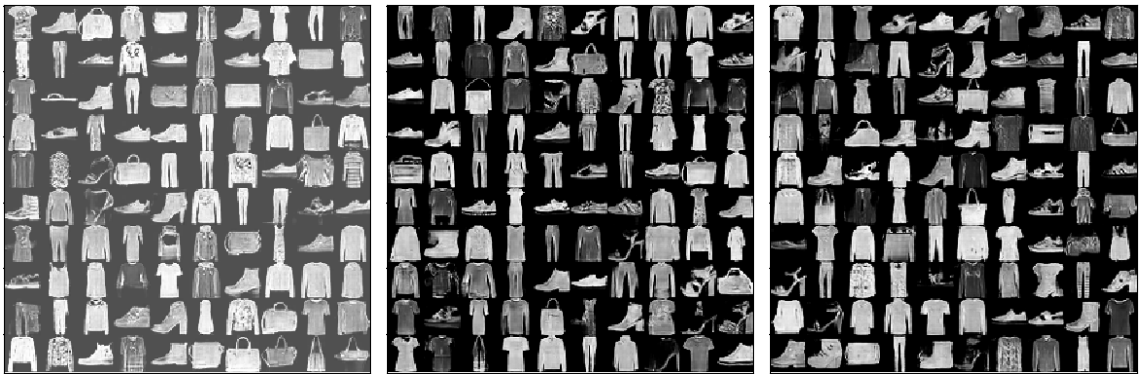


Figure 17: ACVI (Algorithm 1) with 1 EG step on MNIST. Left: samples $\tilde{\boldsymbol{x}}_g \sim p_g$ of the last iterate of the Generator. Right: FID and IS of GDA, depicted in blue and red, respectively. Using step size of 0.001 for $\boldsymbol{x}$ and 0.2 for $\boldsymbol{y}$, and $l = 1$ both for $\boldsymbol{x}$ and $\boldsymbol{y}$.



(*a*) EG (*b*) ACVI+EG 1 step (*c*) ACVI+EG 10 step

Figure 18: Generated images at fixed wall-clock computation time (3000s) by: the baseline EG, and by ACVI with $l \in \{1, 10\}$ on the Fashion-MNIST [86] dataset. See App. F.3 and G.3 for details on the implementation and discussion, resp.

(*a*) GDA (*b*) ACVI+GDA 1 step (*c*) ACVI+GDA 10 step

Figure 19: Generated images at fixed wall-clock computation time (3000s) by: the baseline GDA, and by ACVI with $l \in \{1, 10\}$ on the Fashion-MNIST [86] dataset. See App. F.3 and G.3 for details on the implementation and discussion, resp.