AMORTIZED IMPLICIT DIFFERENTIATION FOR STOCHASTIC BILEVEL OPTIMIZATION

ABSTRACT

We study a class of algorithms for solving bilevel optimization problems in both *stochastic* and *deterministic* settings when the inner-level objective is strongly convex. Specifically, we consider algorithms based on inexact implicit differentiation and we exploit a *warm-start strategy* to amortize the estimation of the exact gradient. We then introduce a unified theoretical framework inspired by the study of *singularly perturbed systems* (Habets, 1974) to analyze such *amortized* algorithms. By using this framework, our analysis shows these algorithms to match the computational complexity of oracle methods that have access to an unbiased estimate of the gradient, thus outperforming many existing results for bilevel optimization. We illustrate these findings on synthetic experiments and demonstrate the efficiency of these algorithms on hyper-parameter optimization experiments involving several thousands of variables.

1 INTRODUCTION

Bilevel optimization refers to a class of algorithms for solving problems with a hierarchical structure involving two levels: an *inner* and an *outer* level. The *inner*-level problem seeks a solution $y^\star(x)$ minimizing a cost $g(x, y)$ over a set $\mathcal{Y}$ given a fixed outer variable x in a set $\mathcal{X}$. The *outer*-level problem minimizes an objective of the form $\mathcal{L}(x) = f(x, y^\star(x))$ over $\mathcal{X}$ for some *upper*-level cost f. When the solution $y^\star(x)$ is unique, the bilevel optimization problem takes the following form:

$$\min_{x \in \mathcal{X}} \mathcal{L}(x) := f(x, y^\star(x)), \qquad \text{such that } y^\star(x) = \arg\min_{y \in \mathcal{Y}} g(x, y). \qquad (1)$$

First introduced in the field of economic game theory by Stackelberg (1934), this problem has recently received increasing attention in the machine learning community (Domke, 2012; Gould et al., 2016; Liao et al., 2018; Blondel et al., 2021; Liu et al., 2021; Shaban et al., 2019). Indeed, many machine learning applications can be reduced to (1) including hyper-parameter optimization (Feurer and Hutter, 2019), meta-learning (Bertinetto et al., 2018), reinforcement learning (Hong et al., 2020b; Liu et al., 2021) or dictionary learning (Mairal et al., 2011; Lecouat et al., 2020a;b).

The hierarchical nature of (1) introduces additional challenges compared to standard optimization problems, such as finding a suitable trade-off between the computational budget for approximating the inner and outer level problems (Ghadimi and Wang, 2018; Dempe and Zemkoho, 2020). These considerations are exacerbated in machine learning applications, where the costs f and g often come as an average of functions over a *large* or *infinite* number of data points (Franceschi et al., 2018). All these challenges highlight the need for methods that are able to control the computational costs inherent to (1) while dealing with the large-scale setting encountered in machine learning.

Gradient-based bilevel optimization methods appear to be viable approaches for solving (1) in large-scale settings (Lorraine et al., 2020). They can be divided into two categories: Iterative differentiation (ITD) and Approximate implicit differentiation (AID). *ITD approaches* approximate the map $y^\star(x)$ by a differentiable optimization algorithm $\mathcal{A}(x)$ viewed as a function of x. The resulting surrogate loss $\tilde{\mathcal{L}}(x) = f(x, \mathcal{A}(x))$ is optimized instead of $\mathcal{L}(x)$ using reverse-mode automatic differentiation (see Baydin et al., 2018). *AID approaches* (Pedregosa, 2016) rely on an expression of the gradient $\nabla \mathcal{L}$ resulting from the *implicit function theorem* (Lang, 2012, Theorem 5.9). Unlike ITD, AID avoids differentiating the algorithm approximating $y^\star(x)$ and, instead, approximately

solves a linear system using only Hessian and Jacobian-vector products to estimate the gradient $\nabla \mathcal{L}$ (Rajeswaran et al., 2019). These methods can also rely on stochastic approximation to increase scalability (Franceschi et al., 2018; Grazzi et al., 2020; 2021).

In the context of machine-learning, Ghadimi and Wang (2018) provided one of the first comprehensive studies of the computational complexity for a class of bilevel algorithms based on AID approaches. Subsequently, Hong et al. (2020b); Ji et al. (2021); Ji and Liang (2021); Yang et al. (2021) proposed different algorithms for solving (1) and obtained improved overall complexity by achieving a better trade-off between the cost of the inner and outer level problems. Still, the question of whether these complexities can be improved by better exploiting the structure of (1) through heuristics such as *warm-start* remains open (Grazzi et al., 2020). Moreover, these studies proposed separate analysis of their algorithms depending on the convexity of the loss $\mathcal{L}$ and whether a *stochastic* or *deterministic* setting is considered. This points out to a lack of unified and systematic theoretical framework for analyzing bilevel problems, which is what the present work addresses.

We consider the **Am**ortized **I**mplicit **G**radient **O**ptimization (AmIGO) algorithm, a bilevel optimization algorithm based on Approximate Implicit Differentiation (AID) approaches that exploits a warm-start strategy when estimating the gradient of $\mathcal{L}$. We then propose a unified theoretical framework for analyzing the convergence of AmIGO when the inner-level problem is strongly convex in both stochastic and deterministic settings. The proposed framework is inspired from the early work of Habets (1974) on *singularly perturbed systems* and analyzes the effect of warm start by viewing the iterates of AmIGO algorithm as a dynamical system. The evolution of such system is described by a total energy function which allows to recover the convergence rates of *unbiased oracle methods* which have access to an unbiased estimate of $\nabla \mathcal{L}$ (c.f. Table 1). To the best of our knowledge, this is the first time a bilevel optimization algorithm based on a warm-start strategy provably recovers the rates of *unbiased oracle methods* across a wide range of settings including the stochastic ones.

2 RELATED WORK

Singularly perturbed systems (SPS) are continuous-time deterministic dynamical systems of coupled variables $(x(t), y(t))$ with two time-scales where $y(t)$ evolves much faster than $x(t)$. As such, they exhibit a hierarchical structure similar to (1). The early work of Habets (1974); Saberi and Khalil (1984) provided convergence rates for SPS towards equilibria by studying the evolution of a single scalar energy function summarizing these systems. The present work takes inspiration from these works to analyze the convergence of AmIGO which involves three time-scales.

Two time-scale Stochastic Approximation (TTSA) can be viewed as a discrete-time stochastic version of SPS. (Kaledin et al., 2020) showed that TTSA achieves a finite-time complexity of $O(\epsilon^{-1})$ for linear systems while Doan (2020) obtained a complexity of $O(\epsilon^{-3/2})$ for general non-linear systems by extending the analysis for SPS. Hong et al. (2020b) further adapted the non-linear TTSA for solving (1). In the present work, we obtain faster rates by taking into account the dynamics of a third variable z_k appearing in AmIGO, thus resulting in a three time-scale dynamics.

Warm-start in bilevel optimization. Ji et al. (2021); Ji and Liang (2021) used a warm-start for the inner-level algorithm to obtain an improved computational complexity over algorithms without warm-start. In the deterministic non-convex setting, Ji et al. (2021) used a warm-start strategy when solving the linear system appearing in AID approaches to obtain improved convergence rates. However, it remained open whether using a warm-start when solving both inner-level problem and linear system arising in AID approaches can yield faster algorithms in the more challenging stochastic setting (Grazzi et al., 2020). In the present work, we provide a positive answer to this question.

3 AMORTIZED IMPLICIT GRADIENT OPTIMIZATION

3.1 GENERAL SETTING AND MAIN ASSUMPTIONS

Notations. In all what follows, $\mathcal{X}$ and $\mathcal{Y}$ are Euclidean spaces. For a differentiable function $h(x, y) : \mathcal{X} \times \mathcal{Y} \to \mathbb{R}$, we denote by ∇h its gradient w.r.t. (x, y), by $\partial_x h$ and $\partial_y h$ its partial derivatives w.r.t. x and y and by $\partial_{xy} h$ and $\partial_{yy} h$ the partial derivatives of $\partial_y h$ w.r.t x and y, respectively.

Geometries	Setting	Algorithms	Complexity
Strongly Convex	Deterministic	BA (Ghadimi and Wang, 2018)	$O(\kappa_{\mathcal{L}}^2 \vee \kappa_g^2 \log^2 \epsilon^{-1})$
		AccBio (Ji and Liang, 2021)	$O(\kappa_{\mathcal{L}}^{\frac{1}{2}} \kappa_g \log^2 \epsilon^{-1})$
		AmIGO (Corrolary 1)	$O(\kappa_{\mathcal{L}} \kappa_g \log \epsilon^{-1})$
	Stochastic	BSA (Ghadimi and Wang, 2018)	$O(\epsilon^{-2})$
		TTSA (Hong et al., 2020b)	$O(\epsilon^{-3/2} \log \epsilon^{-1})$
		AmIGO (Corrolary 2)	$O(\epsilon^{-1} \log \epsilon^{-1})$
Non-Convex	Deterministic	BA (Ghadimi and Wang, 2018)	$O(\kappa_{\mathcal{L}}^5 \epsilon^{-5/4})$
		AID-BiO (Ji et al., 2021)	$O(\kappa_{\mathcal{L}}^4 \epsilon^{-1})$
		AmIGO (Corrolary 3)	$O(\kappa_{\mathcal{L}}^4 \epsilon^{-1})$
	Stochastic	BSA (Ghadimi and Wang, 2018)	$O(\epsilon^{-3})$
		TTSA (Hong et al., 2020b)	$O(\epsilon^{-5/2} \log \epsilon^{-1})$
		stocBiO (Ji et al., 2021)	$O(\epsilon^{-2} \log \epsilon^{-1})$
		MRBO/VRBO* (Yang et al., 2021)	$O(\epsilon^{-3/2} \log \epsilon^{-1})$
		AmIGO (Corrolary 4)	$O(\epsilon^{-2})$

Table 1: Comparison of the cost of finding an ϵ-accurate solution as measured by $\mathbb{E}[\mathcal{L}(x_k) - \mathcal{L}^\star]$ when $\mathcal{L}$ is strongly-convex and $\frac{1}{k}\sum_{i=1}^{k}\mathbb{E}\left[\|\nabla\mathcal{L}(x_i)\|^2\right]$ when $\mathcal{L}$ is non-convex. The cost corresponds to the total number of gradients, Jacobian and Hessian-vector products used by the algorithm. $\kappa_{\mathcal{L}}$ and κ_g are the conditioning numbers of $\mathcal{L}$ and g whenever applicable. The rate of MRBO/VRBO is obtained under the additional *mean-squared smoothness* assumption (Arjevani et al., 2019).

To ensure that (1) is well-defined, we consider the setting where the inner-level problem is strongly convex so that the solution $y^\star(x)$ is unique as stated by the following assumption:

Assumption 1. *For any $x \in \mathcal{X}$, the function $y \mapsto g(x,y)$ is L_g-smooth and μ_g-strongly convex.*

Assumption 1 holds in the context of hyper-parameter selection when the inner-level is a kernel regression problem (Franceschi et al., 2018), or when the variable y represents the last *linear* layer of a neural network as in many meta-learning tasks (Ji et al., 2021). Under Assumption 1 and additional smoothness assumptions on f and g, the next proposition shows that $\mathcal{L}$ is differentiable:

Proposition 1. *Let g be a twice differentiable function satisfying Assumption 1. Assume that f is differentiable and consider the quadratic problem:*

$$\min_{z \in \mathbb{R}^{d_y}} Q(x,y,z) := \frac{1}{2} z^\top (\partial_{yy} g(x,y)) z + z^\top \partial_y f(x,y). \tag{2}$$

Then, (2) admits a unique minimizer $z^\star(x,y)$ for any (x,y) in $\mathcal{X} \times \mathcal{Y}$. Moreover, $y^\star(x)$ is unique and well-defined for any x in $\mathcal{X}$ and $\mathcal{L}$ is differentiable with gradient given by:

$$\nabla\mathcal{L}(x) = \partial_x f(x, y^\star(x)) + \partial_{xy} g(x, y^\star(x)) z^\star(x, y^\star(x)). \tag{3}$$

Proposition 1 follows by application of the *implicit function theorem* (Lang, 2012, Theorem 5.9) and provides an expression for $\mathcal{L}$ solely in terms of partial derivatives of f and g evaluated at $(x, y^\star(x))$. Following Ghadimi and Wang (2018), we further make two smoothness assumptions on f and g:

Assumption 2. *There exist positive constants L_f and B such that for all $x, x' \in \mathcal{X}$ and $y, y' \in \mathcal{Y}$:*

$$\|\nabla f(x,y) - \nabla f(x',y')\| \le L_f \|(x,y) - (x',y')\|, \qquad \|\partial_y f(x,y)\| \le B.$$

Assumption 3. *There exit positive constants L'_g, M_g such that for any $x, x' \in \mathcal{X}$ and $y, y' \in \mathcal{Y}$:*

$$\max\{\|\partial_{xy} g(x,y) - \partial_{xy} g(x',y')\|, \|\partial_{yy} g(x,y) - \partial_{yy} g(x',y')\|\} \le M_g \|(x,y) - (x',y')\|$$
$$\|\partial_y g(x,y) - \partial_y g(x',y)\| \le L'_g \|x - x'\|.$$

Assumptions 1 to 3 allow a control of the variations of $y^\star$ and $z^\star$ and ensure $\mathcal{L}$ is L-smooth for some positive constant L as shown in Proposition 6 of Appendix B.2. As an L-smooth function, $\mathcal{L}$ is necessarily weakly convex (Davis et al., 2018), meaning that $\mathcal{L}$ satisfies the inequality $\mathcal{L}(x) - \mathcal{L}(y) \le$

$\nabla \mathcal{L}(x)^\top (x - y) - \frac{\mu}{2} \|x - y\|^2$ for some fixed $\mu \in \mathbb{R}$ with $|\mu| \le L$. In particular, $\mathcal{L}$ is convex when $\mu \ge 0$, strongly convex when $\mu > 0$ and generally non-convex when $\mu < 0$. We thus consider two cases for $\mathcal{L}$, the *strongly convex* case ($\mu > 0$) and the *non-convex* case ($\mu < 0$). When $\mathcal{L}$ is convex, we denote by $\mathcal{L}^\star$ its minimum value achieved at a point $x^\star$ and define $\kappa_\mathcal{L} = L/\mu$ when $\mu > 0$.

Stochastic/deterministic settings. We consider the general setting where $f(x, y)$ and $g(x, y)$ are expressed as an expectation of stochastic functions $\hat{f}(x, y, \xi)$ and $\hat{g}(x, y, \xi)$ over a noise variable ξ. We recover the deterministic setting as a particular case when the variable ξ has zero variance, thus allowing us to treat both *stochastic* and *deterministic* settings in a unified framework. As often in machine-learning, we assume we can always draw a new batch $\mathcal{D}$ of i.i.d. samples of the noise variable ξ with size $|\mathcal{D}| \ge 1$ and use it to compute stochastic approximations of f and g defined by abuse of notation as $\hat{f}(x, y, \mathcal{D}) := \frac{1}{|\mathcal{D}|} \sum_{\xi \in \mathcal{D}} \hat{f}(x, y, \xi)$ and $\hat{g}(x, y, \mathcal{D}) := \frac{1}{|\mathcal{D}|} \sum_{\xi \in \mathcal{D}} \hat{g}(x, y, \xi)$. We make the following noise assumptions which are implied by those in Ghadimi and Wang (2018):

Assumption 4. *For any batch $\mathcal{D}$, $\nabla \hat{f}(x, y, \mathcal{D})$ and $\partial_y \hat{g}(x, y, \mathcal{D})$ are unbiased estimator of $\nabla f(x, y)$ and $\partial_y g(x, y)$ with a uniformly bounded variance, i.e. for all $x, y \in \mathcal{X} \times \mathcal{Y}$:*

$$\mathbb{E}\left[\left\| \nabla \hat{f}(x, y, \mathcal{D}) - \nabla f(x, y) \right\|^2\right] \le \tilde{\sigma}_f^2 |\mathcal{D}|^{-1}, \qquad \mathbb{E}\left[\left\| \partial_y \hat{g}(x, y, \mathcal{D}) - \partial_y g(x, y) \right\|^2\right] \le \tilde{\sigma}_g^2 |\mathcal{D}|^{-1}.$$

Assumption 5. *For any batch $\mathcal{D}$, the matrices $F_1(x, y, \mathcal{D}) := \partial_{xy} \hat{g}(x, y, \mathcal{D}) - \partial_{xy} g(x, y)$ and $F_2(x, y, \mathcal{D}) := \partial_{yy} \hat{g}(x, y, \mathcal{D}) - \partial_{yy} g(x, y)$ have zero mean and satisfy for all $x, y \in \mathcal{X} \times \mathcal{Y}$:*

$$\left\| \mathbb{E}\left[F_1(x, y, \mathcal{D})^\top F_1(x, y, \mathcal{D}) \right] \right\|_{op} \le \tilde{\sigma}_{g_{xy}}^2 |\mathcal{D}|^{-1}, \quad \left\| \mathbb{E}\left[F_2(x, y, \mathcal{D})^\top F_2(x, y, \mathcal{D}) \right] \right\|_{op} \le \tilde{\sigma}_{g_{yy}}^2 |\mathcal{D}|^{-1}.$$

For conciseness, we will use the notations $\sigma_f^2 := \tilde{\sigma}_f^2 |\mathcal{D}|^{-1}$, $\sigma_g^2 := \tilde{\sigma}_g^2 |\mathcal{D}|^{-1}$, $\sigma_{g_{xy}}^2 := \tilde{\sigma}_{g_{xy}}^2 |\mathcal{D}|^{-1}$ and $\sigma_{g_{yy}}^2 := \tilde{\sigma}_{g_{yy}}^2 |\mathcal{D}|^{-1}$, without explicit reference to the batch $\mathcal{D}$. Next, we describe the algorithm.

3.2 Algorithms

Amortized Implicit Gradient Optimization (AmIGO) is an iterative algorithm for solving (1). It constructs iterates x_k, y_k and z_k such that x_k approaches a stationary point of $\mathcal{L}$ while y_k and z_k track the quantities $y^\star(x_k)$ and $z^\star(x_k, y_k)$. AmIGO computes the iterate x_{k+1} using an update equation $x_{k+1} = x_k - \gamma_k \hat{\psi}_k$ for some given step-size γ_k and a stochastic estimate $\hat{\psi}_k$ of $\nabla \mathcal{L}(x_k)$ based on (3) and defined according to (4) below for some new batches of samples $\mathcal{D}_f$ and $\mathcal{D}_{g_{xy}}$.

$$\hat{\psi}_k := \partial_x \hat{f}(x_k, y_k, \mathcal{D}_f) + \partial_{x,y} \hat{g}(x_k, y_k, \mathcal{D}_{g_{xy}})^\top z_k. \tag{4}$$

Algorithm 1 AmIGO

1: Inputs: x_0, y_{-1}, z_{-1}.
2: Parameters: γ_k, K.
3: **for** $k \in \{0, ..., K\}$ **do**
4: $\quad y_k \leftarrow \mathcal{A}_k(x_k, y_{k-1})$
5: $\quad$ Sample batches $\mathcal{D}_f, \mathcal{D}_g$.
6: $\quad (u_k, v_k) \leftarrow \nabla \hat{f}(x_k, y_k, \mathcal{D}_f)$.
7: $\quad z_k \leftarrow \mathcal{B}_k(x_k, y_k, v_k, z_{k-1})$
8: $\quad w_k \leftarrow \partial_{xy} \hat{g}(x_k, y_k, \mathcal{D}_{g_{xy}}) z_k$
9: $\quad \hat{\psi}_{k-1} \leftarrow u_k + w_k$
10: $\quad x_k \leftarrow x_{k-1} - \gamma_k \hat{\psi}_{k-1}$
11: **end for**
12: Return x_K.

AmIGO computes $\hat{\psi}_k$ in 4 steps given iterates x_k, y_{k-1} and z_{k-1}. A first step computes an approximation y_k to $y^\star(x_k)$ using a stochastic algorithm $\mathcal{A}_k$ *initialized* at y_{k-1}. A second step computes unbiased estimates $u_k = \partial_x \hat{f}(x_k, y_k, \mathcal{D}_f)$ and $v_k = \partial_y \hat{f}(x_k, y_k, \mathcal{D}_f)$ of the partial derivatives of f w.r.t. x and y. A third step computes an approximation z_k to $z^\star(x_k, y_k)$ using a second stochastic algorithm $\mathcal{B}_k$ for solving (2) *initialized* at z_{k-1}. To increase efficiency, algorithm $\mathcal{B}_k$ uses the pre-computed vector v_k for approximating the partial derivative $\partial_y f$ in (2). Finally, the stochastic estimate $\hat{\psi}_k$ is computed using (4) by summing the pre-computed vector u_k with the jacobian-vector product $w_k = \partial_{xy} \hat{g}(x_k, y_k, \mathcal{D}_{g_{xy}}) z_k$. AmIGO is summarized in Algorithm 1.

Algorithms $\mathcal{A}_k$ and $\mathcal{B}_k$. While various choices for $\mathcal{A}_k$ and $\mathcal{B}_k$ are possible, such as adaptive algorithms (Kingma and Ba, 2015), or accelerated stochastic algorithms (Ghadimi and Lan, 2012), we focus on simple stochastic gradient descent algorithms with a pre-defined number of iterations T and N. These algorithms compute intermediate iterates y^t and z^n optimizing the functions $y \mapsto g(x_k, y)$ and $z \mapsto Q(x_k, y_k, z)$ starting from some initial values y^0 and z^0 and returning the last iterates y^T

and z^N as described in Algorithms 2 and 3. Algorithm $\mathcal{A}_k$ updates the current iterate y^{t-1} using a stochastic gradient $\partial_y \hat{g}(x_k, y^{t-1}, \mathcal{D}_g)$ for some new batch of samples $\mathcal{D}_g$ and a fixed step-size α_k. Algorithm $\mathcal{B}_k$ updates the current iterate z^{t-1} using a stochastic estimate of $\partial_z Q(x_k, y_k, z^{t-1})$ with step-size β_k. The stochastic gradient is computed by evaluating the Hessian-vector product $\partial_{yy} \hat{g}(x_k, y_k, \mathcal{D}_{g_{yy}}) z^{t-1}$ for some new batch of samples $\mathcal{D}_{g_{yy}}$ and summing it with a vector v_k approximating $\partial_y \hat{f}(x_k, y_k)$ provided as input to algorithm $\mathcal{B}_k$.

Warm-start for y^0 and z^0. Following the intuition that $y^\star(x_k)$ remains close to $y^\star(x_{k-1})$ when $x_k \simeq x_{k-1}$, and assuming that y_{k-1} is an accurate approximation to $y^\star(x_{k-1})$, it is natural to initialize $\mathcal{A}_k$ with the iterate y_{k-1}. The same intuition applies when initializing $\mathcal{B}_k$ with z_{k-1}. Next, we introduce a framework for analyzing the effect of warm-start on the convergence speed of AmIGO.

Algorithm 2 $\mathcal{A}_k(x, y^0)$	**Algorithm 3** $\mathcal{B}_k(x, y, v, z^0)$
1: Parameters: α_k, T	1: Parameters: β_k, N.
2: **for** $t \in \{1, ..., T\}$ **do**	2: **for** $n \in \{1, ..., N\}$ **do**
3: Sample batch $\mathcal{D}_{t,k}^g$.	3: Sample batch $\mathcal{D}_{n,k}^{g_{yy}}$.
4: $y^t \leftarrow y^{t-1} - \alpha_k \partial_y \hat{g}\left(x, y^{t-1}, \mathcal{D}_{t,k}^g\right)$.	4: $z^n \leftarrow z^{n-1} - \beta_k \left(\partial_{yy} \hat{g}\left(x, y, \mathcal{D}_{n,k}^{g_{yy}}\right) z^{n-1} + v\right)$.
5: **end for**	5: **end for**
6: Return y^T.	6: Return z^N.

4   ANALYSIS OF AMORTIZED IMPLICIT GRADIENT OPTIMIZATION

The proposed approach consists in three main steps: (1) Analysis of the outer-level problem , (2) Analysis of the inner-level problem and (3) Analysis of the joint dynamics of both levels.

Outer-level problem. We consider a quantity E_k^x describing the evolution of x_k defined as follows:

$$E_k^x := \begin{cases} \frac{\delta_k}{2\gamma_k} \mathbb{E}\left[\|x_k - x^\star\|^2\right] + (1-u)\mathbb{E}[\mathcal{L}(x_k) - \mathcal{L}^\star], & \mu \geq 0 \\ \frac{\delta_k}{2\gamma_k L^2} \mathbb{E}\left[\|\nabla \mathcal{L}(x_k)\|^2\right], & \mu < 0. \end{cases} \quad (5)$$

where δ_k is a positive sequence defined later in Section 4.1 that determines the convergence rate of the outer-level problem and $u \in [0,1]$ is a fixed parameter set to 1 in the *stochastic* setting and to 0 in the *deterministic* one. In Section 4.1, we control E_k^x using bias and variance error E_{k-1}^ψ and V_{k-1}^ψ of $\hat{\psi}_k$ given by (6) below where $\mathbb{E}_k$ denotes expectation conditioned on (x_k, y_k, z_{k-1}).

$$E_k^\psi := \mathbb{E}\left[\left\|\mathbb{E}_k\left[\hat{\psi}_k\right] - \nabla \mathcal{L}(x_k)\right\|^2\right], \qquad V_k^\psi := \mathbb{E}\left[\left\|\hat{\psi}_k - \mathbb{E}_k\left[\hat{\psi}_k\right]\right\|^2\right]. \quad (6)$$

Inner-level problems. We consider the mean-squared errors E_k^y and E_k^z between initializations ($y^0 = y_{k-1}$ and $z^0 = z_{k-1}$) and stationary values ($y^\star(x_k)$ and $z^\star(x_k, y_k)$) of algorithms $\mathcal{A}_k$ and $\mathcal{B}_k$:

$$E_k^y := \mathbb{E}\left[\left\|y_k^0 - y^\star(x_k)\right\|^2\right], \qquad E_k^z := \mathbb{E}\left[\left\|z_k^0 - z^\star(x_k, y_k)\right\|^2\right].$$

In Section 4.2, we show that the *warm-start strategy* allows to control E_k^y and E_k^z in terms of previous iterates E_{k-1}^y and E_{k-1}^z as well as the bias and variance errors in (6). We further prove that such bias and variance errors are, in turn, controlled by E_k^y and E_k^z.

Joint dynamics. In Section 4.3, following Habets (1974), we will consider, an aggregate error E_k^{tot} defined as a linear combination of E_k^x, E_k^y and E_k^z with carefully selected coefficients a_k and b_k:

$$E_k^{tot} = E_k^x + a_k E_k^y + b_k E_k^z. \quad (7)$$

As such E_k^{tot} represents the dynamics of the whole system. We then provide an error bound on E_k^{tot} from which we are able to deduce convergence rates for Algorithm 1 in various settings.

4.1   ANALYSIS OF THE OUTER-LEVEL PROBLEM

We start by defining the sequence δ_k appearing in (5) recursively in terms of the outer-level step-size γ_k and convexity constant μ. To this end, we choose γ_k to be a non-increasing sequence with

$\gamma_0 \leq \frac{1}{L}$ and define η_k and δ_k such that $\gamma_0^{-1} \geq \eta_0 \geq \max(\mu, 0)$ and

$$\eta_{k+1} = \eta_k \frac{1}{1 + \gamma_{k+1}(\eta_k - \mu)}, \qquad \delta_k := \begin{cases} \eta_k \gamma_k, & \mu \geq 0 \\ L\gamma_k, & \mu < 0 \end{cases}.$$

With this choice for δ_k and by setting $u = 1$ in (5), E_k^x recovers the quantity considered in the *stochastic estimate sequences* framework of Kulunchakov and Mairal (2020) to analyze the convergence of stochastic optimization algorithms when $\mathcal{L}$ is convex. When $\mathcal{L}$ is non-convex, E_k^x recovers a standard measure of stationarity (Davis and Drusvyatskiy, 2018). To allow a unified analysis of the dynamics of E_k^x in both convex and non-convex settings, we define F_k and s_k as follows:

$$F_k := u\delta_k \mathbb{E}[\mathcal{L}(x_k) - \mathcal{L}^\star]\mathbb{1}_{\mu > 0} + \left(\mathbb{E}[\mathcal{L}(x_k) - \mathcal{L}(x_{k-1})] + E_{k-1}^x - E_k^x\right)\mathbb{1}_{\mu < 0}$$

$$s_k := \frac{1}{2}((1 + u)\delta_k + 1 - u)\mathbb{1}_{\mu > 0} + \frac{\delta_k}{2}\mathbb{1}_{\mu < 0}.$$

The following proposition, with a proof in Appendix C.1, provides recursive inequality on E_k^x involving the errors in (6) due to the inexact gradient $\hat{\psi}_k$:

Proposition 2. *Let ρ_k be a non-increasing sequence with $0 < \rho_k < 1$. Assumptions 1 to 3 ensure that:*

$$F_k + E_k^x \leq (1 - (1 - \rho_k)\delta_k)E_{k-1}^x + \gamma_k s_k V_{k-1}^\psi + \gamma_k\left(s_k + (2\rho_k)^{-1}\right)E_{k-1}^\psi. \tag{9}$$

When the bias E_k^ψ vanishes, (9) simplifies to (Kulunchakov and Mairal, 2019, Proposition 1) which recovers standard convergence rates for stochastic gradient methods in the convex case. Next, we analyze the case when y_k and z_k are inexact thus yielding a non-zero bias E_k^ψ.

4.2 ANALYSIS OF THE INNER-LEVEL PROBLEM

We start with a standard result controlling the expected mean squared errors between iterates y_k and z_k and their limiting values $y^\star(x_k)$ and $z^\star(x_k, y_k)$ which is essential in our analysis.

Proposition 3. *Let α_k and β_k be two positive sequences with $\alpha_k \leq L_g^{-1}$ and $\beta_k \leq L_g^{-1}\left(1 + \mu_g^{-2}\sigma_{g_{yy}}^2\right)$ and define $\Lambda_k := (1 - \alpha_k\mu_g)^T$ and $\Pi_k := \left(1 - \frac{\beta_k\mu_g}{2}\right)^N$. Under Assumptions 1, 4 and 5, it holds that:*

$$\mathbb{E}\left[\|y_k - y^\star(x_k)\|^2\right] \leq \Lambda_k E_k^y + R_k^y, \qquad \mathbb{E}\left[\|z_k - z^\star(x_k, y_k)\|^2\right] \leq \Pi_k E_k^z + R_k^z, \tag{10}$$

where $R_k^y = O\left(\kappa_g\sigma_g^2\right)$ and $R_k^z = O\left(\kappa_g^3\sigma_{g_{yy}}^2 + \kappa_g^2\sigma_f^2\right)$ are defined in (15) of Appendix A.2.

While Proposition 3 is specific to the choice of the algorithms $\mathcal{A}_k$ and $\mathcal{B}_k$ in Algorithms 2 and 3, our analysis directly extends to other algorithms satisfying inequalities similar to (10). We can now formalize the intuition discussed in Section 3.2 about the effect of the warm-start by providing a recursive inequality satisfied by E_k^y and E_k^z in Proposition 4 with a proof in Appendix C.2.

Proposition 4. *Let $\zeta_k > 0$, a 2×2 matrix $\boldsymbol{P_k}$, two vectors $\boldsymbol{U_k}$ and $\boldsymbol{V_k}$ in $\mathbb{R}^2$ all independent of x_k, y_k and z_k be as defined in Proposition 8 of Appendix A.2. Under Assumptions 1 to 5, it holds that:*

$$\begin{pmatrix} E_k^y \\ E_k^z \end{pmatrix} \leq \boldsymbol{P_k} \begin{pmatrix} \Lambda_{k-1}E_{k-1}^y + R_{k-1}^y \\ \Pi_{k-1}E_{k-1}^z + R_{k-1}^z \end{pmatrix} + 2\left(E_{k-1}^\psi + V_{k-1}^\psi + \zeta_k E_{k-1}^x\right)\boldsymbol{U_k} + \boldsymbol{V_k}. \tag{11}$$

Proposition 4 follows from a first order expansion of E_k^y and E_k^z around previous iterates E_{k-1}^y and E_{k-1}^z with the residual error controlled by bounding the difference between two successive iterates of x_k and y_k. Finally, the following proposition, proved in Appendix C.1, controls the bias and variance terms V_k^ψ and E_k^ψ appearing in both (9) and (19) in terms of E_k^y and E_k^z:

Proposition 5. *Under Assumptions 1 to 5, the following inequalities hold:*

$$E_k^\psi \leq 2L_\psi^2(\Lambda_k E_k^y + \Pi_k E_k^z + R_k^y), \qquad V_k^\psi \leq w_x^2 + \sigma_x^2\Pi_k E_k^z, \tag{12}$$

where $w_x^2 = O\left(\kappa_g^2\left(\sigma_f^2 + \sigma_{g_{xy}}^2\right) + \kappa_g^3\sigma_{g_{yy}}^2\right)$, $\sigma_x^2 = O\left(\sigma_{g_{xy}}^2 + \kappa_g^2\sigma_{g_{yy}}^2\right)$ and $L_\psi = O(\kappa_g^2)$ are positive constants defined in (14) and (17) of Appendix A.2 with L_ψ controlling the variations of $\mathbb{E}_k[\hat{\psi}_k]$.

4.3 GENERAL ERROR BOUND

The joint dynamics of (E_k^x, E_k^y, E_k^z) is challenging to analyze directly since it is described by a coupled system of inequalities in (9), (12) and (19). Instead, we analyze the dynamics of the aggregate error E_k^{tot} defined in (7) by linearly combining the set of inequalities satisfied by (E_k^x, E_k^y, E_k^z). The following proposition provides an error bound for E_k^{tot} in both convex and non-convex settings for a suitable choice of the coefficients a_k and b_k provided that T and N are large enough:

Theorem 1. *Choose constant step-sizes* $\alpha_k = \frac{1}{L_g}$, $\beta_k = \left(L_g \left(1 + \mu_g^{-2} \sigma_{g_{yy}}^2 \right) \right)^{-1}$ *and* $\gamma_k = \frac{1}{L}$. *Set the coefficients* a_k *and* b_k *to be* $a_k := \gamma \Lambda_k^{1/2}$ *and* $b_k := \gamma \Pi_k^{1/2}$ *and set the number of iterations* T *and* N *of Algorithms 2 and 3 to be of order* $T = O(\kappa_g)$ *and* $N = O(\kappa_g(1 + \kappa_g^2 \sigma_{g_{yy}}^2))$ *up to a logarithmic dependence on* κ_g. *Let* $\hat{x}_k = u(1 - \delta_k)\hat{x}_{k-1} + (1 - u(1 - \delta_k))x_k$, *with* $\hat{x}_0 = x_0$. *Then, under Assumptions 1 to 5, E_k^{tot} satisfies:*

$$
\begin{cases}
\mathbb{E}[\mathcal{L}(\hat{x}_k) - \mathcal{L}^\star] + E_k^{tot} \leq \left(1 - (2\kappa_\mathcal{L})^{-1} \right)^k \left[E_0^{tot} + u^k \mathbb{E}[\mathcal{L}(x_0) - \mathcal{L}^\star] \right] + 2\gamma \kappa_\mathcal{L} \mathcal{W}, & \mu \geq 0 \\[2ex]
\dfrac{1}{k} \sum_{t=1}^k E_t^{tot} \leq \dfrac{2}{L\gamma k} (\mathbb{E}[\mathcal{L}(x_0) - \mathcal{L}^\star] + \gamma(E_0^y + E_0^z)) + \dfrac{2}{L} \mathcal{W}, & \mu < 0,
\end{cases}
$$

where $\mathcal{W}$ is defined in (22) of Appendix A.2 and satisfies:

$$
\mathcal{W} = O\left(\kappa_g^5 |\mathcal{D}_g|^{-1} \tilde{\sigma}_g^2 + \kappa_g^3 |\mathcal{D}_{g_{yy}}|^{-1} \tilde{\sigma}_{g_{yy}}^2 + \kappa_g^2 |\mathcal{D}_{g_{xy}}|^{-1} \tilde{\sigma}_{g_{xy}}^2 + \kappa_g^2 |\mathcal{D}_f|^{-1} \tilde{\sigma}_f^2 \right),
$$

with $\tilde{\sigma}_f^2$, $\tilde{\sigma}_g^2$, $\tilde{\sigma}_{g_{xy}}^2$ and $\tilde{\sigma}_{g_{yy}}^2$ being the variances defined in Assumptions 4 and 5.

We provide a proof outline in Appendix A describing the general approach and a detailed proof in Appendix A.2 which proceeds by combining the error bounds in (9), (12) and (19) to get a general error bound on E_k^{tot} in Proposition 9 expressed in terms of Λ_k and Π_k. Then, Proposition 10 of Appendix A.2 allows to deduce conditions on T and N so that the error bound on E_k^{tot} holds. The choice of a_k and b_k ensures that all terms defining E_k^{tot} are of the same order as k increases. Their dependence on Λ_k and Π_k ensures that the contributions of E_k^y and E_k^z to E_k^{tot} matter less as the algorithms $\mathcal{A}_k$ and $\mathcal{B}_k$ become more accurate. Next, we apply Theorem 1 to obtain the worst case complexity of AmIGO in various settings. First, we define the measure of complexity for AmIGO.

Measure of complexity. Let the number of iterations k, T and N and sizes of the batches $\mathcal{D}_g$, $\mathcal{D}_f$, $\mathcal{D}_{g_{xy}}$ and $\mathcal{D}_{g_{yy}}$, be such that AmIGO algorithm achieves a precision ϵ according to some pre-defined criterion. We define the *complexity* $\mathcal{C}(\epsilon)$ to be the total number of queries to the gradients of f and g, Jacobian and hessian vector products performed by the AmIGO algorithm

$$
\mathcal{C}(\epsilon) = k \left(T|\mathcal{D}_g| + N|\mathcal{D}_{g_{yy}}| + |\mathcal{D}_{g_{xy}}| + |\mathcal{D}_f| \right). \tag{13}
$$

We provide now the complexity of AmIGO in the deterministic setting when $\mathcal{L}$ is strongly convex.

Corrolary 1 (Case $\mu > 0$ and $\mathcal{W} = 0$). *Use batches of size 1. Achieving an error $\mathcal{L}(x_k) - \mathcal{L}^\star \leq \epsilon$ requires* $\mathcal{C}(\epsilon) = O\left(\kappa_\mathcal{L} \kappa_g \log \left(\frac{E_0^{tot}}{\epsilon} \right) \right)$.

This result improves over those in Table 1 in terms of the dependence on ϵ. It is possible to improve the dependence on κ_g to $\kappa_g^{1/2}$ by accelerating algorithms $\mathcal{A}_k$ and $\mathcal{B}_k$ as we illustrate in Section 5, or using generic acceleration methods such as Catalyst (Lin et al., 2018). The following result deals with the stochastic setting:

Corrolary 2 (Case $\mu > 0$ and $\mathcal{W} > 0$). *Choose* $|\mathcal{D}_f| = O\left(\frac{\kappa_\mathcal{L} \tilde{\sigma}_f^2}{\epsilon} \right)$, $|\mathcal{D}_g| = O\left(\frac{\kappa_\mathcal{L} \kappa_g^2 \tilde{\sigma}_g^2}{\epsilon} \right)$, $|\mathcal{D}_{g_{xy}}| = O\left(\frac{\kappa_\mathcal{L} \tilde{\sigma}_{g_{xy}}^2}{\epsilon} \right)$ *and* $|\mathcal{D}_{g_{yy}}| = O\left(\frac{\kappa_\mathcal{L} \tilde{\sigma}_{g_{yy}}^2}{\epsilon} \right)$. *Achieving an error* $\mathbb{E}[\mathcal{L}(\hat{x}_k) - \mathcal{L}^\star] \leq \epsilon$ *requires:*

$$
\mathcal{C}(\epsilon) = O\left(\kappa_\mathcal{L}^2 \left(\kappa_g^3 \tilde{\sigma}_g^2 + \kappa_g \tilde{\sigma}_{g_{yy}}^2 + \tilde{\sigma}_{g_{xy}}^2 + \tilde{\sigma}_f^2 \right) \frac{1}{\epsilon} \log \left(\frac{E_0^{tot} + \mathbb{E}[\mathcal{L}(x_0) - \mathcal{L}^\star]}{\epsilon} \right) \right).
$$

Corrolary 2 improves over those in Table 1 and recovers the dependence on ϵ up to a logarithmic factor of stochastic gradient descent for smooth and strongly convex functions. The following result provides the complexity of AmIGO in the deterministic setting when $\mathcal{L}$ is non-convex:

Corrolary 3 (Case $\mu < 0$ and $\mathcal{W} = 0$)**.** *Choose batches of size* 1. *Achieving an error* $\frac{1}{k}\sum_{i=1}^{k}\|\nabla\mathcal{L}(x_i)\|^2 \leq \epsilon$ *requires* $\mathcal{C}(\epsilon) = \mathcal{O}\big(\big(\kappa_g^4(\mathcal{L}(x_0) - \mathcal{L}^\star) + \kappa_g(E_0^y + E_0^z)\big)\frac{1}{\epsilon}\big)$.

Corrolary 3 recovers the complexity of AID-BiO (Ji et al., 2021) which is expected since their proposed algorithm also exploited warm-start for both $\mathcal{A}_k$ and $\mathcal{B}_k$. The following result provides the complexity in the stochastic setting for non-convex loss $\mathcal{L}$:

Corrolary 4 (Case $\mu < 0$ and $\mathcal{W} > 0$)**.** *Choose* $|\mathcal{D}_{g_{xy}}| = O\big(\frac{\kappa_g^2}{\epsilon}\tilde{\sigma}_{g_{xy}}^2\big)$, $|\mathcal{D}_{g_{yy}}| = O\big(\frac{\kappa_g^3}{\epsilon}\tilde{\sigma}_{g_{yy}}^2\big)$, $|\mathcal{D}_f| = O\big(\frac{\kappa_g^2\tilde{\sigma}_f^2}{\epsilon}\big)$ *and* $|\mathcal{D}_g| = O\big(\frac{\kappa_g^5\tilde{\sigma}_g^2}{\epsilon}\big)$. *Achieving an error* $\frac{1}{k}\sum_{i=1}^{k}\mathbb{E}\big[\|\nabla\mathcal{L}(x_i)\|^2\big] \leq \epsilon$ *requires:*

$$\mathcal{C}(\epsilon) = O\Bigg(\Big(\kappa_g^6\tilde{\sigma}_g^2 + \kappa_g^4\tilde{\sigma}_{g_{yy}}^2 + \kappa_g^2(\tilde{\sigma}_{g_{xy}}^2 + \tilde{\sigma}_f^2)\Big)\big(\kappa_g^3\mathbb{E}[\mathcal{L}(x_k) - \mathcal{L}^\star] + E_0^y + E_0^z\big)\frac{1}{\epsilon^2}\Bigg)$$

Corrolary 4 recovers the optimal dependence on ϵ of $O(\frac{1}{\epsilon^2})$ achieved by stochastic gradient descent in the smooth non-convex case (Arjevani et al., 2019, Theorem 1). It also improves over the results in (Ji et al., 2021) which involve an additional logarithmic factor $\log(\epsilon^{-1})$ since N is required to be $O(\kappa_g\log(\epsilon^{-1}))$. In our case, N remains constant since $\mathcal{B}_k$ benefits from warm-start. The faster rates of MRBO/VRBO* (Yang et al., 2021) are obtained under the additional *mean-squared smoothness* assumption (Arjevani et al., 2019), which we do not investigate in the present work. Such assumption allows to achieve the improved complexity of $O(\epsilon^{-3/2}\log(\epsilon^{-1}))$. However, these algorithm still require $N = O(\log(\epsilon^{-1}))$, indicating that the use of warm-start in $\mathcal{B}_k$ could further reduce the complexity to $O(\epsilon^{-3/2})$ which would be an interesting direction for future work.

5 Experiments

5.1 Synthetic problem

To study the behavior of AmIGO in a controlled setting, we consider a synthetic problem where both inner and outer level losses are quadratic functions, thus admitting closed form expression for the minimizers. More precisely, the functions f and g are chosen as $f(x,y) := \frac{1}{2}x^\top A_f x + y^\top C_f$ and $g(x,y) := \frac{1}{2}y^\top A_g y + y^\top B_g x$ where A_f and A_g are symmetric definite positive matrices of size $d_x \times d_x$ and $d_y \times d_y$, B_g is a $d_y \times d_x$ matrix and C_f is a d_y vector with $d_x = 2000$ and $d_y = 1000$.

We generate A_f and A_g so that the smoothness constants L and L_g are fixed to 1, $\kappa = 10$ and κ_g taking values in $\{10^i, i \in \{0,..,7\}\}$. We then solve each problem using 8 different methods. We first consider AmIGO with either gradient descent (AmIGO-GD) or conjugate gradient (AmIGO-CG) for algorithm $\mathcal{B}_k$. We then compare these algorithms with AID methods without warm-start for $\mathcal{B}_k$ which we refer to as (AID-GD) and (AID-CG). We also consider other variants using either a fixed-point algorithm (AID-FP) (Grazzi et al., 2020) or Neumann series expansion (AID-N) (Lorraine et al., 2020) for $\mathcal{B}_k$. Finally, we consider two algorithms based on iterative differentiation which we refer to as (ITD) (Grazzi et al., 2020) and (Reverse) (Franceschi et al., 2017). For all methods, we use warm-start in algorithm $\mathcal{A}_k$, however only AmIGO exploits warm-start in $\mathcal{B}_k$, the other AID based methods intializing $\mathcal{B}_k$ with $z^0 = 0$. We fix the step-sizes to $\gamma_k = 1/L$ and $\alpha_k = \beta_k = 1/L_g$ and perform a grid-search on the number of iterations T and M of algorithms $\mathcal{A}_k$ and $\mathcal{B}_k$ from $\in \{10^i, i \in 0, 1, 2, 3\}$. For AID methods without warm-start in $\mathcal{B}_k$, we consider an additional setting where M increases logarithmically with k, as suggested in Ji et al. (2021), with $M = \lfloor 10^3\log(k)\rfloor$. Similarly, for (ITD) and (Reverse), we additionally use an increasing T of the same form. Figure 1(a) shows the complexity $\mathcal{C}(\epsilon)$ needed to reach 10^{-6} relative error amongst the best choice for T and M as the conditioning number κ_g increases. AmIGO-CG achieves the lowest time and is followed by AID-CG thus showing a favorable effect of warm-start for $\mathcal{B}_k$. The same conclusion holds for AmIGO-GD compared to AID-GD. Note that AID-CG is still faster than AmIGO-CG for larger values of κ_g highlighting the advantage of using algorithms $\mathcal{B}_k$ with $O(\kappa_g^{-1/2})$ complexity such as (CG) instead non-accelerated ones with $O(\kappa_g^{-1})$ such as (GD). Figure 1(b) shows the relative error after 10s and maintains the same conclusions. For moderate values of κ_g, only AmIGO and AID-CG reach an error of 10^{-20} as shown in Figure 1(c). We refer to Figures 2 and 3 of Appendix F for additional results on the effect of the choice of T and M showing that AmIGO consistently performs well for a wide range of values of T and M.

Figure 1: Top row: performance on the synthetic task. The relative error is defined as a ratio between current and initial errors $(\mathcal{L}(x_k) - \mathcal{L}^\star)/(\mathcal{L}(x_0) - \mathcal{L}^\star)$. The complexity $\mathcal{C}(\epsilon)$ as defined in (13). Bottom row: performance on the hyper-parameter optimization task.

5.2 HYPER-PARAMETER OPTIMIZATION

To illustrate the effectiveness of the method on high dimensional problems, we consider a classification task on the `20Newsgroup` dataset using a linear model. The inner-level cost g is a regularized logistic loss over a training set with y being the regression matrix mapping from an input data space of dimension $d = 101631$ to a target space of dimension $p = 20$ and x is a vector of regularization parameters of dimension d. Each component of x regularizes a particular input dimension y. The outer-level cost f is an un-regularized regression loss over a validation set. We consider again both versions of AmIGO (GD and CG) and compare them to the methods discussed in Section 5.1. We also compare with BSA algorithm (Ghadimi and Wang, 2018), TTSA algorithm (Hong et al., 2020a) and stocBiO (Ji et al., 2021). We also compared with VRBO (Yang et al., 2021) using the implementation available online and noticed instabilities for large values of T and N, as reported by the authors, but also a drop in performance compared to stocBiO for smaller T and N due to inexact estimates of the gradient. We consider two evaluations settings: A *default* setting and a *grid-seach* search setting. Note that AmIGO relies on the same parameters as other AID methods. Thus, for the *default* setting, we use the well-chosen parameters reported in Grazzi et al. (2020); Ji et al. (2021) where $\alpha_k=\gamma_k=100$, $\beta_k=0.5$, and $T=N=10$. For the grid-search setting, we perform a grid-search near the default values of β_k, T and N as detailed in Appendix F.1. We also vary the batch-size from $10^3 * \{0.1, 1, 2, 4\}$ and report the best performing choices for each method. Figure 1(d,e,f) show AmIGO-CG to be the fastest, achieving the lowest error and highest validation and test accuracies. The test accuracy of AmIGO-CG decreases after exceeding 80% indicating a potential overfitting as also observed in Franceschi et al. (2018). Similarly, AmIGO-GD outperformed all other methods that uses an algorithm $\mathcal{B}_k$ with $O(\kappa_g)$ complexity. Moreover, all remaining methods achieved comparable performance matching those reported in Ji et al. (2021), thus indicating that the warm-start in $\mathcal{B}_k$ and acceleration in $\mathcal{B}_k$ were the determining factors for obtaining an improved performance. Additionally, Figure 4 of Appendix F report similar results for each choice of the batch-size and shows a behavior that is consistent with the observations of Figure 1(d,e,f).

6 CONCLUSION

We studied AmIGO, an algorithm for bilevel optimization based on amortized implicit differentiation and introduced a unified framework for analyzing its convergence. Our analysis showed that AmIGO achieves the same complexity as *unbiased oracle methods*, thus achieving improved rates compared to methods without warm-start in various settings. We then illustrated empirically such improved convergence in both synthetic and a hyper-optimization experiments. A future research direction consists in extending the proposed framework to non-smooth objectives and analyzing acceleration in both inner and outer level problems as well as variance reduction techniques.

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A   CONVERGENCE OF AMIGO ALGORITHM

In this section, we provide a proof of Theorem 1 as well as its corollaries Corrolaries 1 to 4. In Appendix A.2 we provide an outline of the proof Theorem 1 that states the main intermediary results needed for the proof and provides explicit expressions for the quantities needed throughout the rest of the paper. The proofs of the intermediary results are deferred to Appendices B and C.

A.1   PROOF OUTLINE OF THEOREM 1

The proof of Theorem 1 proceeds in 8 steps as discussed bellow.

Step 1: Smoothness of properties This step consists in characterizing the smoothness of $\nabla \mathcal{L}$, $y^\star$, $z^\star$ as well as the conditional expectation $\mathbb{E}[\hat{\psi}_k | x_k, y_k, z_k]$ knowing x_k, y_k and z_k. For this purpose, we consider the function $\Psi : \mathcal{X} \times \mathcal{Y} \times \mathcal{Y} \to \mathcal{X}$ defined as follows:

$$\Psi(x, y, z) := \partial_x f(x, y) + \partial_{xy} g(x, y) z.$$

Hence, by definition of $\hat{\psi}_k$, it is easy to see that $\mathbb{E}[\hat{\psi}_k | x_k, y_k, z_k] = \Psi(x_k, y_k, z_k)$. The following proposition controls the smoothness of $\nabla \mathcal{L}$, $y^\star$, $z^\star$ and Ψ and is adapted from (Ghadimi and Wang, 2018, Lemma 2.2). We provide a proof in Appendix B.2 for completeness.

Proposition 6. *Under Assumptions 1 to 3, $\mathcal{L}$, Ψ, $y^\star$ and $z^\star$ satisfy:*

$$\|y^\star(x) - y^\star(x')\| \le L_y \|x - x'\|, \qquad \|z^\star(x, y) - z^\star(x', y')\| \le L_z(\|x - x'\| + \|y - y'\|)$$
$$\|\nabla \mathcal{L}(x) - \nabla \mathcal{L}(x')\| \le L \|x - x'\|, \quad \|\Psi(x, y, z) - \nabla \mathcal{L}(x)\| \le L_\psi [\|y - y^\star(x)\| + \|z - z^\star(x, y)\|]$$
$$\|z^\star(x, y^\star(x))\| \le \mu_g^{-1} B,$$

where L_y, L_z, L_ψ and L are given by:

$$L_y := \mu_g^{-1} L_g', \qquad L_z := \mu_g^{-2} M_g B + \mu_g^{-1} L_f, \tag{14}$$
$$L_\psi := \max\left((L_f + M_g \mu_g^{-1} B + L_g' L_z), L_g'\right)$$
$$L := \left(L_f + \mu_g^{-2} L_g' M_g B + \mu_g^{-1}\left(L_g' L_f + M_g B\right)\right)\left(1 + \mu_g^{-1} L_g'\right)$$

The expressions of L_y, L_z, L_ψ and L suggests the following dependence on the conditioning κ_g of the inner-level problem which will be useful for the complexity analysis: $L_y = O(\kappa_g)$, $L_\psi = O(\kappa_g^2)$, $L_z = O(\kappa_g^2)$ and $L = O(\kappa_g^3)$.

Step 2: Convergence of the inner-level iterates. In this step, we control the mean squared errors $\mathbb{E}[\|y_k - y^\star(x_k)\|]$ and $\mathbb{E}[\|z_k - z^\star(x_k, y_k)\|]$ as stated in Proposition 3. In fact we prove a slightly stronger version stated below:

Proposition 7. *Let α_k and β_k be two positive sequences with $\alpha_k \le L_g^{-1}$ and $\beta_k \le L_g^{-1}\left(1 + \mu_g^{-2} \sigma_{g_{yy}}^2\right)$ and define $\Lambda_k := (1 - \alpha_k \mu_g)^T$ and $\Pi_k := \left(1 - \frac{\beta_k \mu_g}{2}\right)^N$. Denote by $\bar{z}_k$ the conditional expectation of z_k knowing x_k, y_k and z_k^0. Let R_k^y and R_k^z be defined as:*

$$R_k^y := 2\alpha_k^2 \sigma_g^2 \left(T \wedge \frac{1}{\alpha_k \mu_g}\right), \quad R_k^z := \beta_k^2 \left(B^2 \mu_g^{-2} \sigma_{g_{yy}}^2 + 3\left(N \wedge \frac{1}{\beta_k \mu_g}\right)\sigma_f^2\right)\left(N \wedge \frac{1}{\beta_k \mu_g}\right). \tag{15}$$

Then, under Assumptions 1, 4 and 5 the iterates z_k and $\bar{z}_k$ satisfy:

$$\mathbb{E}\left[\|y_k - y^\star(x_k)\|^2\right] \le \Lambda_k E_k^y + R_k^y$$
$$\mathbb{E}\left[\|z_k - z^\star(x_k, y_k)\|^2\right] \le \Pi_k \mathbb{E}\left[\|z_k^0 - z^\star(x_k, y_k)\|^2\right] + R_k^z,$$
$$\mathbb{E}\left[\|\bar{z}_k - z^\star(x_k, y_k)\|^2\right] \le \Pi_k \mathbb{E}\left[\|z_k^0 - z^\star(x_k, y_k)\|^2\right]$$
$$\mathbb{E}\left[\|z_k - \bar{z}_k\|^2\right] \le 4\sigma_g^2 \mu_g^{-2} \Pi_k \mathbb{E}\left[\|z_k^0 - z^\star(x_k, y_k)\|^2\right] + 2R_k^z. \tag{16a}$$

It is easy to see from the above expressions that $R_k^y = O(\kappa_g \sigma_g^2)$ while $R_k^z = O\left(\kappa_g^3 \sigma_{g_y y}^2 + \kappa_g^2 \sigma_f^2\right)$ as stated in Proposition 3. While controlling $\mathbb{E}[\|y_k - y^\star(x_k)\|]$ follows by standard results on SGD (Kulunchakov and Mairal, 2020, Corollary 31), since the iterates of Algorithm 2 uses i.i.d. samples, the error terms $\mathbb{E}[\|z_k - z^\star(x_k, y_k)\|]$ is more delicate since Algorithm 3 uses the same sample $\partial_y \hat{f}(x_k, y_k)$ for updating the iterates, therefore introducing additional correlations between these iterates. We defer the proof of Proposition 7 to Appendix B.3 which relies on a general result for stochastic linear systems with correlated noise that we provide in Appendix E.

Step 3: Controlling the bias and variance errors V_k^ψ and E_k^ψ is achieved by Proposition 5. The bias E_k^ψ is controlled simply by using the smoothness of the potential Ψ near the point $(x_k, y^\star(x_k), z^\star(x_k, y^\star(x_k)))$ as shown in Proposition 6. The variance term $V_k^\star$ is more delicate to control due to the multiplicative noise resulting from the Jacobian-vector product between $\partial_{xy} \hat{g}(x_k, y_k, \mathcal{D}_{g_{xy}}) z_k$. We defer the proof to Appendix B.4 and provide below explicit expressions for the constants σ_x^2 and w_x^2:

$$w_x^2 := \left(1 + 2L_g' \mu_g^{-1} + 6\left(\sigma_{g_{xy}}^2 + (L_g')^2\right) \mu_g^{-2}\right) \sigma_f^2$$
$$+ 2\left(\sigma_{g_{xy}}^2 + (L_g')^2\right) B^2 L_g^{-1} \mu_g^{-3} \sigma_{g_{yy}}^2 + 2B^2 \mu_g^{-2} \sigma_{g_{xy}}^2.$$
$$\sigma_x^2 := 2\sigma_{g_{xy}}^2 + 2(L_g')^2 \mu_g^{-2} \sigma_{g_{yy}}^2.$$

Note that $w_x^2 = O\left(\kappa_g^2\left(\sigma_f^2 + \sigma_{g_{xy}}^2\right) + \kappa_g^3 \sigma_{g_{yy}}^2\right)$ and $\sigma_x^2 = O\left(\sigma_{g_{xy}}^2 + \kappa_g^2 \sigma_{g_{yy}}^2\right)$, as stated in Proposition 5.

Step 4: Outer-level error bound. This step consists in obtaining the inequality in Proposition 2 which extends the result of (Kulunchakov and Mairal, 2020, Proposition 1) to biased gradients and to the non-convex case. We defer the proof of such result to Appendix C.1.

Step 5: Inner-level error bound. This step consists in proving Proposition 4. For clarity, we provide a second statement with explicit expressions for the quantities of interest:

Proposition 8. *Let r_k and θ_k be two positive non-increasing sequences no greater than 1. For any $0 \leq v \leq 1$, denote by ϕ_k and $\tilde{R}_k$ the following non-negative scalars:*

$$\phi_k := (1-v) L_g^2 T^2 \alpha_k^2 + 2v, \qquad \tilde{R}_k^y := 2(1-v) L_g \mu_g^{-1} T^2 \sigma_g^2 \alpha_k^2 + v R_k^y$$
$$\zeta_k := 2L\left(\min\left((1-u)^{-1}, L\eta_k^{-1}\right) \mathbb{1}_{\mu \geq 0} + \mathbb{1}_{\mu < 0}\right)$$

Finally, consider the following matrices and vectors:

$$\boldsymbol{P_k} := \begin{pmatrix} 1 + r_k, & 0 \\ 16 L_z^2 \frac{\phi_k}{\theta_k}, & 1 + \theta_k \end{pmatrix}, \quad \boldsymbol{U_k} := \begin{pmatrix} 2L_y^2 \frac{\gamma_k^2}{r_k} \\ 4L_z^2 \frac{\gamma_k^2}{\theta_k}\left(1 + 4L_y^2 \frac{\phi_k}{r_k}\right) \end{pmatrix}, \quad \boldsymbol{V_k} := \tilde{R}_k^y \begin{pmatrix} 0 \\ 8L_z^2 \theta_k^{-1} \end{pmatrix}.$$

$$(18)$$

Then, under Assumptions 1 to 5, the following holds:

$$\begin{pmatrix} E_k^y \\ E_k^z \end{pmatrix} \leq \boldsymbol{P_k} \begin{pmatrix} \Lambda_{k-1} E_{k-1}^y + R_{k-1}^y \\ \Pi_{k-1} E_{k-1}^z + R_{k-1}^z \end{pmatrix} + 2\left(E_{k-1}^\psi + V_{k-1}^\psi + \zeta_k E_{k-1}^x\right) \boldsymbol{U_k} + \boldsymbol{V_k}. \qquad (19)$$

We defer the proof of the above result to Appendix C.2.

Step 6: General error bound. By combining the inequalities in Propositions 2 and 8 resulting from the analysis of both outer and inner levels, we obtain a general error bound on E_k^{tot} that holds without any condition on the error made by Algorithms 2 and 3:

Proposition 9. *Choose the step-sizes α_k and β_k such that they are non-increasing in k and choose r_k and θ_k such that $\gamma_k r_k^{-1}$ and $\gamma_k \theta_k^{-1}$ are non-increasing sequences. Choose the coefficients a_k and b_k defining E_k^{tot} in (7) to be of the form $a_k = \gamma_k r_k^{-1} \Lambda_k^s$ and $b_k = \gamma_k \theta_k^{-1} \Pi_k^s$ for some $0 < s < 1$. Fix a scalar $0 < \kappa < \frac{1}{2}$ and a non-increasing sequence $0 < \rho_k < 1$. Then, under Assumptions 1 to 5 E_k^{tot} satisfies:*

$$F_k + E_k^{tot} \leq \|A_k\|_\infty E_{k-1}^{tot} + V_k^{tot}.$$

where V_k and $\boldsymbol{A_k}$ are given by:

$$A_k := \begin{pmatrix} 1 - (1-\kappa\rho_k)\delta_k + 2\zeta_k\lambda_k u_k^I \\ \Lambda_k^{1-s}\left(1 + r_k\left(1 + 16L_z^2\Pi_k^s\theta_k^{-2}\phi_k + 2L_\psi^2\left(2u_k^I + (s_k + (2\kappa\rho_k)^{-1})\right)\right)\right) \\ \Pi_k^{1-s}\left(1 + \theta_k\left(1 + \left(2L_\psi^2 + \sigma_x^2\right)\left(2u_k^I + (s_k + (2\kappa\rho_k)^{-1})\right)\right)\right) \end{pmatrix} \tag{20}$$

$$V_k^{tot} := \gamma_k\left(\Lambda_k^s(1 + r_k^{-1}) + 16L_z^2\phi_k\theta_k^{-2}\Pi_k^s + 2L_\psi^2\left(2u_k^I + (s_k + (2\kappa\rho_k)^{-1})\right)\right)R_{k-1}^y$$
$$+ \gamma_k\left((1 + \theta_k^{-1})\Pi_k^s R_{k-1}^z + 8L_z^2\theta_k^{-2}\Pi_k^s\tilde{R}_k^y\right) + \gamma_k\left(s_k + u_k^I\right)w_x^2$$

where we introduced $u_k^I = a_k U_k^{(1)} + b_k U_k^{(2)}$ for conciseness with $U_k^{(1)}$ and $U_k^{(2)}$ being the components of the vector $\boldsymbol{U_k}$ defined in Proposition 8.

The general form of a_k and b_k allows to account for potentially decreasing step-sizes γ_k, α_k and β_k. Hence, the result of Proposition 9 could be used to analyze algorithms with decreasing step-sizes. In the present work we will restrict to the constant step-size for ease of presentation as we discuss next.

Step 7: Controlling the precision of the inner-level algorithms. In this step, we provide conditions on T and N to that $\|A\|_\infty \leq 1 - (1-\rho_k)\delta_k$ in the constant step-size case. Under such condition, we also obtain a simpler upper-bound on the variance V_k^{tot}. For this purpose, we consider the following constants which will appear in the conditions on T and N:

$$A_1 := 2\log(2 + 2L_\psi^2(\kappa\rho_0)^{-1} + 2(1 + 6L_\psi^2)) \tag{21a}$$

$$A_2 := 2\log\left(2 + 32^2 L^{-2}\rho_0^{-1}(1-\kappa)^{-1}\tilde{\zeta}\right) \tag{21b}$$

$$A_1' := 2\log(2 + 2(L_\psi^2 + 2^{-1}\sigma_x^2)(\kappa\rho_0)^{-1} + 2(1 + 6L_\psi^2)) \tag{21c}$$

$$A_2' := 2\log\left(1 + L_z^2 L_y^{-2}\left(1 + 16L_y^2\right)\right), \tag{21d}$$

where $\tilde{\zeta}$ is given by $\tilde{\zeta} := \mathbb{1}_{\mu<0} + \mathbb{1}_{\mu>0}\mu^{-1}L\min((1-u)^{-1}, \mu^{-1}L)$. We then have the following proposition, which we prove in Appendix D:

Proposition 10. *Let Assumptions 1 to 5 hold. Choose $\rho_k = \rho_0$ for some $0 < \rho_0 < 1$, the step-sizes to be constant: $\alpha_k = \alpha \leq \frac{1}{L_g}$, $\beta_k = \beta \leq \frac{1}{L_g\left(1 + \mu_g^{-1}\sigma_{g_{yy}}^2\right)}$ and $\gamma_k = \gamma \leq \frac{1}{L}$. Moreover, set $s = \frac{1}{2}$, $r_k = \theta_k = 1$. Finally, choose T and N as follows:*

$$T = \lfloor\alpha^{-1}\mu_g^{-1}\max(1, A_1, A_2)\rfloor + 1, \quad N = \lfloor2\beta^{-1}\mu_g^{-1}(\max(1, A_1, A_2) + \max(1, A_1', A_2'))\rfloor + 1$$

with A_1, A_2, A_1' and A_2' defined in (21a) to (21d). Then, $\|A\|_\infty \leq 1 - (1-\rho_k)\delta_k$ and $V_k^{tot} \leq \gamma\mathcal{W}$, with $\mathcal{W}$ given by:

$$\mathcal{W} := \gamma\left(\frac{1-u}{2} + \left(\frac{1+u}{2}\delta_0 + 4L_y^2 L^{-2}\right)\right)w_x^2 + 2\gamma\mu_g^{-3}\left(B^2\beta\sigma_{g_{yy}}^2 + 3\mu_g\sigma_f^2\right) \tag{22}$$
$$+ 2\gamma\alpha\mu_g^{-1}\left(3 + 2L_\psi^2\left(3 + (2\kappa\rho_0)^{-1}\right) + 8L_z^2\right)\sigma_g^2$$

It is easy to see from Proposition 10 that when $\alpha = \frac{1}{L_g}$ and $\beta = \frac{1}{L_g(1 + \mu_g^{-2}\sigma_{g_{yy}}^2)}$, we obtain that $T = O(\kappa_g)$ and $N = O\left(\kappa_g\left(1 + \kappa_g^2\sigma_{g_{yy}}^2\right)\right)$, where the big-$O$ notation hides a logarithmic dependence in κ_g coming from the constants A_1, A_2, A_1' and A_2'.

Step 8: Proving the main inequalities. The final step combines Propositions 9 and 10 to get the desired inequality. We provide a full proof in Appendix A.2 assuming Propositions 9 and 10 hold.

A.2 Proof of Theorem 1

In order to prove Appendix A.2 in the convex case, we need the following *averaging strategy lemma*, a generalization of (Kulunchakov and Mairal, 2020, Lemma 30):

Lemma 1. *Let $\mathcal{L}$ be a convex function on $\mathcal{X}$. Let x_k be a (potentially stochastic) sequence of iterates in $\mathcal{X}$. Let $(E_k)_{k\geq 0}$, $(V_k)_{k\geq 0}$ and $(\delta_k)_{k\geq 0}$ be non-negative sequences such that $\delta_k \in (0, 1)$. Fix*

some non-negative number $u \in [0,1]$ and define the following averaged iterates $\hat{x}_k$ recursively by $\hat{x}_k = u(1-\delta_k)\hat{x}_{k-1}+(1-(1-\delta_k)u)x_k$ and starting from any initial point $\hat{x}_0$. Assume the iterates $(x_k)_{k \geq 1}$ satisfy the following relation for all $k \geq 1$:

$$(1-u(1-\delta_k))\mathbb{E}[\mathcal{L}(x_k)-\mathcal{L}^\star]+E_k \leq (1-\delta_k)(E_{k-1}+(1-u)\mathbb{E}[\mathcal{L}(x_{k-1})-\mathcal{L}^\star])+V_k. \quad (23)$$

Let $\Gamma_k := \prod_{t=1}^{k}(1-\delta_k)$. Then the averaged iterates $(\hat{x}_k)_{k \geq 1}$ satisfy the following:

$$\mathbb{E}[\mathcal{L}(\hat{x}_k)-\mathcal{L}^\star]+E_k \leq \Gamma_k\big(E_0+\mathbb{E}\big[\mathcal{L}(x_0)-\mathcal{L}^\star+u^k(\mathcal{L}(\hat{x}_0)-\mathcal{L}^\star)\big]\big)+\Gamma_k \sum_{1 \leq t \leq k}\Gamma_t^{-1}V_t.$$

Proof. For simplicity, we write $F_k = \mathbb{E}[\mathcal{L}(x_k)-\mathcal{L}^\star]$ and $\hat{F}_k = \mathbb{E}[\mathcal{L}(\hat{x}_k)-\mathcal{L}^\star]$. We first multiply (23) by Γ_k^{-1} and sum the resulting inequalities for all $1 \leq k \leq K$ to get:

$$\sum_{k=1}^{K}\Gamma_k^{-1}(1-u(1-\delta_k))F_k+\Gamma_k^{-1}T_k \leq \sum_{k=1}^{K}\Gamma_k^{-1}(1-\delta_k)(T_{k-1}+(1-u)F_{k-1})+\sum_{k=1}^{K}\Gamma_k^{-1}V_k.$$

Grouping the terms in F_k together and recalling that $\Gamma_k^{-1}(1-\delta_k) = \Gamma_{k-1}^{-1}$ yields:

$$\sum_{k=1}^{K}\Gamma_k^{-1}F_k-\Gamma_{k-1}^{-1}F_{k-1}+u\sum_{k=1}^{K}\Gamma_{k-1}^{-1}(F_{k-1}-F_k) \leq \sum_{k=1}^{K}\left(\Gamma_{k-1}^{-1}T_{k-1}-\Gamma_k^{-1}T_k\right)+\sum_{k=1}^{K}\Gamma_k^{-1}V_k.$$

Simplifying the telescoping sums and multiplying by Γ_k, we get:

$$F_K+u\Gamma_K\sum_{k=1}^{K}\Gamma_{k-1}^{-1}(F_{k-1}-F_k) \leq -T_K+\Gamma_K\left(F_0+T_0+\sum_{k=1}^{K}\Gamma_k^{-1}V_k\right). \quad (24)$$

Consider now the quantity $\hat{F}_k$. Recalling that $\mathcal{L}$ is convex and by definition of the iterates $\hat{x}_k$, by Jensen's inequality, it holds that:

$$\hat{F}_k \leq u(1-\delta_k)\hat{F}_{k-1}+(1-u(1-\delta_k))F_k.$$

By iteratively applying the above inequality, we get that:

$$\hat{F}_K \leq u^K\Gamma_K\hat{F}_0+\Gamma_K\sum_{k=1}^{K}u^{K-k}\Gamma_k^{-1}(1-u(1-\delta_k))F_k$$

$$=u^K\Gamma_K\hat{F}_0+\Gamma_K\sum_{k=1}^{K}u^{K-k}\Gamma_k^{-1}F_k-\sum_{k=1}^{K}u^{K-k+1}\Gamma_{k-1}F_k$$

$$=u^K\Gamma_K\hat{F}_0+F_K+\Gamma_K\sum_{k=1}^{K-1}u^{K-k}\Gamma_k^{-1}(F_k-F_{k+1})$$

We can therefore apply (24) to the above inequality to get the desired result. $\square$

We can now proceed to prove Theorem 1.

Proof of Theorem 1. By application of Proposition 9 and using the choice of T and N given by Proposition 10, the following inequality holds:

$$F_k+E_k^{tot} \leq (1-(1-\rho_k)\delta_k)E_{k-1}^{tot}+\gamma\mathcal{W}. \quad (25)$$

with $\mathcal{W}$ defined in Proposition 10. We then distinguish two cases depending on the sign of μ:

Case $\mu \geq 0$. In this case, we recall that F_k and E_k^x are given by:

$$F_k = u\delta_k\mathbb{E}[\mathcal{L}(x_k)-\mathcal{L}^\star], \qquad E_k^x = \frac{\delta_k}{2\gamma_k}\|x_k-x^\star\|^2+(1-u)\mathbb{E}[\mathcal{L}(x_k)-\mathcal{L}^\star]$$

Since $\mu > 0$, $\mathcal{L}$ is a convex function and we can apply Lemma 1 with $V_k=\gamma\mathcal{W}$ and $E_k=\frac{\delta_k}{2\gamma_k}\|x_k-x^\star\|^2+a_kE_k^y+b_kE_k^z$. The result follows by noting that the sum $\Gamma_k\sum_{t=1}^{k}\Gamma_t^{-1} \leq \kappa_\mathcal{L}$.

Case $\mu < 0$. In this case, we recall that F_k and E_k^x are given by:

$$F_k = \mathbb{E}[\mathcal{L}(x_k) - \mathcal{L}(x_{k-1})] + E_{k-1}^x - E_k^x, \qquad E_k^x = L^{-1}\mathbb{E}\Big[\|\nabla\mathcal{L}(x_k)\|^2\Big].$$

We then sum (25) for all iterations $0 \le t \le k$ which, by telescoping, yields:

$$\mathbb{E}[\mathcal{L}(x_k) - \mathcal{L}(x_0)] + E_0^x - E_k^x + E_k^{tot} - E_0^{tot} + \sum_{1 \le t \le k}(1-\rho_t)\delta_t E_t^{tot} \le k\gamma\mathcal{W}.$$

Using that $\mathbb{E}[\mathcal{L}(x_k) - \mathcal{L}^\star] + E_k^{tot} - E_k^x$ is non-negative since $E_k^{tot} - E_k^x = a_k E_k^y + b_k E_k^z$, we get:

$$\sum_{1 \le t \le k}(1-\rho_t)\delta_t E_t^x \le (\mathbb{E}[\mathcal{L}(x_0) - \mathcal{L}^\star] + a_0 E_0^y + b_0 E_0^z) + k\gamma\mathcal{W}.$$

Finally, since $\rho_t = \frac{1}{2}$, $\delta_t = L\gamma$, the result follows after dividing both sides by $\frac{1}{2}kL\gamma$. $\qquad\square$

A.3 PROOF OF CORROLARIES 1 TO 4

Proof of Corrolary 1. We choose $u = 0$ which implies that $E_k^x = \frac{\mu}{2}\|x_k - x^\star\|^2 + \mathcal{L}(x_k) - \mathcal{L}^\star \le E_k^{tot}$. We can then apply Theorem 1 for $\mu > 0$ which, after noting that $\mathcal{W} = 0$, yields the following:

$$\frac{\mu}{2}\|x_k - x^\star\| + \mathcal{L}(x_k) - \mathcal{L}^\star \le \left(1 - (2\kappa_{\mathcal{L}})^{-1}\right)^k E_0^{tot}.$$

Hence, to achieve an error $\mathcal{L}(x_k) - \mathcal{L}^\star \le \epsilon$ for some $\epsilon > 0$, we need $k = O\Big(\kappa_{\mathcal{L}}\log\Big(\frac{E_0^{tot}}{\epsilon}\Big)\Big)$. Since $T = O(\kappa_g)$ and $N = O(\kappa_g(1 + \kappa_g^2\sigma_{g_{yy}}^2))$ as required by Theorem 1. Moreover, since $\sigma_{g_{yy}}^2 = 0$ in the deterministic setting, it holds that $N = O(\kappa_g)$. Using batches of size 1, yields the desired complexity. $\qquad\square$

Proof of Corrolary 2. Here we choose $u = 1$ and apply Theorem 1 for $\mu > 0$ which yields:

$$\mathbb{E}[\mathcal{L}(\hat{x}_k) - \mathcal{L}^\star] + E_k^{tot} \le \left(1 - (2\kappa_{\mathcal{L}})^{-1}\right)^k \left(E_0^{tot} + \mathbb{E}[\mathcal{L}(x_0) - \mathcal{L}^\star]\right) + 2\kappa_{\mathcal{L}}L^{-1}\mathcal{W}$$

Hence, to achieve an error $\mathbb{E}[\mathcal{L}(\hat{x}_k) - \mathcal{L}^\star] \le \epsilon$, we need $k = O\Big(\kappa_{\mathcal{L}}\log\Big(\frac{E_0^{tot}+\mathbb{E}[\mathcal{L}(x_0)-\mathcal{L}^\star]}{\epsilon}\Big)\Big)$ to guarantee that the first term in the l.h.s. of the above inequality is $O(\epsilon)$. Moreover, we recall that $L^{-1} = O(\kappa_g^{-3})$ from Proposition 6 and that Theorem 1 ensure the variance $\mathcal{W}$ satisfies:

$$\mathcal{W} = O\Big(\kappa_g^5\sigma_g^2 + \kappa_g^3\sigma_{g_{yy}}^2 + \kappa_g^2\sigma_{g_{xy}}^2 + \kappa_g^2\sigma_f^2\Big).$$

Hence, ensuring the second term $2L^{-1}\kappa_{\mathcal{L}}\mathcal{W} = O(\epsilon)$ is achieved by choosing the size of the batches as follows:

$$|\mathcal{D}_f| = O\Big(\frac{\kappa_{\mathcal{L}}}{\epsilon}\Big), \quad |\mathcal{D}_g| = O\Big(\frac{\kappa_{\mathcal{L}}\kappa_g^2}{\epsilon}\Big), \quad |\mathcal{D}_{g_{xy}}| = O\Big(\frac{\kappa_{\mathcal{L}}}{\epsilon}\Big), \quad |\mathcal{D}_{g_{yy}}| = O\Big(\frac{\kappa_{\mathcal{L}}}{\epsilon}\Big)$$

Recall that $T = O(\kappa_g)$ and $N = O\Big(\kappa_g(1 + \kappa_g^2\sigma_{g_{yy}}^2)\Big)$ as required by, Theorem 1. Moreover, since $|\mathcal{D}_{g_{yy}}| = O(\frac{\kappa_{\mathcal{L}}}{\epsilon})$, it follows that $1 + \kappa_g^2\sigma_{g_{yy}}^2 = O(1)$ and hence that $N = O(\kappa_g)$, which yields the desired result. $\qquad\square$

Proof of Corrolary 3. In the non-convex deterministic case, recall that $E_k^x = \frac{1}{L}\|\nabla\mathcal{L}(x_k)\|^2 \le E_k^{tot}$. We thus apply Theorem 1 for $\mu < 0$, multiply by L and recall that $\mathcal{W} = 0$ in the deterministic to get:

$$\frac{1}{k}\sum_{t=1}^k \|\nabla\mathcal{L}(x_t)\|^2 \le \frac{2L}{k}\big(\mathcal{L}(x_0) - \mathcal{L}^\star + L^{-1}(E_0^y + E_0^z)\big)$$

Therefore, recalling that $L = O(\kappa_g^3)$ from Proposition 6, to achieve an error of order $\min 1 \le t \le k\|\nabla\mathcal{L}(x_t)\|^2 \le \epsilon$, it suffice to choose $k = O\Big(\frac{(\kappa_g^3\mathcal{L}(x_0)-\mathcal{L}^\star+(E_0^y+E_0^z))}{\epsilon}\Big)$. Thus using batches of size 1 and T and N of order κ_g (since $\sigma_{g_{yy}}^2 = 0$) we directly get the desired complexity. $\qquad\square$

Proof of Corrolary 4. In the non-convex stochastic case, $E_k^x = \frac{1}{L}\mathbb{E}\Big[\|\nabla\mathcal{L}(x_k)\|^2\Big] \leq E_k^{tot}$. We thus apply Theorem 1 for $\mu < 0$, multiply by L to get:

$$\frac{1}{k}\sum_{t=1}^{k}\mathbb{E}\Big[\|\nabla\mathcal{L}(x_t)\|^2\Big] \leq \frac{2L}{k}\big(\mathbb{E}[\mathcal{L}(x_0) - \mathcal{L}^\star] + L^{-1}(E_0^y + E_0^z)\big) + 2\mathcal{W}.$$

to achieve an error of order $\min 1 \leq t \leq k\mathbb{E}\Big[\|\nabla\mathcal{L}(x_t)\|^2\Big] \leq \epsilon$, we need to ensure each term in the l.h.s. of the above inequality is of order ϵ. For the first term, similarly to the deterministic setting Corrolary 3, we simply need $k = O\bigg(\frac{(\kappa_g^3\mathcal{L}(x_0) - \mathcal{L}^\star + (E_0^y + E_0^z))}{\epsilon}\bigg)$. For the second term, we need to have $\mathcal{W} + O(\epsilon)$, which is achieved using the following choice for the sizes of the batches:

$$|\mathcal{D}_f| = O\left(\frac{\kappa_g^2}{\epsilon}\right), \quad |\mathcal{D}_g| = O\left(\frac{\kappa_g^5}{\epsilon}\right), \quad |\mathcal{D}_{g_{xy}}| = O\left(\frac{\kappa_g^2}{\epsilon}\right), \quad |\mathcal{D}_{g_{yy}}| = O\left(\frac{\kappa_g^3}{\epsilon}\right)$$

Finally, as required by Theorem 1, we set $T = O(\kappa_g)$ and $N = O(\kappa_g(1 + \kappa_g\sigma_{g_{yy}}^2))$ and recall that $1 + \kappa_g\sigma_{g_{yy}}^2 = O(1)$ since $\sigma_{g_{yy}}^2 = O(|\mathcal{D}_{g_{yy}}|^{-1}) = O(\epsilon\kappa_g^{-3})$, thus yielding the desired complexity. $\square$

B PRELIMINARY RESULTS

B.1 EXPRESSION OF THE GRADIENT

We provide a proof of Proposition 1 which shows that $\mathcal{L}$ is differentiable and provides an expression of its gradient.

Proof. Assumption 1 ensures that $y \mapsto g(x, y)$ admits a unique minimizer $y^\star(x)$ defined as the unique solution to the implicit equation $\partial_y g(x, y^\star(x)) = 0$. Moreover, since g is twice continuously differentiable and strongly convex, it follows that $\partial_{yy}g(x, y^\star(x))$ is invertible for any $x \in \mathcal{X}$. Therefore the implicit function theorem (Lang, 2012, Theorem 5.9), ensures that $x \mapsto y^\star(x)$ is continuously differentiable with Jacobian given by $\nabla y^\star(x) = -\partial_{xy}g(x, y^\star(x))\partial_{yy}g(x, y^\star(x))^{-1}$. Hence, by composition of differentiable functions, $\mathcal{L}$ is also differentiable with gradient given by:

$$\nabla\mathcal{L}(x) = \partial_x f(x, y^\star(x)) - \partial_{xy}g(x, y^\star(x))\partial_{yy}g(x, y^\star(x))^{-1}\partial_y f(x, y^\star(x)).$$

We can thus define $z^\star(x, y) = -\partial_{yy}g(x, y)^{-1}\partial_y f(x, y)$ to get the desired expression for $\nabla\mathcal{L}(x)$ and note that $z^\star$ is the solution to (2). $\square$

B.2 SMOOTHNESS PROPERTIES OF $\mathcal{L}$, $\mathcal{Y}^\star$, $z^\star$ AND Ψ

Proof of Proposition 6. **Lipschitz continuity of $x \mapsto y^\star(x)$.** By Assumptions 1 and 3, the implicit function theorem (Lang, 2012, Theorem 5.9) ensures $y^\star(x)$ is differentiable with Jacobian given by:

$$\nabla y^\star(x) = -\partial_{xy}g(x, y^\star(x))(\partial_{yy}g(x, y^\star(x)))^{-1}.$$

Moreover, by Assumption 3, we know that $\partial_y g(x, y)$ is L_g'-Lipschitz in x for any $y \in \mathcal{Y}$, hence, $\|\partial_{xy}g(x, y^\star(x))\|_{op}$ is upper-bounded by L_g'. Moreover, by Assumption 1, g is μ_g-strongly convex in y uniformly on $\mathcal{X}$. Therefore, it holds that $\big\|\partial_{yy}g(x, y^\star(x))^{-1}\big\|_{op} \leq \mu_g^{-1}$. This allows to deduce that $\|\nabla y^\star(x)\|_{op} \leq \mu_g^{-1}L_g'$, any by application of the fundamental theorem of calculus that:

$$\|y^\star(x) - y^\star(x')\| \leq \mu_g^{-1}L_g'\|x - x'\|.$$

This shows that $y^\star$ is L_y-Lipschitz continuous with $L_y := \mu_g^{-1}L_g'$.

Lipchitz continuity of $x \mapsto z^\star(x, y)$. Let (x, y) and (x', y') be two points in $\mathcal{X} \times \mathcal{Y}$. Recalling the definition of $z^\star(x, y)$ in Proposition 1, it is easy to see that $z^\star(x, y)$ admits the following expression:

$$z^\star(x, y) = -\partial_{yy}g(x, y)^{-1}\partial_y f(x, y). \tag{26}$$

Recalling the expression of $z^\star(x, y)$, the following holds:

$$
\begin{aligned}
z^\star(x, y) - z^\star(x', y') =&\partial_{yy}g(x', y')^{-1}\partial_y f(x', y') - \partial_{yy}g(x, y)^{-1}\partial_y f(x, y) \\
=&\left(\partial_{yy}g(x', y')^{-1} - \partial_{yy}g(x, y)^{-1}\right)\partial_y f(x', y') \\
&+ \partial_{yy}g(x, y)^{-1}(\partial_y f(x', y') - \partial_y f(x, y)) \\
=&\partial_{yy}g(x', y')^{-1}(\partial_{yy}g(x, y) - \partial_{yy}g(x', y'))\partial_{yy}g(x, y)^{-1}\partial_y f(x', y') \\
&+ \partial_{yy}g(x, y)^{-1}(\partial_y f(x', y') - \partial_y f(x, y))
\end{aligned}
$$

Hence, by taking the norm of the above quantity a triangular inequality followed by operator inequalities, it follows that:

$$
\begin{aligned}
\|z^\star(x, y) - z^\star(x', y')\| \leq& \left\|H_2^{-1}\right\|_{op}\|H_1 - H_2\|_{op}\left\|H_1^{-1}\right\|_{op}\|\partial_y f(x', y')\| \\
&+ \left\|H_1^{-1}\right\|\|\partial_y f(x', y') - \partial_y f(x, y)\|.
\end{aligned}
$$

where we introduced $H_1 = \partial_{yy}g(x, y)$ and $H_2 = \partial_{yy}g(x', y')$ for conciseness. Using Assumption 1, we can upper-bound $\left\|H_1^{-1}\right\|_{op}$ and $\left\|H_2^{-1}\right\|_{op}$ by μ_g^{-1}. By Assumption 3, we know that $\|H_1 - H_2\|_{op} \leq M_g\|(x, y) - (x', y')\|$. Finally by Assumption 2, we also have that $\|\partial_y f(x', y') - \partial_y f(x, y)\| \leq L_f\|(x, y) - (x', y')\|$ and that $\|\partial_y f(x', y')\| \leq B$ ensuring that:

$$
\|z^\star(x, y) - z^\star(x', y')\| \leq \left(\mu_g^{-2}M_g B + \mu_g^{-1}L_f\right)\|(x, y) - (x', y')\|.
$$

Hence, we conclude that $z^\star$ is L_z-Lipchitz continuous with L_z defined as in (14).

boundedness of $z^\star(x, y)$ Recalling the expression of $z^\star$ in (26), it is easy to see that $\|z^\star(x, y)\|$ is upper-bounded by $\mu_g^{-1}B$ since $\partial_{yy}g(x, y)$ is μ_g-strongly convex in y by Assumption 1 and $\partial_y f(x, y)$ is bounded by B by Assumption 2.

Regularity of Ψ.

$$
\begin{aligned}
\Psi(x, y, z) - \Psi(x', y', z') =&\partial_x f(x, y) - \partial_x f(x', y') + \partial_{xy}g(x, y)z - \partial_{xy}g(x', y')z' \\
=&\partial_x f(x, y) - \partial_x f(x', y') + \partial_{xy}g(x, y)(z - z') \\
&+ (\partial_{xy}g(x, y) - \partial_{xy}g(x', y'))z'.
\end{aligned}
$$

By taking the norm of the above expression and applying a triangular inequality followed by operator inequalities, it follows that:

$$
\begin{aligned}
\|\Psi(x, y, z) - \Psi(x', y', z')\| \leq& \|\partial_x f(x, y) - \partial_x f(x', y')\| + \|\partial_{xy}g(x, y)\|_{op}\|z - z'\| \quad (27) \\
&+ \|\partial_{xy}g(x, y) - \partial_{xy}g(x', y')\|_{op}\|z'\| \\
\leq& L_f(\|x - x'\| + \|y - y'\|) + L_g'\|z - z'\| \\
&+ M_g\|z'\|(\|x - x'\| + \|y - y'\|).
\end{aligned}
$$

To get the first term of the last inequality above, we used that $\partial_x f$ is L_f-Lipschitz by Assumption 2. To get the second term, we used that $\partial_{xy}g(x, y)$ is bounded since $\partial_y g(x, y)$ is L_g'-Lipschitz by Assumption 3. Finally, for the last term, we used that $\partial_{xy}g(x, y)$ is M_g-Lipschitz by Assumption 3.

By choosing $x' = x$, $y' = y^\star(x)$ and $z' = z^\star(x, y^\star(x))$, it is easy to see from Proposition 1 that $\Psi(x, y^\star(x), z^\star(x, y^\star(x))) = \nabla\mathcal{L}(x)$. Hence, applying the above inequality yields:

$$
\begin{aligned}
\|\Psi(x, y, z) - \nabla\mathcal{L}(x)\| \leq& (L_f + M_g\|z^\star(x, y^\star(x))\|)\|y - y^\star(x)\| + L_g'\|z - z^\star(x, y^\star(x))\| \\
\leq& (L_f + M_g\|z^\star(x, y^\star(x))\|)\|y - y^\star(x)\| + L_g'\|z - z^\star(x, y)\| \\
&+ L_g'\|z^\star(x, y) - z^\star(x, y^\star(x))\|
\end{aligned}
$$

As shown earlier, $\|z^\star(x, y^\star(x))\|$ is upper-bounded by $\mu_g^{-1}B$, while $\|z^\star(x, y) - z^\star(x, y^\star(x))\|$ is bounded by $L_z\|y - y^\star(x)\|$. This allows to conclude that $\|\Psi(x, y, z) - \nabla\mathcal{L}(x)\| \leq L_\psi$ with L_ψ defined in (14).

Lipschitz continuity of $x \mapsto \nabla\mathcal{L}(x)$. We apply (27) with $(y, z) = (y^\star(x), z^\star(x, y^\star(x)))$ and $(y', z') = (y^\star(x'), z^\star(x, y^\star(x')))$ which yields:

$$
\begin{aligned}
\|\nabla\mathcal{L}(x) - \nabla\mathcal{L}(x')\| \leq& (L_f + M_g\|z^\star(x', y^\star(x'))\|)(\|x - x'\| + \|y^\star(x) - y^\star(x')\|) \\
&+ L_g'\|z^\star(x, y^\star(x)) - z^\star(x', y^\star(x'))\| \\
\leq& \left(L_f + M_g\mu_g^{-1}B + L_g'L_z\right)(1 + L_y)\|x - x'\|,
\end{aligned}
$$

where we used that $\|z^\star(x', y^\star(x'))\|$ is upper-bounded by $\mu_g^{-1}B$, $z^\star$ is L_z-Lipschitz and $y^\star$ is L_y-Lipschitz. Hence, $\nabla\mathcal{L}$ is L-Lipschitz continuous, with L as given by (14). $\qquad\square$

B.3 Convergence of the iterates of algorithms $\mathcal{A}_k$ and $\mathcal{B}_k$

Proof. **Controlling the iterates y^t of $\mathcal{A}_k$.**

Consider a new batch $\mathcal{D}_g$ of samples ξ. We have by definition of the update equation of y^t that:

$$
\begin{aligned}
\left\|y^t - y^\star(x_k)\right\|^2 =& \left\|y^{t-1} - y^\star(x_k)\right\|^2 + \alpha_k^2\left\|\partial_y\hat{g}\big(x_k, y^{t-1}, \mathcal{D}_g\big)\right\|^2 \\
& - 2\alpha_k\partial_y\hat{g}\big(x_k, y^{t-1}, \mathcal{D}_g\big)^\top\big(y^{t-1} - y^\star(x_k)\big)
\end{aligned}
$$

Taking the expectation conditionally on x_k and y^{t-1}, we get:

$$
\begin{aligned}
\mathbb{E}\Big[\left\|y^t - y^\star(x_k)\right\|^2\Big|x_k, y^{t-1}\Big] =& (1 - \alpha_k\mu_g)\left\|y^{t-1} - y^\star(x_k)\right\|^2 \\
& + \alpha_k^2\mathbb{E}\Big[\left\|\partial_y\hat{g}\big(x_k, y^{t-1}, \mathcal{D}_g\big) - \partial_y g\big(x_k, y^{t-1}\big)\right\|^2\Big|x_k, y^{t-1}\Big] \\
& - 2\alpha_k\partial_y g\big(x_k, y^{t-1}\big)^\top\Big(y^{t-1} - y^\star(x_k) - \frac{\alpha_k}{2}\partial_y g\big(x_k, y^{t-1}\big)\Big) \\
\leq& (1 - \alpha_k\mu_g)\left\|y^{t-1} - y^\star(x_k)\right\|^2 + 2\alpha_k^2\sigma_g^2
\end{aligned}
$$

The second equality uses that $\partial_y\hat{g}\big(x_k, y^{t-1}, \mathcal{D}_g\big)$ is an unbiased estimator of $\partial_y g\big(x_k, y^{t-1}\big)$. For the second inequality, we use Assumption 4 which allows to upper-bound the variance of $\partial_y\hat{g}$ by σ_g^2. Moreover, since g is convex and L_g-smooth and since $\alpha_k \leq L_g^{-1}$, it follows that the last term in the above inequality is non-positive and can thus be upper-bounded by 0. By unrolling the resulting inequality recursively for $1 < t \leq k$, we obtain the desired result.

Controlling the iterates z^n of $\mathcal{B}_k$. The poof follows by direct application of Proposition 14 with $\beta = \beta_k$ and the following choices for $A_n,$, A, $\hat{b}$, b:

$$
\begin{aligned}
A_n =& \partial_{yy}\hat{g}(x_k, y_k, \mathcal{D}_{g_{yy}}), & \hat{b} =& \partial_y\hat{f}(x_k, y_k, \mathcal{D}_f) \\
A =& \partial_{yy}g(x_k, y_k) & b =& \partial_y f(x_k, y_k).
\end{aligned}
$$

This directly yields the following inequalities:

$$
\begin{aligned}
\mathbb{E}\Big[\|z_k - z^\star(x_k, y_k)\|^2\Big] \leq& \tilde{\Pi}_k\mathbb{E}\Big[\left\|z_k^0 - z^\star(x_k, y_k)\right\|^2\Big] + \tilde{R}_k^z, \\
\mathbb{E}\Big[\|\bar{z}_k - z^\star(x_k, y_k)\|^2\Big] \leq& \tilde{\Pi}_k\mathbb{E}\Big[\left\|z_k^0 - z^\star(x_k, y_k)\right\|^2\Big].
\end{aligned}
$$

where $\tilde{\Pi}_k$ and $\tilde{R}_k^z$ are given by:

$$
\tilde{\Pi}_k := (1 - \beta_k\mu_g)^N, \qquad \tilde{R}_k^z := \beta_k^2\bigg(\sigma_g^2\mathbb{E}\Big[\|z^\star(x_k, y_k)\|^2\Big] + 3\bigg(N \wedge \frac{1}{\beta_k\mu_g}\bigg)\sigma_f^2\bigg)\bigg(N \wedge \frac{1}{\beta_k\mu_g}\bigg).
$$

First we have that $\tilde{\Pi}_k \leq \Pi_k$. Moreover, Proposition 6, we have that $\|z^\star(x_k, y_k)\| \leq B\mu_g^{-1}$ hence, $\tilde{R}_k^z \leq R_k^z$ thus yielding the desired inequalities. Finally (16a) also follows similarly using (45) from Proposition 14. $\qquad\square$

B.4 Controlling the bias and variance E_k^ψ and V_k^ψ

Proof of Proposition 5 . Recall that the expressions of E_k^ψ and V_k^ψ in (6) involves the conditional expectation $\mathbb{E}_k\Big[\hat{\psi}_k\Big]$ knowing x_k, y_k and z_{k-1}. This can be also expressed using Ψ as follows:

$$
\begin{aligned}
\mathbb{E}_k\Big[\hat{\psi}_k\Big] =& \mathbb{E}\Big[\mathbb{E}\Big[\hat{\psi}_k|x_k, y_k, z_k\Big]|x_k, y_k, z_{k-1}\Big] \\
=& \mathbb{E}[\Psi(x_k, y_k, z_k)|x_k, y_k, z_{k-1}] \\
=& \mathbb{E}[\Psi(x_k, y_k, \mathbb{E}[z_k||x_k, y_k, z_{k-1}])]
\end{aligned}
$$

where we used the tower property for conditional expectations in the first line, then the fact that the expectation of ψ_k conditionally on x_k, y_k and z_k is simply $\Psi(x_k, y_k, z_k)$. Finally, for the last line, we use the independence of the noise and the linearity of Ψ w.r.t. the last variable. In all what follows, we write $\bar{z}_k = \mathbb{E}[z_k | x_k, y_k, z_{k-1}]$ which is the same object as defined in Proposition 7. We then treat E_k^ψ and V_k^ψ separately.

Bounding E_k^ψ. Using Propositions 6 and 7 we directly get the desired inequality:

$$
\begin{aligned}
E_k^\psi \leq & 2L_\psi^2 \Big(\mathbb{E}\big[\|y_k - y^\star(x_k)\|^2 \big] + \mathbb{E}\big[\|z_k - z^\star(x_k, y_k)\|^2 \big] \Big) \\
\leq & 2L_\psi^2 (\Lambda_k E_k^y + \Pi_k E_k^z + R_k^y)
\end{aligned}
$$

Bound on V_k^ψ. We decompose V_k^ψ into a sum of three terms W_k, W_k' and W_k'' given by:

$$
W_k := \mathbb{E}\Big[\big\| \partial_x \hat{f}(x_k, y_k, \mathcal{D}_f) - \partial_x f(x_k, y_k) \big\|^2 \Big],
$$

$$
W_k' := \mathbb{E}\Big[\big\| \partial_{xy}\hat{g}\big(x_k, y_k, \tilde{\xi}_{N+1,k}\big) z_k - \partial_{xy} g(x_k, y_k) \bar{z}_k \big\|^2 \Big]
$$

$$
W_k'' := \mathbb{E}\Big[\big(\partial_x \hat{f}(x_k, y_k, \mathcal{D}_f) - \partial_x f(x_k, y_k) \big)^\top \partial_{xy} g(x_k, y_k)(z_k - \bar{z}_k) \Big].
$$

where we used that $\tilde{\xi}_{N+1,k}$ is independent from z_k and $\mathcal{D}_f$ to get the last term. Hence, using Assumption 4 to bound the first term of the above relation, we get $V_k^\psi \leq \sigma_f^2 + W_k' + 2W_k''$. Thus, it remains to control each of W_k' and W_k''.

Bound on W_k''. Using that $\mathcal{D}_f$ is independent from $\tilde{\xi}_{n,k}$, we can apply Proposition 13 to write:

$$
W_k'' = \beta_k \mathbb{E}\Big[\partial_x(\hat{f}-f)(x_k, y_k, \mathcal{D}_f)^\top \partial_{xy} g(x_k, y_k) \Big(\sum_{t=1}^N (I - \beta_k A)^{N-t} \Big) \partial_y(\hat{f}-f)(x_k, y_k, \mathcal{D}_f) \Big]
$$

where we used the simplifying notion $(\hat{f}-f)(x_k, y_k, \mathcal{D}_f) = \hat{f}(x_k, y_k, \mathcal{D}_f) - f(x_k, y_k)$. Using Assumption 3 to bound $\|\partial_{xy} g(x_k, y_k)\|_{op}$ by L_g', Assumption 1 to upper-bound $\big\| \big(\sum_{t=1}^N (I - \beta_k A)^{N-t} \big) \big\|_{op}$ by $\big(\sum_{t=1}^N (1 - \beta_k \mu_g)^{N-t} \big)$ we get

$$
\begin{aligned}
|W_k''| \leq & \beta_k L_g' \sum_{t=0}^{N-1} (1 - \beta_k \mu_g)^t \mathbb{E}\Big[\big| \partial_x(\hat{f}-f)(x_k, y_k, \mathcal{D}_f)^\top \partial_y(\hat{f}-f)(x_k, y_k, \mathcal{D}_f) \big| \Big] \\
\leq & L_g' \mu_g^{-1} \mathbb{E}\Big[\big| \partial_x(\hat{f}-f)(x_k, y_k, \mathcal{D}_f)^\top \partial_y(\hat{f}-f)(x_k, y_k, \mathcal{D}_f) \big| \Big] \\
\leq & L_g' \mu_g^{-1} \mathbb{E}\Big[\big\| \partial_x(\hat{f}-f)(x_k, y_k, \mathcal{D}_f)^2 \big\| \Big]^{\frac{1}{2}} \mathbb{E}\Big[\big\| \partial_y(\hat{f}-f)(x_k, y_k, \mathcal{D}_f) \big\|^2 \Big]^{\frac{1}{2}} \leq \quad L_g' \mu_g^{-1} \sigma_f^2
\end{aligned}
$$

where we used that $\sum_{t=0}^{N-1}(1 - \beta_k \mu_g) \leq \frac{1}{\beta_k \mu_g}$ for the second line, Cauchy-Schwarz inequality to get the third line and Assumption 4 to get the last line.

bound on W'_k Using that $\tilde{\xi}_{N+1,k}$ is independent from z_k, we write:

$$
\begin{aligned}
W'_k =& \mathbb{E}\left[\left\|\partial_{xy}(\hat{g}-g)\Big(x_k,y_k,\tilde{\xi}_{N+1,k}\Big)z_k\right\|^2\right] + \mathbb{E}\left[\|\partial_{xy}g(x_k,y_k)(z_k-\bar{z}_k)\|^2\right] \\
&\overset{(i)}{\leq} \sigma^2_{g_{xy}}\mathbb{E}\left[\|z_k\|^2\right] + (L'_g)^2\mathbb{E}\left[\|z_k-\bar{z}_k\|^2\right] \\
&\overset{(ii)}{\leq} 2\sigma^2_{g_{xy}}\left(\mathbb{E}\left[\|z_k-z^\star(x_k,y_k)\|^2\right] + \mathbb{E}\left[\|z^\star(x_k,y_k)\|^2\right]\right) + (L'_g)^2\mathbb{E}\left[\|z_k-\bar{z}_k\|^2\right] \\
&\overset{(iii)}{\leq} 2\sigma^2_{g_{xy}}\mathbb{E}\left[\|z_k-z^\star(x_k,y_k)\|^2\right] + (L'_g)^2\mathbb{E}\left[\|z_k-\bar{z}_k\|^2\right] + 2\sigma^2_{g_{xy}}B^2\mu_g^{-2} \\
&\overset{(iv)}{\leq} 2\sigma^2_{g_{xy}}\left(\Pi_k\mathbb{E}\left[\|z_k^0-z^\star(x_k,y_k)\|^2 + R_k^z\right]\right) \\
&\quad + (L'_g)^2\left(4\sigma^2_{g_{yy}}\mu_g^{-2}\Pi_k\mathbb{E}\left[\|z_k^0-z^\star(x_k,y_k)\|^2\right] + 2R_k^z\right) + 2\sigma^2_{g_{xy}}B^2\mu_g^{-2} \\
&\leq 2\left(\sigma^2_{g_{xy}} + 2(L'_g)^2\mu_g^{-2}\sigma^2_{g_{yy}}\right)\Pi_k\mathbb{E}\left[\|z_k^0-z^\star(x_k,y_k)\|^2\right] \\
&\quad + 2\left(\sigma^2_{g_{xy}} + (L'_g)^2\right)R_k^z + 2\sigma^2_{g_{xy}}B^2\mu_g^{-2}
\end{aligned}
$$

(i) follows from Assumptions 3 and 5, (ii) uses that $\|z_k\|^2 \leq 2\left(\|z_k-z\|^2 + \|z\|^2\right)$, (iii) uses that $\|z^\star(x_k,y_k)\| \leq B\mu_g^{-1}$ by Proposition 6. Finally (iv) follows by application of Proposition 7. We further have by definition of R_k^z that:

$$
R_k^z \leq B^2 L_g^{-1}\mu_g^{-3}\sigma^2_{g_{yy}} + 3\mu_g^{-2}\sigma^2_f \tag{30}
$$

Combining the inequalities on W'_k, W''_k and (30), we get that $V_k^\psi \leq w_x^2 + \sigma_x^2\Pi_k E_k^z$, with w_x^2 and σ_x^2 given by (17). □

C GENERAL ANALYSIS OF AMIGO

C.1 ANALYSIS OF THE OUTER-LOOP

Proof of Proposition 2. We treat both cases $\mu \geq 0$ and $\mu < 0$ separately. For simplicity we denote by $\mathbb{E}_k$ the conditional expectation knowing the iterates x_k, y_k and z_{k-1} and write $\psi_k = \mathbb{E}_k\left[\hat{\psi}_k\right]$.

Case $\mu \geq 0$. Recall that E_k^x is given by:

$$
E_k^x = \frac{\eta_k}{2}\mathbb{E}\left[\|x_k-x^\star\|^2\right] + (1-u)\mathbb{E}[\mathcal{L}(x_k)-\mathcal{L}^\star].
$$

For simplicity define $\epsilon_k = u\delta_k + (1-u)$, $e_k = (1-u)(\mathcal{L}(x_k)-\mathcal{L}^\star) + \frac{\eta_k}{2}\|x_k-x^\star\|^2$ and $e'_k = u\delta_k(\mathcal{L}(x_k)-\mathcal{L}^\star)$. It is then easy to see that $\mathbb{E}[e_k]$ is equal to the l.h.s of (9), i.e. $\mathbb{E}[e_k] = E_k^x$. We

will start by bounding the difference between two successive iterates of e_k:

$$
\begin{aligned}
e'_k + e_k - e_{k-1} \leq\; & u\delta_k(\mathcal{L}(x_k) - \mathcal{L}^\star) + (1-u)(\mathcal{L}(x_k) - \mathcal{L}(x_{k-1})) \\
& + \frac{\eta_k}{2}\|x_k - x^\star\|^2 - \frac{\eta_{k-1}}{2}\|x_{k-1} - x^\star\|^2 \\
\overset{(i)}{\leq}\; & u\delta_k(\mathcal{L}(x_k) - \mathcal{L}^\star) + (1-u)(\mathcal{L}(x_k) - \mathcal{L}(x_{k-1})) \\
& - \frac{\delta_k \eta_{k-1}}{2}\|x_{k-1} - x^\star\|^2 + \delta_k\left(\frac{\mu}{2}\|x_{k-1} - x^\star\|^2 - \nabla\mathcal{L}(x_{k-1})^\top(x_{k-1} - x^\star)\right) \\
& + \frac{\eta_k}{2}\left(\gamma_k^2\|\hat{\psi}_{k-1}\|^2 - 2\gamma_k\left(\hat{\psi}_{k-1} - \nabla\mathcal{L}(x_{k-1})\right)^\top(x_{k-1} - x^\star)\right) \\
\overset{(ii)}{\leq}\; & u\delta_k(\mathcal{L}(x_k) - \mathcal{L}^\star) + (1-u)(\mathcal{L}(x_k) - \mathcal{L}(x_{k-1})) \\
& - \frac{\delta_k \eta_{k-1}}{2}\|x_{k-1} - x^\star\|^2 - \delta_k(\mathcal{L}(x_{k-1}) - \mathcal{L}^\star) \\
& + \frac{\eta_k}{2}\left(\gamma_k^2\|\hat{\psi}_{k-1}\|^2 - 2\gamma_k\left(\hat{\psi}_{k-1} - \nabla\mathcal{L}(x_{k-1})\right)^\top(x_{k-1} - x^\star)\right) \\
\leq\; & (u\delta_k + (1-u))(\mathcal{L}(x_k) - \mathcal{L}(x_{k-1})) \\
& - \frac{\delta_k \eta_{k-1}}{2}\|x_{k-1} - x^\star\|^2 - \delta_k(1-u)(\mathcal{L}(x_{k-1}) - \mathcal{L}^\star) \\
& + \frac{\eta_k}{2}\left(\gamma_k^2\|\hat{\psi}_{k-1}\|^2 - 2\gamma_k\left(\hat{\psi}_{k-1} - \nabla\mathcal{L}(x_{k-1})\right)^\top(x_{k-1} - x^\star)\right) \\
\leq\; & \epsilon_k(\mathcal{L}(x_k) - \mathcal{L}(x_{k-1})) - \delta_k e_{k-1} \\
& + \frac{\eta_k}{2}\left(\gamma_k^2\|\hat{\psi}_{k-1}\|^2 - 2\gamma_k\left(\hat{\psi}_{k-1} - \nabla\mathcal{L}(x_{k-1})\right)^\top(x_{k-1} - x^\star)\right) \\
\overset{(iii)}{\leq}\; & -\delta_k e_{k-1} + \epsilon_k \nabla\mathcal{L}(x_{k-1})^\top(x_k - x_{k-1}) + \frac{\epsilon_k L}{2}\|x_k - x_{k-1}\|^2 \\
& + \frac{\eta_k}{2}\left(\gamma_k^2\|\hat{\psi}_{k-1}\|^2 - 2\gamma_k\left(\hat{\psi}_{k-1} - \nabla\mathcal{L}(x_{k-1})\right)^\top(x_{k-1} - x^\star)\right) \\
=\; & -\delta_k e_{k-1} - \gamma_k \epsilon_k \nabla\mathcal{L}(x_{k-1})^\top\hat{\psi}_{k-1} + \frac{\epsilon_k \gamma_k^2 L}{2}\|\hat{\psi}_{k-1}\|^2 \\
& + \frac{\eta_k}{2}\left(\gamma_k^2\|\hat{\psi}_{k-1}\|^2 - 2\gamma_k\left(\hat{\psi}_{k-1} - \nabla\mathcal{L}(x_{k-1})\right)^\top(x_{k-1} - x^\star)\right).
\end{aligned}
$$

(i) follows from the update expression $x_k = x_{k-1} - \gamma_k\hat{\psi}_{k-1}$, (ii) follows from the convexity of $\mathcal{L}$ and (iii) follows by L-smoothness of $\mathcal{L}$. Taking the expectation conditionally on the randomness at iteration $k-1$ and using that $\mathbb{E}_{k-1}\left[\hat{\psi}_{k-1}\right] = \psi_{k-1}$, we therefore get

$$
\begin{aligned}
\mathbb{E}_{k-1}[e'_k + e_k - e_{k-1}] \leq\; & -\delta_k e_{k-1} - \gamma_k \epsilon_k \nabla\mathcal{L}(x_{k-1})^\top\psi_{k-1} + \frac{\gamma_k}{2}(\delta_k + \epsilon_k)\mathbb{E}_{k-1}\left[\|\hat{\psi}_{k-1}\|^2\right] \\
& - \delta_k(\psi_{k-1} - \nabla\mathcal{L}(x_{k-1}))^\top(x_{k-1} - x^\star) \\
=\; & -\delta_k e_{k-1} + \gamma_k s_k\left(\mathbb{E}_{k-1}\left[\|\hat{\psi}_{k-1} - \psi_{k-1}\|^2\right] + \|\psi_{k-1} - \nabla\mathcal{L}(x_{k-1})\|^2\right) \\
& - \delta_k(\psi_{k-1} - \nabla\mathcal{L}(x_{k-1}))^\top(x_{k-1} - \gamma_k\nabla\mathcal{L}(x_{k-1}) - x^\star) \\
& - \frac{\gamma_k}{2}(\epsilon_k - \delta_k)\|\nabla\mathcal{L}(x_{k-1})\|^2 \\
\overset{(i)}{\leq}\; & -\delta_k e_{k-1} + \gamma_k s_k\left(\mathbb{E}_{k-1}\left[\|\hat{\psi}_{k-1} - \psi_{k-1}\|^2\right] + \|\psi_{k-1} - \nabla\mathcal{L}(x_{k-1})\|^2\right) \\
& - \delta_k(\psi_{k-1} - \nabla\mathcal{L}(x_{k-1}))^\top(x_{k-1} - \gamma_k\nabla\mathcal{L}(x_{k-1}) - x^\star)
\end{aligned}
$$

where (i) follows from $\delta_k \leq \epsilon_k$ since by construction $\delta_k \leq 1$. Taking the expectation w.r.t. all the randomness yields the following inequality:

$$u\delta_k(\mathcal{L}(x_k) - \mathcal{L}^\star) + E_k^x \leq (1 - \delta_k)E_{k-1}^x + \gamma_k s_k \left(V_{k-1}^\psi + E_{k-1}^\psi\right) \tag{31}$$
$$- \delta_k \left(E_{k-1}^\psi\right)^{\frac{1}{2}} \mathbb{E}\left[\|x_{k-1} - \gamma_k \nabla \mathcal{L}(x_{k-1}) - x^\star\|^2\right]^{\frac{1}{2}}.$$

Since $\mathcal{L}$ is convex, we have the inequality: $\|x_{k-1} - \gamma_k \nabla \mathcal{L}(x_{k-1}) - x^\star\|^2 \leq \|x_{k-1} - x^\star\|^2$. Hence, we can deduce that:

$$\delta_k \|x_{k-1} - \gamma_k \nabla \mathcal{L}(x_{k-1}) - x^\star\|^2 \leq \delta_k \|x_{k-1} - x^\star\|^2 \leq 2\gamma_k \eta_k \eta_{k-1}^{-1} E_{k-1}^x \leq 2\gamma_k E_{k-1}^x,$$

where we used that η_k is non-increasing by construction. Combining the above inequality with (31) yields the desired result.

Case $\mu < 0$. Recall that for $\mu < 0$, we set $E_k^x = \frac{1}{L}\mathbb{E}\left[\|\nabla \mathcal{L}(x_k)\|^2\right]$. Using that $\mathcal{L}$ is L-smooth, we have that:

$$\mathcal{L}(x_k) - \mathcal{L}(x_{k-1}) \leq \nabla \mathcal{L}(x_{k-1})^\top (x_k - x_{k-1}) + \frac{L}{2}\|x_k - x_{k-1}\|^2$$
$$\leq -\gamma_k \nabla \mathcal{L}(x_{k-1})^\top \hat{\psi}_{k-1} + \frac{L\gamma_k^2}{2}\left\|\hat{\psi}_{k-1}\right\|^2$$
$$\leq -\gamma_k \|\nabla \mathcal{L}(x_{k-1})\|^2 - \gamma_k \nabla \mathcal{L}(x_{k-1})^\top \left(\hat{\psi}_{k-1} - \nabla \mathcal{L}(x_{k-1})\right)$$
$$+ \frac{L\gamma_k^2}{2}\left(\left\|\hat{\psi}_{k-1} - \psi_{k-1}\right\|^2 + 2\left(\hat{\psi}_{k-1} - \psi_{k-1}\right)^\top \psi_{k-1} + \|\psi_{k-1}\|^2\right).$$

Taking the expectation w.r.t. all randomness in the algorithm in the above inequality, we get:

$$\mathbb{E}[\mathcal{L}(x_k) - \mathcal{L}(x_{k-1})] \leq -\gamma_k \mathbb{E}\left[\|\nabla \mathcal{L}(x_{k-1})\|^2\right] - \gamma_k \mathbb{E}\left[\nabla \mathcal{L}(x_{k-1})^\top (\psi_{k-1} - \nabla \mathcal{L}(x_{k-1}))\right]$$
$$+ \frac{L\gamma_k^2}{2}\left(\mathbb{E}\left[\left(\left\|\hat{\psi}_{k-1} - \psi_{k-1}\right\|^2\right)\right] + \mathbb{E}\left[\|\psi_{k-1}\|^2\right]\right)$$
$$= -\gamma_k(1 - \frac{L\gamma_k}{2})\mathbb{E}\left[\|\nabla \mathcal{L}(x_{k-1})\|^2\right] + \frac{L\gamma_k^2}{2}\left(V_{k-1}^\psi + E_{k-1}^\psi\right)$$
$$- \gamma_k(1 - L\gamma_k)\mathbb{E}\left[\nabla \mathcal{L}(x_{k-1})^\top (\psi_{k-1} - \nabla \mathcal{L}(x_{k-1}))\right]$$
$$\leq -\frac{\gamma_k}{2}\mathbb{E}\left[\|\nabla \mathcal{L}(x_{k-1})\|^2\right] + \frac{L\gamma_k^2}{2}\left(V_{k-1}^\psi + E_{k-1}^\psi\right)$$
$$+ \gamma_k \mathbb{E}\left[\|\nabla \mathcal{L}(x_{k-1})\|^2\right]^{\frac{1}{2}}\left(E_{k-1}^\psi\right)^{\frac{1}{2}}.$$
$$= -\delta_k E_{k-1}^x + \frac{\delta_k \gamma_k}{2}\left(V_{k-1}^\psi + E_{k-1}^\psi\right) + \sqrt{2}\delta_k^{\frac{1}{2}}\gamma_k^{\frac{1}{2}}\left(E_{k-1}^x\right)^{\frac{1}{2}}\left(E_{k-1}^\psi\right)^{\frac{1}{2}}.$$

where we used that $1 - \frac{L\gamma_k}{2} \geq \frac{1}{2}$ and $0 \leq 1 - L\gamma_k \leq 1$ to get the last inequality. We get the desired result using Young's inequality, to upper-bound the last term in the r.h.s. of the above inequality. More precisely, we use that for any $0 < \rho_k < 1$:

$$\sqrt{2}\gamma_k^{\frac{1}{2}}\delta_k^{\frac{1}{2}}\left(E_{k-1}^\psi\right)^{\frac{1}{2}}\left(E_{k-1}^x\right)^{\frac{1}{2}} \leq \rho_k \delta_k E_{k-1}^x + (2\rho_k)^{-1}\gamma_k E_{k-1}^\psi.$$

$\square$

C.2 INNER-LEVEL ERROR BOUND

In this section we prove Proposition 4 which controls the evolutions of the warm-start errors E_k^y and E_k^z. As a first step, in Proposition 11, we provide a result controlling the mean squared error between two successive iterates x_{k-1}, x_k and y_{k-1}, y_k which will be used in the proof of Proposition 4.

Proposition 11 (Control of the increments of x_k and y_k). *Consider ζ_k, ϕ_k and $\tilde{R}_k^y$ as defined in Proposition 8 for some fixed $0 \le v \le 1$. Then, the following holds:*

$$\mathbb{E}\Big[\|x_k - x_{k-1}\|^2\Big] \le \gamma_k^2\Big(V_{k-1}^\psi + 2E_{k-1}^\psi + 2\zeta_k E_{k-1}^x\Big)$$

$$\mathbb{E}\Big[\|y_k - y_{k-1}\|^2\Big] \le 2\phi_k E_k^y + 2\tilde{R}_k^y$$

Proof. Proof of Proposition 11 We prove each inequality separately.

Increments of x_k. By the update equation, we have that $x_k = x_{k-1} - \gamma_k \hat{\psi}_k$, hence we only need to control $\mathbb{E}\Big[\big\|\hat{\psi}_k\big\|^2\Big]$. We have the following:

$$\mathbb{E}\Big[\big\|\hat{\psi}_k\big\|^2\Big] \le \mathbb{E}\Big[\big\|\hat{\psi}_k - \psi_k\big\|^2\Big] + 2\mathbb{E}\Big[\|\psi_k - \nabla\mathcal{L}(x_k)\|^2\Big] + 2\mathbb{E}\Big[\|\nabla\mathcal{L}(x_k)\|^2\Big]$$
$$= V_k^\psi + 2E_k^\psi + 2\mathbb{E}\Big[\|\nabla\mathcal{L}(x_k)\|^2\Big].$$

In the case $(\mu < 0)$, we have $E_k^x = \frac{1}{2L}\mathbb{E}\Big[\|\nabla\mathcal{L}(x_k)\|^2\Big]$, hence by setting $\zeta_k := 2L$, we get the desired inequality. In the convex case $(\mu \ge 0)$, since $\mathcal{L}$ is L-smooth, we have that:

$$\|\nabla\mathcal{L}(x_k)\|^2 \le 2L(\mathcal{L}(x_k) - \mathcal{L}^\star) \le 2L(1-u)^{-1}E_k^x,$$

provided that $u < 1$. We also have that $(\mathcal{L}(x_k) - \mathcal{L}^\star) \le \frac{L}{2}\|x_k - x^\star\|^2 \le L\eta_k^{-1}E_k^x$ which yields $\|\nabla\mathcal{L}(x_k)\|^2 \le 2L^2\eta_k^{-1}E_k^x$. Hence, we can set $\zeta_k = 2L\min\big((1-u)^{-1}, L\eta_k^{-1}\big)$.

Increments of y_k. Denoting by $\mathcal{D}_g^t$ a batch of samples at time iteration t of algorithm $\mathcal{A}_k$ and using the update equation of y^t we get the following inequality by application of the triangular inequality:

$$\mathbb{E}\Big[\|y_k - y_{k-1}\|^2\Big]^{\frac{1}{2}} \le \alpha_k \sum_{t=0}^{T-1} \mathbb{E}\Big[\big\|\partial_y \hat{g}(x_k, y^t, \mathcal{D}_g^t)\big\|^2\Big]^{\frac{1}{2}} \le \alpha_k \sum_{t=0}^{T-1}\Big(\sigma_g^2 + L_g^2\mathbb{E}\Big[\|y^t - y^\star(x_k)\|^2\Big]\Big)^{\frac{1}{2}}$$
$$\le \alpha_k \sum_{t=0}^{T-1}\big(\sigma_g^2 + L_g^2 2\alpha_k \sigma_g^2 \mu_g^{-1}\big)^{\frac{1}{2}} + \alpha_k T L_g (E_k^y)^{\frac{1}{2}},$$

where we applied Proposition 7 for every $0 < t \le T-1$ to get the last inequality recalling that R_k^y satisfies $R_{t,k}^y \le 2\alpha_k \sigma_g^2 \mu_g^{-1}$. This directly implies the following bound:

$$\mathbb{E}\Big[\|\tilde{y}_k - y_{k-1}\|^2\Big] \le 2\alpha_k^2 T^2 \big(\sigma_g^2\big(1 + 2L_g^2 \mu_g^{-1}\alpha_k\big) + L_g^2 E_k^y\big). \tag{32}$$

On the other hand, using a triangular inequality and applying Proposition 7, we also have that:

$$\mathbb{E}\Big[\|y_k - y_{k-1}\|^2\Big] \le 2\mathbb{E}\Big[\|y_k - y^\star(x_k)\|^2\Big] + 2\mathbb{E}\Big[\|y_{k-1} - y^\star(x_k)\|^2\Big] \le (4E_k^y + 2R_k^y). \tag{33}$$

The result follows by combining (32) and (33) using coefficients $1-v$ and v. $\qquad\square$

C.3 PROOF OF PROPOSITION 4

Proof of Proposition 4 . We will control each of E_k^y and E_k^z separately.

Upper-bound on E_k^y. Let r_k be a non-increasing sequences between 0 and 1. The following holds:

$$
\begin{aligned}
E_k^y =& \mathbb{E}\Big[\|y_{k-1} - y^\star(x_{k-1})\|^2\Big] + \mathbb{E}\Big[\|y^\star(x_k - y^\star(x_{k-1})\|^2\Big] \\
&+ 2\mathbb{E}\Big[(y_{k-1} - y^\star(x_{k-1}))^\top (y^\star(x_k - y^\star(x_{k-1})))\Big] \\
\overset{(i)}{\leq}& (1+r_k)\mathbb{E}\Big[\|y_{k-1} - y^\star(x_{k-1})\|^2\Big] + \big(1 + r_k^{-1}\big)\mathbb{E}\Big[\|y^\star(x_k - y^\star(x_{k-1})\|^2\Big] \\
\overset{(ii)}{\leq}& (1+r_k)\big(\Lambda_{k-1}E_{k-1}^y + R_{k-1}^y\big) + 2r_k^{-1}\mathbb{E}\Big[\|y^\star(x_k) - y^\star(x_{k-1})\|^2\Big] \\
\overset{(iii)}{\leq}& (1+r_k)\big(\Lambda_{k-1}E_{k-1}^y + R_{k-1}^y\big) + 2L_y^2 r_k^{-1}\mathbb{E}\Big[\|x_k - x_{k-1}\|^2\Big] \\
\overset{(iv)}{\leq}& (1+r_k)\big(\Lambda_{k-1}E_{k-1}^y + R_{k-1}^y\big) + 2L_y^2 r_k^{-1}\gamma_k^2\mathbb{E}\Big[\big\|\hat{\psi}_{k-1}\big\|^2\Big]
\end{aligned}
\tag{34a}
$$

(i) follows by Young's inequality, (ii) uses Proposition 7 to bound the first term and that $(1 + r_k^{-1}) \leq 2r_k^{-1}$ for the second term, (iii) uses that $y^\star$ is L_y-Lipschitz by Proposition 6 and (iv) uses the update equation $x_k = x_{k-1} - \gamma_k\hat{\psi}_k$.

Upper-bound on E_k^z. Similarly, for a non-increasing sequence $0 < \theta_k \leq 1$, we have that:

$$
\begin{aligned}
E_k^z =& \mathbb{E}\Big[\|z_{k-1} - z^\star(x_{k-1}, y_{k-1})\|^2\Big] + \mathbb{E}\Big[\big\|z^\star(x_k, y_k) - z^\star(x_{k-1}, y_{k-1})^2\big\|\Big] \\
&+ 2\mathbb{E}\Big[(z_{k-1} - z^\star(x_{k-1}, y_{k-1}))^\top (z^\star(x_k, y_k) - z^\star(x_{k-1}, y_{k-1}))\Big] \\
\overset{(i)}{\leq}& (1+\theta_k)\mathbb{E}\Big[\|z_{k-1} - z^\star(x_{k-1}, y_{k-1})\|^2\Big] + (1 + \theta_k^{-1})\mathbb{E}\Big[\|z^\star(x_k, y_k) - z^\star(x_{k-1}, y_{k-1})\|^2\Big] \\
\overset{(ii)}{\leq}& (1+\theta_k)\big(\Pi_{k-1}E_{k-1}^z + R_{k-1}^z\big) + 2\theta_k^{-1}\mathbb{E}\Big[\|z^\star(x_k, y_k) - z^\star(x_{k-1}, y_{k-1})\|^2\Big] \\
\overset{(iii)}{\leq}& (1+\theta_k)\big(\Pi_{k-1}E_{k-1}^z + R_{k-1}^z\big) + 4L_z^2\theta_k^{-1}\Big(\mathbb{E}\Big[\|x_k - x_{k-1}\|^2 + \|y_k - y_{k-1}\|^2\Big]\Big) \\
\overset{(iv)}{\leq}& (1+\theta_k)\big(\Pi_{k-1}E_{k-1}^z + R_{k-1}^z\big) + 4L_z^2\theta_k^{-1}\Big(\gamma_k^2\mathbb{E}\Big[\big\|\hat{\psi}_{k-1}\big\|^2\Big] + 2\phi_k E_k^y + 2\tilde{R}_k^y\Big)
\end{aligned}
\tag{35a}
$$

(i) follows by Young's inequality, (ii) uses Proposition 7 to bound the first term and that $(1 + \theta_k^{-1}) \leq 2\theta_k^{-1}$ for the second term, (iii) uses that $z^\star(x, y)$ is L_z-Lipschitz in x and y by Proposition 6 and, finally, (iv) uses the update equation $x_k = x_{k-1} - \gamma_k\hat{\psi}_k$ for the term $\mathbb{E}\Big[\|x_k - x_{k-1}\|^2\Big]$ and Proposition 11 to control the increments $\mathbb{E}\Big[\|y_k - y_{k-1}\|^2\Big]$.

In order to express the upper-bound on E_k^z in terms of E_{k-1}^y instead of E_k^y, we substitute E_k^y in (35a) by its upper-bound in (34a) and use that $(1 + r_k) \leq 2$ to write:

$$
\begin{aligned}
E_k^z \leq& (1+\theta_k)\big(\Pi_{k-1}E_{k-1}^z + R_{k-1}^z\big) + 4L_z^2\theta_k^{-1}\gamma_k^2\big(1 + 4L_y^2\phi_k r_k^{-1}\big)\mathbb{E}\Big[\big\|\hat{\psi}_k\big\|^2\Big] \\
&+ 8L_z^2\theta_k^{-1}\Big(2\phi_k\big(\Lambda_{k-1}E_{k-1}^y + R_{k-1}^y\big) + \tilde{R}_k^y\Big)
\end{aligned}
\tag{36}
$$

We can then express (34a) and (36) jointly in matrix form, thus yielding:

$$
\begin{pmatrix} E_k^y \\ E_k^z \end{pmatrix} \leq \boldsymbol{P_k}\begin{pmatrix} \Lambda_{k-1}E_{k-1}^y + R_{k-1}^y \\ \Pi_{k-1}E_{k-1}^z + R_{k-1}^z \end{pmatrix} + \mathbb{E}\Big[\big\|\hat{\psi}_{k-1}\big\|^2\Big]\boldsymbol{U_k} + \boldsymbol{V_k}
$$

where the $\boldsymbol{P_k}$ is a 2×2 matrix and $\boldsymbol{U_k}$ and $\boldsymbol{V_k}$ are 2-dimensional vectors given by (18). $\qquad\square$

Proof of Proposition 9. Let κ be a positive number such that $0 < \kappa \leq \frac{1}{2}$ and ρ_k a non-increasing sequence with $0 < \rho_k < 1$. By Proposition 2, it follows that E_k^x satisfies the inequality:

$$F_k + E_k^x \leq (1 - (1 - \kappa\rho_k)\delta_k)E_{k-1}^x + \gamma_k s_k V_{k-1}^\psi + \gamma_k\big(s_k + (2\kappa\rho_k)^{-1}\big)E_{k-1}^\psi \tag{37}$$

On the other hand, by Proposition 4 we know that E_k^y and E_k^z satisfy:

$$\begin{pmatrix} E_k^y \\ E_k^z \end{pmatrix} \leq \boldsymbol{P_k} \begin{pmatrix} \Lambda_{k-1}E_{k-1}^y + R_{k-1}^y \\ \Pi_{k-1}E_{k-1}^z + R_{k-1}^z \end{pmatrix} + \mathbb{E}\left[\left\|\hat{\psi}_k\right\|^2\right]\boldsymbol{U_k} + \boldsymbol{V_k} \tag{38}$$

where the $\boldsymbol{P_k}$, $\boldsymbol{U_k}$ and $\boldsymbol{V_k}$ are defined in (18). For conciseness, define $\boldsymbol{S_k}$ and $\boldsymbol{E_k^I}$ to be:

$$\boldsymbol{S_k} := \begin{pmatrix} a_k & 0 \\ 0 & b_k \end{pmatrix}, \qquad \boldsymbol{E_k^I} = \boldsymbol{S_k}\begin{pmatrix} E_k^y \\ E_k^z \end{pmatrix}$$

By assumption, we have that $\gamma_k r_k^{-1} \leq \gamma_{k-1}r_{k-1}^{-1}$ and $\gamma_k \theta_k^{-1} \leq \gamma_{k-1}\theta_{k-1}^{-1}$. Moreover, since α_k and β_k are non-increasing, we also have that $\Lambda_{k-1} \leq \Lambda_k$ and $\Pi_{k-1} \leq \Pi_k$. Hence, it follows that

$$\boldsymbol{S_k P_k S_{k-1}^{-1}}\begin{pmatrix} \Lambda_{k-1} & 0 \\ 0 & \Pi_{k-1} \end{pmatrix} \leq \boldsymbol{S_k P_k S_k^{-1}}\begin{pmatrix} \Lambda_k & 0 \\ 0 & \Pi_k \end{pmatrix} := \tilde{\boldsymbol{P}}_{\boldsymbol{k}}, \tag{39}$$

where (39) holds component-wise. Therefore, multiplying (38) by $\boldsymbol{S_k}$ and using (39) yields:

$$\boldsymbol{E_k^I} \leq \tilde{\boldsymbol{P}}_{\boldsymbol{k}}\boldsymbol{E_{k-1}^I} + \boldsymbol{S_k}\left(\boldsymbol{P_k}\begin{pmatrix} R_{k-1}^y \\ R_{k-1}^z \end{pmatrix} + \mathbb{E}\left[\left\|\hat{\psi}_k\right\|^2\right]\boldsymbol{U_k} + \boldsymbol{V_k}\right) \tag{40}$$

Finally using Proposition 11 to upper-bound $\mathbb{E}\left[\left\|\hat{\psi}_k\right\|^2\right]$, equation (40) becomes:

$$\boldsymbol{E_k^I} \leq \tilde{\boldsymbol{P}}_{\boldsymbol{k}}\boldsymbol{E_{k-1}^I} + \boldsymbol{S_k}\left(\boldsymbol{P_k}\begin{pmatrix} R_{k-1}^y \\ R_{k-1}^z \end{pmatrix} + \left(V_{k-1}^\psi + 2E_{k-1}^\psi + 2\zeta E_{k-1}^x\right)\boldsymbol{U_k} + \boldsymbol{V_k}\right) \tag{41}$$

Furthermore, by Proposition 5 we can bound E_{k-1}^ψ and V_{k-1}^ψ as follows:

$$E_{k-1}^\psi \leq 2L_\psi^2\Big(\Lambda_k(a_k)^{-1}a_{k-1}E_{k-1}^y + \Pi_k(b_k)^{-1}b_{k-1}E_{k-1}^z\Big) + 2L_\psi^2 R_{k-1}^y$$

$$V_{k-1}^\psi \leq w_x^2 + \sigma_x^2\Pi_k(b_k)^{-1}b_{k-1}E_{k-1}^z,$$

where we used again that $\gamma_k r_k^{-1} \leq \gamma_{k-1}r_{k-1}^{-1}$, $\gamma_k \theta_k^{-1} \leq \gamma_{k-1}\theta_{k-1}^{-1}$, $\Lambda_{k-1} \leq \Lambda_k$ and $\Pi_{k-1} \leq \Pi_k$. By summing both inequalities (37) and (41) and substituting all terms E_{k-1}^ψ and V_{k-1}^ψ by their upper-bounds we obtain the following inequality:

$$F_k + E_k^{tot} \leq A_k^x E_{k-1}^x + A_k^y a_{k-1}^y E_{k-1}^y + A_k^z b_{k-1}^z E_{k-1}^z + V_k^{tot}$$

where A_k^x, A_k^y, A_k^z are the components of the vector A_k defined in (20) and V_k^{tot} is the variance term also defined in (20). By taking upper-bounding A_k^x, A_k^y, A_k^z by their maximum value $\|A_k\|_\infty$ the desired inequality follows. $\qquad\square$

D CONTROLLING THE PRECISION OF THE INNER-LEVEL ALGORITHMS.

In this section, we prove Proposition 10. To achieve this, we first provide general conditions on Λ_k and Π_k for controlling the rate $\|A_k\|_\infty$ and which hold regardless of the choice of step-sizes. This is achieved in Proposition 12 of Appendix D.1. Then we prove Proposition 10 in Appendix D.2.

D.1 Controlling Π_k and Λ_k.

We introduce the following quantities:

$$D_k^{(1)} := \frac{1}{1-s}\left[\log\left[\frac{1-(1-\rho_k)\delta_k}{1+2r_k\left[1+2L_\psi^2\left[s_k+2+[4\kappa\rho_k]^{-1}\right]\right]}\right]\right] \tag{42a}$$

$$D_k^{(2)} := \frac{1}{s}\log\left((1-\kappa)(8L_y^2)^{-1}\rho_k\delta_k\zeta_k^{-1}\gamma_k^{-3}r_k^2\right) \tag{42b}$$

$$D_k^{(3)} := -\frac{1}{s}\log\left(4L_y^2\gamma_k^2 r_k^{-2}\right) \tag{42c}$$

$$D_k^{(4)} := \frac{1}{1-s}\log\left[\frac{1-(1-\rho_k)\delta_k}{1+2\theta_k\left[1+\left(2L_\psi^2+\sigma_x^2\right)\left[s_k+2+(4\kappa\rho_k)^{-1}\right]\right]}\right] \tag{42d}$$

$$D_k^{(5)} := -\frac{1}{s}\log\left(16L_z^2\theta_k^{-2}\phi_k\right) \tag{42e}$$

$$D_k^{(6)} := -\frac{1}{s}\log\left(2L_z^2L_y^{-2}\theta_k^{-2}r_k^2\left(1+2L_y^2 r_k^{-1}\phi_k\right)\right) \tag{42f}$$

Proposition 12. *Let $0 < \kappa \leq \frac{1}{2}$ and ρ_k be a non-increasing sequence of positive numbers smaller than 1. Consider Λ_k and Π_k so that:*

$$\log\Lambda_k \leq \min\left(D_k^{(1)}, D_k^{(2)}, D_k^{(3)}\right), \tag{43a}$$

$$\log\Pi_k \leq \min\left(D_k^{(1)}, D_k^{(2)}, \log\Lambda_k + D_k^{(3)}\right). \tag{43b}$$

Then, the following inequalities holds:

$$\|A_k\|_\infty \leq (1-(1-\rho_k)\delta_k), \qquad V_k^{tot} \leq \tilde{V}_k^{tot}.$$

where A_k and V_k^{tot} are defined in (20) of Proposition 9 and $\tilde{V}_k^{tot}$ is defined as:

$$\tilde{V}_k^{tot} := \gamma_k\left(s_k + 4\Lambda_k^s L_y^2\gamma_k^2 r_k^{-2}\right)w_x^2 + \gamma_k\left(2L_\psi^2\left(s_k+(2\kappa\rho_k)^{-1}+2\right)+2+r_k^{-1}\right)R_{k-1}^y$$
$$+ (\Pi_k)^s\left(2\gamma_k\theta_k^{-1}R_{k-1}^z + 8L_z^2\gamma_k\theta_k^{-2}\tilde{R}_k^y\right)$$

Proof. We first prove that $u_k^I \leq u_k^+ \leq 1$ and $p_k^+ \leq 1$ under (43a) and (43b) with u_k^+, p_k^+ given by:

$$u_k^+ := 4L_y^2\gamma_k^2 r_k^{-2}\Lambda_k^s, \qquad p_k^+ = 16L_z^2\Pi_k^s\theta_k^{-2}\phi_k.$$

A direct calculation shows $u_k^+ \leq 1$ whenever (43a) holds. Moreover, recall that $u_k^I = a_k U_k^{(1)} + b_k U_k^{(2)}$ with $U_k^{(1)}$ and $U_k^{(2)}$ being the components of the vector $\boldsymbol{U_k}$ defined in (18). Thus by direct substitution, we get the following expression for u_k^I:

$$u_k^I = 2L_y^2\gamma_k^2 r_k^{-2}\Lambda_k^s\left(1 + 2L_z^2 L_y^{-2}\theta_k^{-2}r_k^2\left(1+2L_y^2 r_k^{-1}\phi_k\right)\frac{\Pi_k^s}{\Lambda_k^s}\right).$$

Therefore, (43b) suffices to ensure that $u_k^I \leq u_k^+$. Finally, (43b) implies directly that $p_k^+ \leq 1$.

We will control each component A_k^x, A_k^y and A_k^z of the vector A_k separately.

Controlling A_k^x. Recalling the expression of A_k^x, the first component of A_k in (20), it holds that:

$$A_k^x = 1 - (1-\kappa\rho_k)\delta_k + 2\zeta_k\lambda_k u_k^I \overset{(i)}{\leq} 1 - (1-\kappa\rho_k)\delta_k + 2\zeta_k\lambda_k u_k^+$$
$$\overset{(ii)}{\leq} 1 - (1-\kappa\rho_k)\delta_k + (1-\kappa)\rho_k\delta_k = 1 - (1-\rho_k)\delta_k.$$

(i) holds since $u_k^I \leq u_k^+$ while (ii) follows from (43a) which ensures that $2\zeta_k\lambda_k u_k^+ \leq (1-\kappa)\rho_k\delta_k$.

Controlling A_k^y. Recall the expression of second component of A_k^y in (20), we have:

$$A_k^y = \Lambda_k^{1-s}\big(1 + r_k\big(1 + 16L_z^2\Pi_k^s\theta_k^{-2}\phi_k + 2L_\psi^2\big(2u_k^I + (s_k + (2\kappa\rho_k)^{-1})\big)\big)\big)$$

$$\overset{(i)}{\le}\Big(1 + 2r_k\Big(1 + 2L_\psi^2(s_k + 2 + (4\kappa\rho_k)^{-1})\Big)\Big)\Lambda_k^{1-s} \overset{(ii)}{\le} 1 - (1 - \kappa\rho_k)\delta_k.$$

(i) holds since $p_k^+ := 16L_z^2\Pi_k^s\theta_k^{-2}\phi_k \le 1$ and $u_k^I \le 1$ while (ii) is a consequence of (43a).

Controlling A_k^z. Similarly, recalling the expression of the third component of A_k we get that:

$$A_k^z = \Pi_k^{1-s}\big(1 + \theta_k\big(1 + \big(2L_\psi^2 + \sigma_x^2\big)\big(2u_k^I + (s_k + (2\kappa\rho_k)^{-1})\big)\big)\big)$$

$$\overset{(i)}{\le}\Pi_k^{1-s}\big(1 + \theta_k\big(1 + \big(2L_\psi^2 + \sigma_x^2\big)\big(2 + (s_k + (2\kappa\rho_k)^{-1})\big)\big)\big) \overset{(ii)}{\le} 1 - (1 - \kappa\rho_k)\delta_k,$$

where (i) uses that $u_k^I \le 1$ and (ii) follows from (43b).

Controlling V_k^{tot}. Recalling the expression of V_k^{tot} from (20), we have that:

$$V_k^{tot} := \gamma_k\big(\Lambda_k^s(1 + r_k^{-1}) + 16L_z^2\phi_k\theta_k^{-2}\Pi_k^s + 2L_\psi^2\big(2u_k^I + (s_k + (2\kappa\rho_k)^{-1})\big)\big)R_{k-1}^y$$

$$+ \gamma_k\Big(\big(1 + \theta_k^{-1}\big)\Pi_k^s R_{k-1}^z + 8L_z^2\theta_k^{-2}\Pi_k^s\tilde{R}_k^y\Big) + \gamma_k\big(s_k + u_k^I\big)w_x^2$$

$$\le \gamma_k\big(3r_k^{-1} + 2L_\psi^2\big(2 + (s_k + (2\kappa\rho_k)^{-1})\big)\big)R_{k-1}^y$$

$$+ \gamma_k\Big(2\theta_k^{-1}\Pi_k^s R_{k-1}^z + 8L_z^2\theta_k^{-2}\Pi_k^s\tilde{R}_k^y\Big) + \gamma_k\big(s_k + u_k^+\big)w_x^2 = \tilde{V}_k^{tot}.$$

where we use $16L_z^2\phi_k\theta_k^{-2}\Pi_k^s \le 1$ and $u_k^I \le 1$ for the first line and $u_k^I \le u_k^+$ for the last line. $\qquad\square$

D.2 Controlling the number of inner-level iterations

We provide now a proof of Proposition 10 which is a consequence of Proposition 12.

Proof of Proposition 10. We first provide conditions on the number of iterations T and N of algorithms $\mathcal{A}_k$ and $\mathcal{B}_k$ to control the rate $\|A_k\|_\infty$ and then provide an upper-bound on V_k^{tot}.

Conditions on T and N. We consider the setting with constant step-size $\gamma_k = \gamma$, $\alpha_k = \alpha$ and $\beta_k = \beta$ and choose $r_k = \theta_k = 1$ and $\delta_k = \delta_0$ for some $0 < \delta_0 < 1$. We also take $v = 1$ so that $\phi_k = 2$ and $\tilde{R}_k^y = R_k^y$. By direct substitution of the parameters r_k, θ_k, ϕ_k, γ_k, δ_k, ρ_k and ζ_k, in the expressions of $D_k^{(1)}$, $D_k^{(2)}$, $D_k^{(3)}$, $D_k^{(4)}$, $D_k^{(5)}$ and $D_k^{(6)}$ defined in (42a) to (42f), we verify that:

$$-A_1 \le D_k^{(1)}, \qquad -A_1' \le D_k^{(4)}, \qquad -A_2 \le \min\Big(D_k^{(2)}, D_k^{(3)}\Big), \qquad -A_2' \le \min D_k^{(5)}, D_k^{(6)}.$$

Hence, we can ensure the conditions of Proposition 12 hold by choosing T and N so that:

$$\log\Lambda_k \le -\max(1, A_1, A_2), \qquad \log\Pi_k \le -\log\Lambda_k - \max(1, A_1', A_2').$$

This is achieved by for the following choice:

$$T = \lfloor\kappa_g\max(1, A_1, A_2)\rfloor + 1, \qquad N = \lfloor2\kappa_g(\max(1, A_1, A_2)) + \max(1, A_1', A_2')\rfloor + 1. \quad (44)$$

Hence, for such choice, we are guaranteed by Proposition 12 that $\|A_k\|_\infty \le 1 - (1 - \rho_k)\delta_k$.

Bound on the variance V_k^{tot}. By choosing T and N as in (44), we know that Λ_k and Π_k satisfy (43), so that the variance term V_k^{tot} is upper-bounded by $\tilde{V}_k^{tot}$. Moreover, by direct substitution of the sequences appearing in the expression of $\tilde{V}_k^{tot}$ by their values, we get:

$$\tilde{V}_k^{tot} \le \gamma\left(\frac{1-u}{2} + \left(\frac{1+u}{2}\delta_0 + 4L_y^2\gamma_0^2\right)\right)w_x^2 + \gamma\Big(2R_{k-1}^z + 8L_z^2\tilde{R}_k^y\Big)$$

$$+ \gamma\big(2L_\psi^2\big(3 + (2\kappa\rho_0)^{-1}\big) + 3\big)R_{k-1}^y.$$

Furthermore, by definition of R_{k-1}^y and R_{k-1}^z, we have that:

$$R_{k-1}^y \le 2\mu_g^{-1}\alpha_0\sigma_g^2, \qquad R_{k-1}^z \le \mu_g^{-3}\Big(B^2\beta\sigma_{g_{yy}}^2 + 3\mu_g\sigma_f^2\Big).$$

Moreover, recall that $\tilde{R}_k^y = R_k^y$ since we chose $v = 1$. Thus $\tilde{R}_k^y \le 2\mu_g^{-1}\alpha_0\sigma_g^2$. This yields the inequality $\tilde{V}_k^{tot} \le \gamma\mathcal{W}$, with $\mathcal{W}$ given by (22). $\qquad\square$

E STOCHASTIC LINEAR DYNAMICAL SYSTEM WITH CORRELATED NOISE

Let A be a positive definite matrix in $\mathbb{R}^d \times \mathbb{R}^d$ satisfying $\mu_g|\sigma_i(A)| \leq L_g$ and b a vector in $\mathbb{R}^d$. We denote by $z^\star = -A^{-1}b$. Consider A_m be a sequence of i.i.d. positive symmetric matrices in $\mathbb{R}^d \times \mathbb{R}^d$ such that $\mathbb{E}[A_m] = A$, and $\hat{b}$ a random vector in $\mathbb{R}^d$ such that $\mathbb{E}\left[\hat{b}\right] = b$, with A_m and $\hat{b}$ being mutually independent. Let β be such that $\beta \leq \frac{1}{L_g}$. Define $\Sigma_A = \mathbb{E}\left[(A_n - A)^\top (A_n - A)\right]$ and denote by σ_A and L_A the largest singular values of Σ_A and $A^{-1}\Sigma_A A^{-1}$. Finally let σ_c^2 be an upper-bound on $\mathbb{E}\left[\left\|\hat{b} - b\right\|^2\right]$. Let z and z' be two vectors in $\mathbb{R}^d$ and define the iterates z^n and $\bar{z}^n$ such that $z^0 = z$ and $\bar{z}^0 = z'$ and using the recursion:

$$z^n = (I - \beta A_n)z^{n-1} - \beta\hat{b}, \bar{z}^n = (I - \beta A)\bar{z}^{n-1} - \beta b.$$

Hence, from the definition of z^n and $\bar{z}^n$ we directly have that:

$$z^n = \prod_{t=1}^n (I - \beta A_n)z - \sum_{t=1}^n \prod_{j=t+1}^n \beta(I - \beta A_j)\hat{b}, \quad \bar{z}^n = \prod_{t=1}^n (I - \beta A)z' - \sum_{t=1}^n \prod_{j=t+1}^n \beta(I - \beta A)b.$$

The next proposition computes the bias $\mathbb{E}\left[z^n - \bar{z}^n \big| \hat{b}\right]$.

Proposition 13. *The following identities hold:*

$$\mathbb{E}\left[z^n - \bar{z}^n \big| \hat{b}\right] = (I - \beta A)^n(z - z') - \beta\left(\sum_{t=1}^n (I - \beta A)^{n-t}\right)\left(\hat{b} - b\right)$$

$$\mathbb{E}[z^n - \bar{z}^n] = (I - \beta A)^n(z - z')$$

Proof. The proof is a consequence of A_n and $\hat{b}$ being i.i.d. and unbiased estimates of A and b. $\qquad\square$

The next proposition controls the mean squared errors $\mathbb{E}\left[\|z^n - z^\star\|^2\right]$ and $\mathbb{E}\left[\|z^n - \bar{z}^n\|^2\right]$.

Proposition 14. *Define $n^\star(n, \beta) = \min(n, \frac{1}{\beta\mu_g})$ Then the following inequalities holds:*

$$\mathbb{E}\left[\|z^n - z^\star\|^2\right] \leq (1 - \beta\mu_g)^n\|z - z^\star\|^2 + \beta^2\left(\sigma_A^2\|z^\star\|^2 + 3\left(n \wedge \frac{1}{\beta\mu_g}\right)\sigma_c^2\right)\left(n \wedge \frac{1}{\beta\mu_g}\right),$$

$$\|\bar{z}^n - z^\star\|^2 \leq (1 - \beta\mu_g)^n\|z' - z^\star\|^2.$$

Moreover, if $z' = z$, then we have:

$$\mathbb{E}\left[\|z^n - \bar{z}^n\|^2\right] \leq 4\mu_g^{-2}\sigma_A^2(1 - \frac{\beta\mu_g}{2})^n\|z - z^\star\|^2 \tag{45}$$

$$+ 2\beta^2\left(\sigma_A^2\|z^\star\|^2 + 3\left(n \wedge \frac{1}{\beta\mu_g}\right)\sigma_c^2\right)\left(n \wedge \frac{1}{\beta\mu_g}\right).$$

Proof. It is straightforward to see that:

$$\|\bar{z}^n - z^\star\|^2 \leq (1 - \beta\mu_g)\|z' - z^\star\|^2.$$

Now, let's control $\mathbb{E}\left[\|z^n - \bar{z}^n\|^2\right]$. The following identity holds by definition of z^n and $\bar{z}^n$:

$$\mathbb{E}\left[\|z^n - \bar{z}^n\|^2\right] = \mathbb{E}\left[(z^{n-1} - \bar{z}^{n-1})^\top\left((I - \beta A)^2 + \beta^2\Sigma_A\right)(z^{n-1} - \bar{z}^{n-1})\right]$$

$$- 2\beta\mathbb{E}\left[(z^{n-1} - \bar{z}^{n-1})^\top(I - \beta A)(\hat{b} - b)\right]$$

$$+ \beta^2\left(\mathbb{E}\left[\left\|\hat{b} - b\right\|^2\right] + \mathbb{E}\left[(\bar{z}^{n-1})^\top\Sigma_A\bar{z}^{n-1}\right]\right)$$

$$= \mathbb{E}\left[(z^{n-1} - \bar{z}^{n-1})^\top\left((I - \beta A)^2 + \beta^2\Sigma_A\right)(z^{n-1} - \bar{z}^{n-1})\right]$$

$$\beta^2\left((\bar{z}^{n-1})^\top\Sigma_A\bar{z}^{n-1} + \mathbb{E}\left[\left(\hat{b} - b\right)^\top(I + 2(I - \beta A)D_n)\left(\hat{b} - b\right)\right]\right)$$

where $\Sigma_A = \mathbb{E}\left[(A_n - A)^\top (A_n - A)\right]$ and $D_n = \sum_{t=0}^{n-1}(I - \beta A)^t$. By simple calculation we can upper-bound the last term by:

$$\beta^2 \mathbb{E}\left[\left(\hat{b} - b\right)^\top (I + 2(I - \beta A)D_n)\left(\hat{b} - b\right)\right] \leq 3\left(n \wedge \frac{1}{\beta\mu_g}\right)\beta^2 \mathbb{E}\left[\left\|\hat{b} - b\right\|^2\right].$$

Moreover, provided that $\beta \leq \frac{1}{L_g(1+L_A)}$, where L_A is the highest eigenvalue of $A^{-1}\Sigma_A A^{-1}$, then we have the following:

$$\mathbb{E}\left[\|z^n - \bar{z}^n\|^2\right] \leq (1 - \beta\mu_g)\mathbb{E}\left[\|z^{n-1} - \bar{z}^{n-1}\|^2\right]$$
$$+ \beta^2\left(\left(\bar{z}^{n-1}\right)^\top \Sigma_A \bar{z}^{n-1} + 3\left(n \wedge \frac{1}{\beta\mu_g}\right)\mathbb{E}\left[\left\|\hat{b} - b\right\|^2\right]\right)$$
$$\leq (1 - \beta\mu_g)\mathbb{E}\left[\|z^{n-1} - \bar{z}^{n-1}\|^2\right] + \beta^2\left(\sigma_A^2\|\bar{z}^{n-1}\|^2 + 3\left(n \wedge \frac{1}{\beta\mu_g}\right)\sigma_c^2\right),$$

Unrolling the recursion, it follows that:

$$\mathbb{E}\left[\|z^n - \bar{z}^n\|^2\right] \leq (1 - \beta\mu_g)^n\|z - z'\|^2 + \beta^2 \sum_{t=1}^{n}(1 - \beta\mu_g)^{n-t}\left(\sigma_A^2\|\bar{z}^{n-1}\|^2 + 3\left(n \wedge \frac{1}{\beta\mu_g}\right)\sigma_c^2\right) \tag{46}$$

In particular, if $z' = z^\star$, then $\bar{z}^n = z^\star$ and we get:

$$\mathbb{E}\left[\|z^n - z^\star\|^2\right] \leq (1 - \beta\mu_g)^n\|z - z^\star\|^2 + \beta^2\left(\sigma_A^2\|z^\star\|^2 + 3\left(n \wedge \frac{1}{\beta\mu_g}\right)\sigma_c^2\right)\left(n \wedge \frac{1}{\beta\mu_g}\right).$$

To get the last inequality, we simply choose $z' = z$ and recall that:

$$\left\|\bar{z}^{t-1}\right\| \leq 2\left((1 - \beta\mu_g)^{t-1}\|z' - z^\star\| + \|z^\star\|\right).$$

Using the above in in (46) yields:

$$\mathbb{E}\left[\|z^n - \bar{z}^n\|^2\right] \leq 2n\beta^2\sigma_A^2(1 - \beta\mu_g)^{n-1}\|z - z^\star\|^2 + 2\beta^2\left(\sigma_A^2\|z^\star\|^2 + 3\left(n \wedge \frac{1}{\beta\mu_g}\right)\sigma_c^2\right)\left(n \wedge \frac{1}{\beta\mu_g}\right).$$

Moreover, by Lemma 2 we know that $n\beta^2\mu_g^2(1 - \beta\mu_g)^{n-1} \leq (1 - \frac{\beta\mu_g}{2})^{n-1}$ and since $\beta\mu_g \leq 1$, we have that $(1 - \frac{\beta\mu_g}{2})^{-1} \leq 2$ so that $(1 - \frac{\beta\mu_g}{2})^{n-1} \leq 2(1 - \frac{\beta\mu_g}{2})^n$. Hence, we can write:

$$\mathbb{E}\left[\|z^n - \bar{z}^n\|^2\right] \leq 4\mu_g^{-2}\sigma_A^2(1 - \frac{\beta\mu_g}{2})^n\|z - z^\star\|^2 + 2\beta^2\left(\sigma_A^2\|z^\star\|^2 + 3\left(n \wedge \frac{1}{\beta\mu_g}\right)\sigma_c^2\right)\left(n \wedge \frac{1}{\beta\mu_g}\right).$$

$\square$

Lemma 2. *Let $0 \leq b < 1$ and $n \geq 1$, then the following inequality holds:*

$$nb^2(1 - b)^{n-1} \leq (1 - \frac{b}{2})^{n-1}.$$

Proof. We consider the function $h(n, b)$ defined by:

$$h(n, b) := (n - 1)\log\left(\frac{1 - \frac{b}{2}}{1 - b}\right) - \log(nb^2).$$

We need to show that $h(n, b)$ is non-negative for any $n \geq 1$ and $0 \leq b < 1$. For this purpose, we fix b and consider the variations of $h(n, b)$ in n:

$$\partial_n h(n, b) = \log\left(\frac{1 - \frac{b}{2}}{1 - b}\right) - \frac{1}{n}.$$

$\partial_n h(n, b)$ is non-negative for $n \geq n^\star = \log\left(\frac{1-\frac{b}{2}}{1-b}\right)^{-1}$ and non-positive for all $n \leq n^\star$. Hence, $h(n, b)$ achieves its minimum value in $n^\star$ over the $(0, +\infty)$. We distinguish two case depending on whether $n^\star$ is greater of smaller than 1.

Case $n^\star \leq 1$. In this case $n \mapsto h(n, b)$ is increasing on the interval $[1, +\infty)$ since $\partial_n h(n, b) \geq 0$ for $n \geq n^\star$. Hence, $h(n, b) \geq h(1, b)$ for all $n \geq 1$. Moreover, since $h(1, b) = -\log(b^2) \geq 0$ the result follows directly.

Case $n^\star > 1$. In this case we still have $h(n, b) \geq h(n^\star, b)$ for all $n \geq 1$, since $n^\star$ achieves the minimum value of h. Thus we only need to show that $h(n^\star, b) \geq 0$. Using the expression of $n^\star$, we have:

$$
\begin{aligned}
h(n^\star, b) &= 1 - \log\left(\frac{1-\frac{b}{2}}{1-b}\right) - \log\left(n^\star b^2\right) \\
&= 1 - \frac{1}{n^\star} - \log\left(n^\star b^2\right)
\end{aligned}
$$

Since $n^\star > 1$, the first term $1 - \frac{1}{n^\star}$ is non-negative, thus we only need to show that $n^\star b^2 \leq 1$ so that the last term is also non-negative. It is easy to see that $n^\star b^2 \leq 1$ is equivalent to having $\tilde{h}(b) \geq 0$, where we define the function $\tilde{h}(b)$ as:

$$
\tilde{h}(b) = \log\frac{1-\frac{b}{2}}{1-b} - b^2.
$$

We can analyze the variations of $\tilde{b}$ be computing its derivative which is given by:

$$
\partial_b \tilde{h}(b) = \frac{1}{(1-b)(2-b)} - 2b.
$$

Hence, we have the following equivalence:

$$
\partial_b \tilde{h}(b) \geq 0 \iff 2b(1-b)(2-b) \leq 1
$$

This is always true for $0 \leq b < 1$ since $b(1-b) \leq \frac{1}{4}$ so that $2b(1-b)(2-b) \leq \frac{2-b}{2} \leq 1$. Thus we have shown that $\tilde{h}$ is increasing over $[0, 1)$ so that $\tilde{h}(b) \geq \tilde{h}(0) = 0$. As discussed above, this is equivalent to having $n^\star b^2 \leq 1$, so that $h(n^\star, b) \geq 0$ which concludes the proof.

$\square$

F   EXPERIMENTS

F.1   EXPERIMENTAL DETAILS FOR LOGISTIC REGRESSION

The inner-level and outer-level cost functions for such task take the following form:

$$
f(x, y) = \frac{1}{|\mathcal{D}_{val}|} \sum_{\xi \in \mathcal{D}_{val}} L(y, \xi), \qquad g(x, y) = \frac{1}{|\mathcal{D}_{tr}|} \sum_{\xi \in \mathcal{D}_{tr}} L(y, \xi) + \frac{1}{pd} \sum_{i=1} \exp(x_i) y_{i,j}^2
$$

We use a setting similar to Grazzi et al. (2020); Ji et al. (2021) where we set the default step-sizes to $\alpha_k = \gamma_k = 100$ and $\beta_k = 0.5$, and $T = N = 10$. When performing a grid-search, we select the best performing parameters T, M and β_k from a grid $\{10, 20\} \times \{5, 10\} \times \{0.5, 10\}$, while the batch-size (chosen to be the same for all steps of the algorithms) varies from $10 * \{0.1, 1, 2, 4\}$.

Figure 2: Evolution of the relative error vs. time in seconds for different AID based methods on the synthetic example. Each column corresponds to a method (AID-CG, AmIGO-CG, AmIGO-GD, AID-N, AID-FP) and each row corresponds to a choice of the conditioning number κ_g. For each method we consider T and N from a grid $\{1, 10, 10^2, 10^3\} \times \{1, 10, 10^2, 10^3\}$. Lightest colors corresponds to smaller values of N while nuances within each color correspond to increasing values of T.

Figure 3: Evolution of the relative error vs. time in seconds for different ITD based methods on the synthetic example. From the left to the right, the first two columns correspond to Reverse and ITD method small conditioning numbers $\kappa_g \in \{1, 10, 10^3\}$, last two column are for higher conditioning numbers $\kappa_g \in \{10^4, 10^5, 10^7\}$. For each method we consider $T \in \{1, 10, 10^2, 10^3\}$. Lightest colors correspond to smaller values of T.

Figure 4: Evolution of the validation loss (left column), validation accuracy (middle column) and test accuracy (right column) in time (s) for different methods on the logistic regression task. Each row correspond to different choices for the size of the batch $|\mathcal{D}| \in \{100, 1000, 2000, 4000\}$ chosen to be the same for all gradient, Hessian and Jacobian-vector products evaluations. Time is reported in seconds.