# Primal-Dual Damping algorithms for optimization 

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#### Abstract

We propose an unconstrained optimization method based on the well-known primal-dual hybrid gradient (PDHG) algorithm. We first formulate the optimality condition of the unconstrained optimization problem as a saddle point problem. We then compute the minimizer by applying generalized primal-dual hybrid gradient algorithms. Theoretically, we demonstrate the continuous-time limit of the proposed algorithm forms a class of second-order differential equations, which contains and extends the heavy ball ODEs and Hessian-driven damping dynamics. Following the Lyapunov analysis of the ODE system, we prove the linear convergence of the algorithm for strongly convex functions. Experimentally, we showcase the advantage of algorithms on several convex and nonconvex optimization problems by comparing the performance with other well-known algorithms, such as Nesterov's accelerated gradient methods. In particular, we demonstrate that our algorithm is efficient in training two-layer and convolution neural networks in supervised learning problems.


Keywords and phrases: Optimization, Primal-dual hybrid gradient algorithms, Primal-dual damping dynamics.

## 1. Introduction

Optimization is one of the essential building blocks in many applications, including scientific computing and machine learning problems. One of the classical algorithms for unconstrained optimization problems is the gradient descent method, which updates the state variable in the negative gradient direction at each step [4]. Nowadays, accelerated gradient descent methods have been widely studied. Typical examples include Nesterov's accelerated gradient method [16], Polyak's heavy ball method [22], and Hessian-driven damping methods [7, 1, 2, 3].

[^0]On the other hand, some first-order methods are introduced to solve linear-constrained optimization problems. Typical examples include the primal-dual hybrid gradient (PDHG) method [6] and the alternating direction method of multipliers (ADMM) [5, 9]. They are designed to solve an infsup saddle point type problem, which updates the gradient descent direction for the minimization variable and applies the gradient ascent direction for the maximization variable. Both PDHG and ADMM are designed for solving optimization problems with affine constraints. Ouyang et al. [18] proposed accelerated linearized ADMM, which incorporates a multi-step acceleration scheme into linearized ADMM. Recently, the PDHG method has been extended into solving nonlinear-constrained minimization problems [27].

In this paper, we study a general class of accelerated first-order methods for unconstrained optimization problems. We reformulate the original optimization problem into an inf-sup type saddle point problem, whose saddle point solves the optimality condition. We then apply a linearized preconditioned primal-dual hybrid gradient algorithm to compute the proposed saddle point problem.

The main description of the algorithm is as follows. Consider the following inf-sup problem for a $\mathcal{C}^{2}$ strongly convex function $f$ over $\mathbb{R}^{d}$

$$
\begin{equation*}
\inf _{\boldsymbol{x} \in \mathbb{R}^{d}} \sup _{\boldsymbol{p} \in \mathbb{R}^{d}}\langle\nabla f(\boldsymbol{x}), \boldsymbol{p}\rangle-\frac{\varepsilon}{2}\|\boldsymbol{p}\|^{2}, \tag{1.1}
\end{equation*}
$$

where $\boldsymbol{p}$ is a constructed "dual variable", $\varepsilon>0$ is a constant, $\langle\cdot, \cdot\rangle$ is an Euclidean inner product, and $\|\cdot\|$ is an Euclidean norm. We later prove that the solution to the saddle point problem (1.1) gives the global minimum of $f$. We propose a linearized preconditioned PDHG algorithm for solving the above inf-sup problem:

$$
\begin{align*}
\boldsymbol{p}^{n+1} & =\boldsymbol{p}^{n}+\sigma \boldsymbol{A}\left(\boldsymbol{x}^{n}\right) \nabla f\left(\boldsymbol{x}^{n}\right)-\sigma \varepsilon \boldsymbol{A}\left(\boldsymbol{x}^{n}\right) \boldsymbol{p}^{n+1},  \tag{1.2a}\\
\tilde{\boldsymbol{p}}^{n+1} & =\boldsymbol{p}^{n+1}+\omega\left(\boldsymbol{p}^{n+1}-\boldsymbol{p}^{n}\right),  \tag{1.2b}\\
\boldsymbol{x}^{n+1} & =\boldsymbol{x}^{n}-\tau \boldsymbol{C}\left(\boldsymbol{x}^{n}\right) \tilde{\boldsymbol{p}}^{n+1} \tag{1.2c}
\end{align*}
$$

where $n=1,2, \cdots$ is the iteration step, $\tau, \sigma>0$ are stepsizes for the updates of $\boldsymbol{x}, \boldsymbol{p}$, respectively, and $\omega>0$ is a parameter. $\tilde{\boldsymbol{p}}^{n+1}$ is an extrapolation of $\boldsymbol{p}^{n+1}$ in the direction of $\boldsymbol{p}^{n+1}-\boldsymbol{p}^{n}$. This extrapolation term is motivated by the PDHG method [6]. In the above algorithm, $\boldsymbol{C}\left(\boldsymbol{x}^{n}\right)=\boldsymbol{B}\left(\boldsymbol{x}^{n}\right) \nabla^{2} f\left(\boldsymbol{x}^{n}\right)$, where $\boldsymbol{A}\left(\boldsymbol{x}^{n}\right) \in \mathbb{R}^{d \times d}$, and $\boldsymbol{B}\left(\boldsymbol{x}^{n}\right) \in \mathbb{R}^{d \times d}$ act as preconditioners on the updates of $\boldsymbol{p}^{n+1}$ and $\boldsymbol{x}^{n+1}$, respectively. This paper will only focus on the simple case where $\boldsymbol{A}(\boldsymbol{x})=A \mathbb{I}$ for some constant $A>0$. Although there
is a second-order term $\nabla^{2} f\left(\boldsymbol{x}^{n}\right)$ in the update of $\boldsymbol{x}^{n+1}$ (hidden in $\boldsymbol{C}\left(\boldsymbol{x}^{n}\right)$ ), our algorithm is still a first-order algorithm by choosing $\boldsymbol{B}\left(\boldsymbol{x}^{n}\right) \nabla^{2} f\left(\boldsymbol{x}^{n}\right)=$ $\boldsymbol{C}\left(\boldsymbol{x}^{n}\right)$ for some $\boldsymbol{C}$ that is easy to compute. For example, we test that $\boldsymbol{C}=\mathbb{I}$ is a very good choice in most of our numerical examples. See empirical choices of parameters in our numerical sections.

Our method forms a class of ordinary differential equation systems in terms of ( $\boldsymbol{x}, \boldsymbol{p}$ ) in the continuous limit $\tau, \sigma \rightarrow 0$. We call it the primal-dual damping (PDD) dynamics. We show that the PDD dynamics form a class of second-order ODEs, which contains and extends the inertia Hessian-driven damping dynamics $[7,1]$. Theoretically, we analyze the convergence property of PDD dynamics. If $f$ is a quadratic function of $\boldsymbol{x}$, with constant $\boldsymbol{A}, \boldsymbol{B}$, the PDD dynamic satisfies a linear ODE system. Under suitable choices of parameters, we obtain a similar convergence acceleration in heavy ball ODE [24]. Moreover, for general nonlinear function $f$, we have the following informal theorem characterizing the convergence speed of our algorithm:

Theorem 1.1 (Informal). Let $f: \mathbb{R}^{d} \rightarrow \mathbb{R}$ be a $\mathcal{C}^{4}$ strongly convex function. Let $\boldsymbol{x}^{*}$ be the global minimum of $f$ and $\boldsymbol{p}^{*}=0$. Then, the iteration $\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)$ produced by (1.2) converges to the saddle point $\left(\boldsymbol{x}^{*}, \boldsymbol{p}^{*}\right)$ if $\tau$, $\sigma$, are small enough. Moreover,

$$
\left\|\boldsymbol{p}^{n}\right\|^{2}+\left\|\nabla f\left(\boldsymbol{x}^{n}\right)\right\|^{2} \leq\left(\left\|\boldsymbol{p}^{0}\right\|^{2}+\left\|\nabla f\left(\boldsymbol{x}^{0}\right)\right\|^{2}\right)\left(1-\frac{\mu^{2}}{M+\delta}\right)^{n},
$$

where $\mu=\min _{\boldsymbol{x}} \lambda_{\min }\left(\nabla^{2} f(\boldsymbol{x}) \boldsymbol{C}(\boldsymbol{x})\right), \boldsymbol{C}(\boldsymbol{x})=\boldsymbol{B}(\boldsymbol{x}) \nabla^{2} f(\boldsymbol{x}), \delta>0$ depends on the initial condition, and $M>0$ depends on $\boldsymbol{C}(\boldsymbol{x})^{T}\left(\nabla^{3} f(\boldsymbol{x}) \nabla f(\boldsymbol{x})+\right.$ $\left.\left(\nabla^{2} f(\boldsymbol{x})\right)^{2}\right) \boldsymbol{C}(\boldsymbol{x}), \tau, \sigma, A, \varepsilon$, and $\omega$. The detailed version is given in Theorem 3.10 .

Numerically, we test the algorithms in both convex and non-convex optimization problems. In convex optimization, we demonstrate the fast convergence results of the proposed algorithm with selected preconditioners, compared with the gradient descent method, Nesterov accelerated gradient method, and Hessian-driven damping method. This justifies the convergence analysis. We also test our algorithm for several well-known non-convex optimization problems. Some examples, such as the Rosenbrock and Ackley functions, demonstrate the potential advantage of our algorithms in converging to the global minimizer. In particular, we compare our algorithms with some popular optimizers including stochastic gradient descent, Nesterov, Adam, and Hessian-driven damping for training two-layer and convolutional neural network functions in supervised learning problems. This showcases the
potential advantage of the proposed methods in terms of convergence speed and test accuracy.

PDHG has been widely used in linear-constrained optimization problems [6]. Recently, Valkonen [27] applied the PDHG for nonlinearly constrained optimization problems. They proved the asymptotic convergence for the nonlinear coupling saddle point problems. It is different from our PDHG algorithm for computing unconstrained optimizations. In our paper, we show linear convergence for a particular nonlinear coupling saddle point problem. Meanwhile, Nesterov accelerated gradient method and Hessiandriven damping algorithms can also be formulated in both discrete-time updates and continuous-time second-order ODEs. Wibisono et al. [28] also introduced the idea of Bregman Lagrangian to study a family of accelerated methods in continuous time limit. It forms a nonlinear second-order ODE. Compared to them, our PDD algorithm induces a generalized second-order ODE system, which contains both heavy ball ODE [24] and Hessian-driven damping dynamics $[7,1,2,3]$. For example, when $C=\mathbb{I}$, algorithm (1.2) can be viewed as a different time discretization of Hessian-driven damping dynamics $[7,1]$. It provides a different update in discrete time update. We only evaluate the gradient of $f$ once, whereas Attouch's algorithm [2] evaluates the gradient of $f$ twice. In numerical experiments, we demonstrate that the proposed algorithm outperforms Nesterov accelerated gradient method and Hessian-driven damping methods in some non-convex optimization problems, including supervised learning problems for training neural network functions.

Our work is also related to preconditioning, an important technique in numerical linear algebra [26] and numerical PDEs [23, 19]. In general, preconditioning aims to reduce the condition number of some operators to improve convergence speed. One famous example would be preconditioning gradient descent by the inverse of the Hessian matrix, which gives rise to Newton's method. In recent years, preconditioning techniques have also been developed in training neural networks [17, 11]. Adam [11] is arguably one of the most popular optimizers in training deep neural networks. It can also be viewed as a preconditioned algorithm using a diagonal preconditioner that approximates the diagonal of the Fisher information matrix [20]. Shortly after Chambolle and Pock [6] developed PDHG for constrained optimization, the same authors also studied preconditioned PDHG method [21], in which they developed a simple diagonal preconditioner that can guarantee convergence without the need to compute step sizes. Liu et al. [15] proposed non-diagonal preconditioners for PDHG and showed close connections between preconditioned PDHG and ADMM. Park et al. [19] studied
the preconditioned Nesterov's accelerated gradient method and proved convergence in the induced norm. Jacobs et al. [10] introduced a preconditioned norm in the primal update of the PDHG method and improved the step size restriction of the PDHG method.

Our paper is organized as follows. In Section 2 we provide some background and derivations of our algorithm. We also provide the ODE formulations for our primal-dual damping dynamics. In Section 3, we prove our main convergence results for the algorithm. In Section 4 we showcase the advantage of our algorithm through several convex and non-convex examples. In particular, we show that our algorithm can train neural networks and is competitive with commonly used optimizers, such as SGD with Nesterov's momentum and Adam. We conclude in Section 5 with more discussions and future directions.

## 2. Primal-dual damping algorithms for optimizations

In this section, we first review PDHG algorithms for constrained optimization problems. We then construct a saddle point problem for the unconstrained optimization problem and apply the preconditioned PDHG algorithm to compute the proposed saddle point problem. We last derive an ODE system, which takes the limit of stepsizes in the PDHG algorithm. It forms a second-order ODE, which generalizes the Hessian-driven damping dynamics. We analyze the convergence properties of the ODE system for quadratic optimization problems.

### 2.1. Review PDHG for constrained optimization

In Chambolle and Pock [6], the following saddle point problem was considered:

$$
\begin{equation*}
\min _{x \in X} \max _{y \in Y}\langle K x, y\rangle+G(x)-F^{*}(y) \tag{2.1}
\end{equation*}
$$

where $X$ and $Y$ are two finite-dimensional real vector spaces equipped with inner product $\langle\cdot, \cdot\rangle$ and norm $\|\cdot\|=\langle\cdot, \cdot\rangle^{1 / 2}$. The map $K: X \rightarrow Y$ is a continuous linear operator. $G: X \rightarrow[0,+\infty]$ and $F^{*}: Y \rightarrow[0,+\infty]$ are proper, convex, lower semi-continuous (l.s.c.) functions. $F^{*}$ is the convex conjugate of a convex l.s.c. function $F$. It is straightforward to verify that (2.1) is the primal-dual formulation of the nonlinear primal problem

$$
\min _{x \in X} F(K x)+G(x) .
$$

Then the PDHG algorithm for saddle point problem (2.1) is given by

$$
\begin{align*}
& y^{n+1}=\left(I+\sigma \partial F^{*}\right)^{-1}\left(y^{n}+\sigma K \tilde{x}^{n}\right),  \tag{2.2a}\\
& x^{n+1}=(I+\tau \partial G)^{-1}\left(x^{n}-\tau K^{*} y^{n+1}\right),  \tag{2.2b}\\
& \tilde{x}^{n+1}=x^{n+1}+\omega\left(x^{n+1}-x^{n}\right), \tag{2.2c}
\end{align*}
$$

where $\tilde{x}^{n+1}$ is an extrapolation of $x^{n+1}$ in the direction of $x^{n+1}-x^{n}$. It is worthwhile noting that $\omega=0$ corresponds to the semi-implicit classical Arrow-Hurwicz algorithm which has been presented as an efficient approach for solving some type of imaging problems [29]. $I$ is the identity operator and $(I+\sigma \partial F)^{-1}$ is the resolvent operator, which is defined the same way as the proximal operator

$$
\begin{aligned}
(I+\tau \partial F)^{-1}(y) & =\underset{x}{\arg \min } \frac{\|x-y\|^{2}}{2 \tau}+F(x) \\
& =\operatorname{prox}_{\tau F}(y)
\end{aligned}
$$

When $\omega=1$, Chambolle and Pock [6] proved convergence if $\tau \sigma\|K\|^{2}<$ 1 , where $\|\cdot\|$ is the induced operator norm. It is worth noting that the convergence analysis requires that $K$ is a linear operator.

### 2.2. Saddle point problem for unconstrained optimization

We consider the problem of minimizing a $\mathcal{C}^{2}$ strongly convex function $f$ : $\mathbb{R}^{d} \rightarrow \mathbb{R}$ over $\mathbb{R}^{d}$. Instead of directly solving for $\nabla f\left(\boldsymbol{x}^{*}\right)=0$, we consider the following saddle point problem:

$$
\begin{equation*}
\inf _{\boldsymbol{x} \in \mathbb{R}^{d}} \sup _{\boldsymbol{p} \in \mathbb{R}^{d}}\langle\nabla f(\boldsymbol{x}), \boldsymbol{p}\rangle, \tag{2.3}
\end{equation*}
$$

due to the following proposition.
Proposition 2.1. Let $f: \mathbb{R}^{d} \rightarrow \mathbb{R}$ be a $\mathcal{C}^{2}$ strongly convex function. Then the saddle point to (2.3) is the unique global minimum of $f$.

Proof. Differentiating the objective function of (2.3) and setting the derivatives to 0 yield

$$
\begin{aligned}
\nabla f\left(\boldsymbol{x}^{*}\right) & =0, \\
\nabla^{2} f\left(\boldsymbol{x}^{*}\right) \boldsymbol{p}^{*} & =0 .
\end{aligned}
$$

By the strong convexity of $f$, we obtain that $\boldsymbol{x}^{*}$ is the unique global minimum and $\boldsymbol{p}^{*}=0$.

Recall that $\boldsymbol{p}^{*}=0$ by the optimality condition. Thus we make the following change to our saddle point formulation. We add a regularization term in (2.3):

$$
\begin{equation*}
\inf _{\boldsymbol{x} \in \mathbb{R}^{d}} \sup _{\boldsymbol{p} \in \mathbb{R}^{d}}\langle\nabla f(\boldsymbol{x}), \boldsymbol{p}\rangle-\frac{\varepsilon}{2}\|\boldsymbol{p}\|^{2}, \tag{2.4}
\end{equation*}
$$

where $\varepsilon>0$ is a constant. This regularization term further drives $\boldsymbol{p}$ to 0 . Similar to Proposition 2.1, we have the following proposition.

Proposition 2.2. Let $f: \mathbb{R}^{d} \rightarrow \mathbb{R}$ be a $\mathcal{C}^{2}$ strongly convex function. Then the saddle point to (2.4) is the unique global minimum of $f$.

Proof. Differentiating the objective function of (2.4) and setting derivatives to 0 yield

$$
\begin{aligned}
\nabla f\left(\boldsymbol{x}^{*}\right) & =\varepsilon \boldsymbol{p}^{*}, \\
\nabla^{2} f\left(\boldsymbol{x}^{*}\right) \boldsymbol{p}^{*} & =0 .
\end{aligned}
$$

Since $f$ is strongly convex, we have $\nabla^{2} f\left(\boldsymbol{x}^{*}\right) \succ 0$ and the second equation implies $\boldsymbol{p}^{*}=0$. Then the first equation implies $\nabla f\left(\boldsymbol{x}^{*}\right)=0$. Since $f$ is strongly convex, we conclude that $\boldsymbol{x}^{*}$ is the unique global minimum.

### 2.3. PDHG for unconstrained optimization

We apply the scheme given by (2.2) to the saddle point problem (2.4) (set $F=G=0$ and identify $K \boldsymbol{x}=\nabla f(\boldsymbol{x})$ in (2.1)). Thus,

$$
\begin{align*}
& \boldsymbol{p}^{n+1}=\underset{\boldsymbol{p}}{\arg \max }\left\langle\nabla f\left(\boldsymbol{x}^{n}\right), \boldsymbol{p}\right\rangle-\frac{\varepsilon}{2}\|\boldsymbol{p}\|^{2}-\frac{\left\|\boldsymbol{p}-\boldsymbol{p}^{n}\right\|_{\boldsymbol{A}\left(\boldsymbol{x}^{n}\right)^{-1}}^{2}}{2 \sigma},  \tag{2.5a}\\
& \widetilde{\boldsymbol{p}}^{n+1}=\omega\left(\boldsymbol{p}^{n+1}-\boldsymbol{p}^{n}\right)+\boldsymbol{p}^{n+1},  \tag{2.5b}\\
& \boldsymbol{x}^{n+1}=\underset{\boldsymbol{x}}{\arg \min }\left\langle\nabla f(\boldsymbol{x}), \widetilde{\boldsymbol{p}}^{n+1}\right\rangle+\frac{\left\|\boldsymbol{x}-\boldsymbol{x}^{n}\right\|_{\boldsymbol{B}\left(\boldsymbol{x}^{n}\right)^{-1}}^{2}}{2 \tau}, \tag{2.5c}
\end{align*}
$$

where we have added symmetric positive definite matrices $\boldsymbol{A}\left(\boldsymbol{x}^{n}\right), \boldsymbol{B}\left(\boldsymbol{x}^{n}\right) \in$ $\mathbb{R}^{d \times d}$, as preconditioners for updates of $\boldsymbol{p}, \boldsymbol{x}$, respectively. We also denote the norm $\|\boldsymbol{h}\|_{\boldsymbol{A}^{-1}}^{2}$ as $\boldsymbol{h}^{T} \boldsymbol{A}^{-1} \boldsymbol{h}$, where $\boldsymbol{h} \in \mathbb{R}^{d}$.

As mentioned, the convergence analysis of PDHG relies on the assumption that $K$ is a linear operator. So we can not apply the same convergence
analysis to (2.5) since $\nabla f(\boldsymbol{x})$ is not necessarily linear in $\boldsymbol{x}$. By taking the optimality conditions of (2.5), we find that $\boldsymbol{p}^{n+1}$ and $\boldsymbol{x}^{n+1}$ solves

$$
\begin{align*}
& \boldsymbol{p}^{n+1}+\sigma \varepsilon \boldsymbol{A}\left(\boldsymbol{x}^{n}\right) \boldsymbol{p}^{n+1}-\boldsymbol{p}^{n}-\sigma \boldsymbol{A}\left(\boldsymbol{x}^{n}\right) \nabla f\left(\boldsymbol{x}^{n}\right)=0,  \tag{2.6a}\\
& \tau \boldsymbol{B}\left(\boldsymbol{x}^{n}\right) \nabla^{2} f\left(\boldsymbol{x}^{n+1}\right)\left[(1+\omega)\left(\mathbb{I}+\sigma \varepsilon \boldsymbol{A}\left(\boldsymbol{x}^{n}\right)\right)^{-1} \sigma \boldsymbol{A}\left(\boldsymbol{x}^{n}\right) \nabla f\left(\boldsymbol{x}^{n}\right)\right. \\
& \text { b) } \left.\quad+\left((\omega+1)\left(\mathbb{I}+\sigma \varepsilon \boldsymbol{A}\left(\boldsymbol{x}^{n}\right)\right)^{-1}-\omega \mathbb{I}\right) \boldsymbol{p}^{n}\right]+\left(\boldsymbol{x}^{n+1}-\boldsymbol{x}^{n}\right)=0, \tag{2.6b}
\end{align*}
$$

where we substitute the update Eq. (2.5b) into update Eq. (2.6b). We use II to represent the identity matrix in Eq. (2.6b).

Note that the update for $\boldsymbol{x}^{n+1}$ in Eq. (2.6b) is implicit, unless $\nabla^{2} f(\boldsymbol{x})$ does not depend on $\boldsymbol{x}$. We also remark that the update for $\boldsymbol{x}^{n+1}$ in Eq. (2.6b) will be explicit if we perform a gradient step instead of a proximal step in Eq. (2.5c). To be more precise, when $\boldsymbol{B}=\mathbb{I}$, the linearized version of Eq. (2.5c) can be written as

$$
\boldsymbol{x}^{n+1}=\operatorname{prox}_{\tau\left\langle\nabla f(\cdot), \widetilde{\boldsymbol{p}}^{n+1}\right\rangle}\left(\boldsymbol{x}^{n}\right) .
$$

Taking a gradient step instead of proximal step yields

$$
\begin{equation*}
\boldsymbol{x}^{n+1}=\boldsymbol{x}^{n}-\tau \nabla^{2} f\left(\boldsymbol{x}^{n}\right) \widetilde{\boldsymbol{p}}^{n+1} . \tag{2.7}
\end{equation*}
$$

For general choice of preconditioner $\boldsymbol{B}\left(\boldsymbol{x}^{n}\right)$, the linearized version of (2.5c) satisfies

$$
\boldsymbol{x}^{n+1}=\boldsymbol{x}^{n}-\tau \boldsymbol{B}\left(x^{n}\right) \nabla^{2} f\left(\boldsymbol{x}^{n}\right) \widetilde{\boldsymbol{p}}^{n+1}=\boldsymbol{x}^{n}-\tau \boldsymbol{C}\left(x^{n}\right) \widetilde{\boldsymbol{p}}^{n+1} .
$$

Here we always denote a matrix function $\boldsymbol{C}$, such that

$$
\boldsymbol{C}\left(\boldsymbol{x}^{n}\right):=\boldsymbol{B}\left(\boldsymbol{x}^{n}\right) \nabla^{2} f\left(\boldsymbol{x}^{n}\right) .
$$

For simplicity of presentation, we only consider the simple case $\boldsymbol{A}\left(\boldsymbol{x}^{n}\right)=A \mathbb{I}$ for some $A>0$. We now summarize the linearized update Eq. (2.6) into the following algorithm.

We note that Algorithm 1 and update Eq. (2.6) are different methods for solving saddle point problem Eq. (2.3). In this paper, we focus on the computation and analysis of Algorithm 1.

### 2.4. PDD dynamics

An approach for analyzing optimization algorithms is by first studying the continuous limit of the algorithm using ODEs [25, 24, 1]. The advantage of

```
Algorithm 1 Linearized Primal-Dual Damping Algorithm
Require: Initial guesses \(\boldsymbol{x}^{0} \in \mathbb{R}^{d}, \boldsymbol{p}^{0} \in \mathbb{R}^{d}\); Stepsizes \(\tau>0, \sigma>0\); Parameters
    \(A>0, \varepsilon>0, \omega>0, C \succ 0, N>0\), tol \(>0\).
    while \(n \leq N\) and \(\left\|\nabla f\left(\boldsymbol{x}^{n}\right)\right\| \geq\) tol do
        \(\boldsymbol{p}^{n+1}=\frac{1}{1+\sigma \varepsilon A} \boldsymbol{p}^{n}+\frac{\sigma A}{1+\sigma \varepsilon A} \nabla f\left(\boldsymbol{x}^{n}\right) ;\)
        \(\tilde{\boldsymbol{p}}^{n+1}=\boldsymbol{p}^{n+1}+\omega\left(\boldsymbol{p}^{n+1}-\boldsymbol{p}^{n}\right) ;\)
        \(\boldsymbol{x}^{n+1}=\boldsymbol{x}^{n}-\tau \boldsymbol{C}\left(\boldsymbol{x}^{n}\right) \tilde{\boldsymbol{p}}^{n+1} ;\)
        \(n=n+1 ;\)
    end while
```

doing so is that ODEs provide insights into the convergence property of the algorithm.

We first reformulate the proposed algorithm Eq. (2.6) into a first-order ODE system.

Proposition 2.3. As $\tau, \sigma \rightarrow 0$ and $\sigma \omega \rightarrow \gamma$, both updates in (2.6) and Algorithm 1 can be formulated as a discrete-time update of the following ODE system.

$$
\begin{align*}
& \dot{\boldsymbol{p}}=\boldsymbol{A}(x) \nabla f(\boldsymbol{x})-\varepsilon \boldsymbol{A}(x) \boldsymbol{p},  \tag{2.8a}\\
& \dot{\boldsymbol{x}}=-\boldsymbol{C}(\boldsymbol{x})(\boldsymbol{p}+\gamma(\boldsymbol{A}(x) \nabla f(\boldsymbol{x})-\varepsilon \boldsymbol{A}(x) \boldsymbol{p})), \tag{2.8b}
\end{align*}
$$

where $\boldsymbol{C}(\boldsymbol{x})=\boldsymbol{B}(\boldsymbol{x}) \nabla^{2} f(\boldsymbol{x})$ and the initial condition satisfies $\boldsymbol{x}(0)=\boldsymbol{x}^{0}$, $\boldsymbol{p}(0)=\boldsymbol{p}^{0}$. Suppose that $\nabla f$ is Lispchitz continuous and each index in matrix $\boldsymbol{A}, \boldsymbol{C}$ is continuous and bounded. Then, there exists a unique solution for the ODE system Eq. (2.8). A stationary state ( $\boldsymbol{x}^{*}, \boldsymbol{p}^{*}$ ) of ODE system Eq. (2.8) satisfies

$$
\nabla f\left(\boldsymbol{x}^{*}\right)=0, \quad \boldsymbol{p}^{*}=0 .
$$

Proof. Rearranging Eq. (2.6a) and Eq. (2.6b), we have

$$
\begin{aligned}
\frac{\boldsymbol{p}^{n+1}-\boldsymbol{p}^{n}}{\sigma}= & \boldsymbol{A}\left(\boldsymbol{x}^{n}\right) \nabla f\left(\boldsymbol{x}^{n}\right)-\varepsilon \boldsymbol{A}\left(\boldsymbol{x}^{n}\right) \boldsymbol{p}^{n+1} \\
\frac{\boldsymbol{x}^{n+1}-\boldsymbol{x}^{n}}{\tau}=- & \boldsymbol{B}\left(\boldsymbol{x}^{n}\right) \nabla^{2} f\left(\boldsymbol{x}^{n+1}\right)\left[(1+\omega)\left(\mathbb{I}+\sigma \varepsilon \boldsymbol{A}\left(\boldsymbol{x}^{n}\right)\right)^{-1} \sigma \boldsymbol{A}\left(\boldsymbol{x}^{n}\right) \nabla f\left(\boldsymbol{x}^{n}\right)\right. \\
& \left.+\left((\omega+1)\left(\mathbb{I}+\sigma \varepsilon \boldsymbol{A}\left(\boldsymbol{x}^{n}\right)\right)^{-1}-\omega \mathbb{I}\right) \boldsymbol{p}^{n}\right] .
\end{aligned}
$$

Taking the limit as $\tau, \sigma \rightarrow 0$ and $\sigma \omega \rightarrow \gamma$, we obtain

$$
\dot{\boldsymbol{p}}=\boldsymbol{A}(x) \nabla f(\boldsymbol{x})-\varepsilon \boldsymbol{A}(x) \boldsymbol{p},
$$

$$
\dot{\boldsymbol{x}}=-\boldsymbol{B}(\boldsymbol{x}) \nabla^{2} f(\boldsymbol{x})(\boldsymbol{p}+\gamma(\boldsymbol{A}(x) \nabla f(\boldsymbol{x})-\varepsilon \boldsymbol{A}(x) \boldsymbol{p})) .
$$

Similarly, the update in Algorithm 1 also converges to the ODE system Eq. (2.8). Clearly, a stationary state satisfies $\boldsymbol{p}^{*}=0, \nabla f\left(\boldsymbol{x}^{*}\right)=0$.

Proposition 2.4 (Primal-dual damping second order ODE). The ODE system Eq. (2.8) satisfies the following second-order ODE

$$
\begin{equation*}
\ddot{\boldsymbol{x}}+\left[\varepsilon \boldsymbol{A}+\gamma \boldsymbol{C} \boldsymbol{A} \nabla^{2} f(\boldsymbol{x})-\dot{\boldsymbol{C}} \boldsymbol{C}^{-1}\right] \dot{\boldsymbol{x}}+\boldsymbol{C} \boldsymbol{A} \nabla f(\boldsymbol{x})=0 . \tag{2.9}
\end{equation*}
$$

Here $\dot{\boldsymbol{C}}=\frac{d}{d t} \boldsymbol{C}(\boldsymbol{x}(t))$.
The proof follows by direct calculations and can be found in Appendix C. We note that the formulation given by Eq. (2.9) includes several important special cases in the literature. In a word, we view Eq. (2.4) as a preconditioned accelerated gradient flow.

Example 2.1. Let $\boldsymbol{C}=\boldsymbol{A}=\mathbb{I}$ and $\gamma \neq 0$. Then equation Eq. (2.4) satisfies

$$
\begin{equation*}
\ddot{\boldsymbol{x}}+\epsilon \dot{\boldsymbol{x}}+\gamma \nabla^{2} f(\boldsymbol{x}) \dot{\boldsymbol{x}}+\nabla f(\boldsymbol{x})=0, \tag{2.10}
\end{equation*}
$$

which is an inertial system with Hessian-driven damping [2].
Remark 2.5. In the case of $\boldsymbol{C}=\boldsymbol{A}=\mathbb{I}$, although the derived second order ODE Eq. (2.9) is the same as the one in Attouch et al. [2] at a continuous time level, our algorithm 1 provides a different time discretization from the one in Attouch et al. [2].

Example 2.2. Let $\boldsymbol{C}=\boldsymbol{A}=\mathbb{I}, \gamma(t)=0$. Then equation Eq. (2.4) satisfies the heavy ball ODE [24]

$$
\begin{equation*}
\ddot{\boldsymbol{x}}+\varepsilon \dot{\boldsymbol{x}}+\nabla f(\boldsymbol{x})=0 . \tag{2.11}
\end{equation*}
$$

Example 2.3. Let $\boldsymbol{C}=\boldsymbol{A}=\mathbb{I}, \gamma(t)=0, \varepsilon(t)=\frac{3}{t}$. Then equation Eq. (2.4) satisfies the Nesterov ODE [25]:

$$
\begin{equation*}
\ddot{\boldsymbol{x}}+\frac{3}{t} \dot{\boldsymbol{x}}+\nabla f(\boldsymbol{x})=0 . \tag{2.12}
\end{equation*}
$$

We next provide a convergence analysis of ODE Eq. (2.8) for quadratic optimization problems. We demonstrate the importance of preconditioners in characterizing the convergence speed of ODE Eq. (2.8).

Theorem 2.6. Suppose $f(\boldsymbol{x})=\frac{1}{2} \boldsymbol{x}^{T} \boldsymbol{Q} \boldsymbol{x}$ for some symmetric positive definite matrix $\boldsymbol{Q} \in \mathbb{R}^{d \times d}$. Assume $\boldsymbol{A}, \boldsymbol{B}$ are constant matrices. In this case, equation Eq. (2.8) satisfies the linear ODE system:

$$
\binom{\dot{\boldsymbol{x}}}{\dot{\boldsymbol{p}}}=\left(\begin{array}{cc}
-\gamma \boldsymbol{B} \boldsymbol{Q} \boldsymbol{A} \boldsymbol{Q} & -\boldsymbol{B} \boldsymbol{Q}(\mathbb{I}-\gamma \varepsilon \boldsymbol{A}) \\
\boldsymbol{A} \boldsymbol{Q} & -\varepsilon \boldsymbol{A}
\end{array}\right)\binom{\boldsymbol{x}}{\boldsymbol{p}} .
$$

Suppose that $\boldsymbol{A}$ commutes with $\boldsymbol{Q}$, such that $\boldsymbol{A} \boldsymbol{Q}=\boldsymbol{Q} \boldsymbol{A}$. Suppose $\boldsymbol{A}$ and $\boldsymbol{B Q A Q}$ are simultaneously diagonalizable and have positive eigenvalues. Let $\mu_{1} \geq \cdots \geq \mu_{n}>0$ be the eigenvalues of $\boldsymbol{B} \boldsymbol{Q} \boldsymbol{A} \boldsymbol{Q}$ and $a_{i}$ the $i$-th eigenvalue of $\boldsymbol{A}$ (not necessarily in descending order) in the same basis. Then
(a) The solution of ODE system (2.8) converges to $\left(\boldsymbol{x}^{*}, \boldsymbol{p}^{*}\right)=(0,0)$ :

$$
\|(\boldsymbol{x}(t), \boldsymbol{p}(t))\| \leq\left\|\left(\boldsymbol{x}_{0}, \boldsymbol{p}_{0}\right)\right\| \exp (\alpha t),
$$

where

$$
\alpha=\max _{i} \frac{1}{2}\left[-\gamma \mu_{i}-\varepsilon a_{i}+\Re\left(\sqrt{\left(\gamma \mu_{i}+\varepsilon\right)^{2}-4 \mu_{i}}\right)\right] .
$$

(b) When $\boldsymbol{A}=\mathbb{I}, \varepsilon=0$, the optimal convergence rate is achieved at $\gamma^{*}=\frac{2 \sqrt{\mu_{1}}}{\sqrt{\mu_{n}\left(2 \mu_{1}-\mu_{n}\right)}}$. The corresponding rate is $\alpha=\frac{-\sqrt{\mu_{n}}}{\sqrt{2-\frac{1}{\kappa}}}$, where $\kappa=$ $\mu_{1} / \mu_{n}>1$.
(c) Moreover, when $\gamma=\varepsilon=0$, the system will not converge for any initial data $\left(\boldsymbol{x}_{0}, \boldsymbol{p}_{0}\right) \neq(0,0)$.
(d) If $\boldsymbol{A}=\mathbb{I}, \gamma \leq \frac{1}{\sqrt{\mu_{1}}}, \varepsilon=2 \sqrt{\mu^{\prime}}-\gamma \mu^{\prime}$ for some $\mu^{\prime} \leq \mu_{n}$, then

$$
\alpha=-\sqrt{\mu^{\prime}}-\frac{\gamma}{2}\left(\mu_{n}-\mu^{\prime}\right) \leq-\sqrt{\mu^{\prime}} .
$$

We defer the proof to Appendix B.
Remark 2.7. If $\omega$ is bounded, then we have $\gamma=\mathcal{O}(\sigma)$. Then, in the limit as $\sigma \rightarrow 0$, we also have that $\gamma \rightarrow 0$. By Theorem 2.6 (c), the ODE system (2.8) does not converge for any initial data.
Remark 2.8. If $\mu^{\prime}$ is an estimate of the smallest eigenvalue $\mu_{n}$, then the convergence speed for the solution of heavy ball ODE is $\exp \left(-\sqrt{\mu^{\prime}} t\right)$. In Theorem 2.6 (d), if $\gamma=0$ and $\mu^{\prime}=\mu_{n}$, then $\alpha=-\sqrt{\mu_{n}}$ which is the same as the convergence rate of the heavy ball ODE [24]. However, if $\gamma>0$ and $\mu^{\prime}<\mu_{n}$, then we have $\alpha=-\sqrt{\mu^{\prime}}-\gamma\left(\mu_{n}-\mu^{\prime}\right)<-\sqrt{\mu^{\prime}}$, which converges faster than the heavy ball ODE.

## 3. Lyapunov Analysis

In this section, we present the main theoretical result of this paper. We provide the convergence analysis for general objective functions in both continuous-time ODEs Eq. (2.8) and discrete-time Algorithm 1. From now on, we make the following two assumptions for the convergence analysis.
Assumption 3.1. There exists two constants $L \geq \mu>0$ such that $\mu \mathbb{I} \preceq$ $\boldsymbol{C}_{0}(\boldsymbol{x}) \preceq L \mathbb{I}$ for all $\boldsymbol{x}$, where $\boldsymbol{C}_{0}(\boldsymbol{x})=\nabla^{2} f(\boldsymbol{x}) \boldsymbol{B}(\boldsymbol{x}) \nabla^{2} f(\boldsymbol{x})$, and $\mu \leq 1$.

Assumption 3.2. There exists a constant $L^{\prime}>0$ such that

$$
\begin{equation*}
\boldsymbol{C}(\boldsymbol{x})^{T}\left(\nabla^{3} f(\boldsymbol{x}) \nabla f(\boldsymbol{x})+\left(\nabla^{2} f(\boldsymbol{x})\right)^{2}\right) \boldsymbol{C}(\boldsymbol{x}) \preceq L^{\prime} \mathbb{I} \tag{3.1}
\end{equation*}
$$

for all $\boldsymbol{x}$, where $\boldsymbol{C}(\boldsymbol{x})=\boldsymbol{B}(\boldsymbol{x}) \nabla^{2} f(\boldsymbol{x})$.

### 3.1. Continuous time Lyapunov analysis

In this subsection, we establish convergence results of the ODE system Eq. (2.8).

Theorem 3.3. Consider the ODE system Eq. (2.8) with an initial condition $(x(0), p(0)) \in \mathbb{R}^{2 d}$. Define the functional

$$
\begin{equation*}
\mathcal{I}(\boldsymbol{x}, \boldsymbol{p})=\frac{1}{2}\left(\|\boldsymbol{p}\|^{2}+\|\nabla f(\boldsymbol{x})\|^{2}\right) \tag{3.2}
\end{equation*}
$$

Suppose Assumption 3.1 holds, we have

$$
\begin{equation*}
\mathcal{I}(\boldsymbol{x}(t), \boldsymbol{p}(t)) \leq \mathcal{I}(\boldsymbol{x}(0), \boldsymbol{p}(0)) \exp (-2 \lambda t) \tag{3.3}
\end{equation*}
$$

where

$$
\begin{aligned}
& \lambda=\min \left\{\mu \gamma A-\frac{1}{2}|A-\mu(1-\varepsilon \gamma A)|, L \gamma A-\frac{1}{2}|A-L(1-\varepsilon \gamma A)|\right. \\
&\left.\varepsilon A-\frac{1}{2}|A-\mu(1-\varepsilon \gamma A)|, \varepsilon A-\frac{1}{2}|A-L(1-\varepsilon \gamma A)|\right\}
\end{aligned}
$$

In particular, when $\gamma=\frac{1}{\mu}, \varepsilon=1, A=\frac{\mu+L}{2+(\mu+L) \varepsilon \gamma}$, then $\lambda=\frac{\mu}{2}$.
Proof. It is straightforward to compute the following

$$
\frac{\mathrm{d} \mathcal{I}}{\mathrm{~d} t}=\langle\boldsymbol{p}, \dot{\boldsymbol{p}}\rangle+\left\langle\nabla f, \nabla^{2} f \dot{\boldsymbol{x}}\right\rangle
$$

$$
\begin{equation*}
=-\nabla f^{T} \boldsymbol{C}_{0} \gamma \boldsymbol{A} \nabla f-\boldsymbol{p}^{T} \varepsilon \boldsymbol{A} \boldsymbol{p}+\nabla f^{T}\left(\boldsymbol{A}-\boldsymbol{C}_{0}(\mathbb{I}-\varepsilon \gamma \boldsymbol{A})\right) \boldsymbol{p} \tag{3.4}
\end{equation*}
$$

We shall find $\lambda$ such that $\frac{\mathrm{d} \mathcal{I}}{\mathrm{d} t}+2 \lambda \mathcal{I} \leq 0$. Then we obtain the exponential convergence by Gronwall's inequality, i.e.,

$$
\mathcal{I}(\boldsymbol{x}(t), \boldsymbol{p}(t)) \leq \mathcal{I}(\boldsymbol{x}(0), \boldsymbol{p}(0)) \exp (-2 \lambda t)
$$

We can compute

$$
\begin{align*}
\frac{\mathrm{d} \mathcal{I}}{\mathrm{~d} t}+2 \lambda \mathcal{I}=\nabla & f^{T}\left(-\boldsymbol{C}_{0} \gamma \boldsymbol{A}+\lambda \mathbb{I}\right) \nabla f+\boldsymbol{p}^{T}(-\varepsilon \boldsymbol{A}+\lambda \mathbb{I}) \boldsymbol{p} \\
& +\nabla f^{T}\left(\boldsymbol{A}-\boldsymbol{C}_{0}(\mathbb{I}-\varepsilon \gamma \boldsymbol{A})\right) \boldsymbol{p} \tag{3.5}
\end{align*}
$$

By Lemma A.1, we obtain the following sufficient conditions for $\frac{\mathrm{d} \mathcal{I}}{\mathrm{d} t}+2 \lambda \mathcal{I} \leq 0$

$$
\begin{align*}
& -\varepsilon A+\lambda+\frac{1}{2}\left|\xi_{i}(1-\varepsilon \gamma A)-A\right| \leq 0  \tag{3.6a}\\
& \lambda-\xi_{i} \gamma A+\frac{1}{2}\left|\xi_{i}(1-\varepsilon \gamma A)-A\right| \leq 0 \tag{3.6b}
\end{align*}
$$

where $\xi_{i}(\boldsymbol{x})$ is the eigenvalue of $\boldsymbol{C}_{0}(\boldsymbol{x})$. By our assumptions, we have $L \geq$ $\xi_{1}(\boldsymbol{x}) \geq \cdots \geq \xi_{n}(\boldsymbol{x}) \geq \mu$. Eq. (3.6) give two upper bounds on $\lambda$. Define $g_{1}(\xi)=\varepsilon A-\frac{1}{2}|\xi(1-\varepsilon \gamma A)-A|$, and $g_{2}(\xi)=\xi \gamma A-\frac{1}{2}|\xi(1-\varepsilon \gamma A)-A|$ on the interval $[\mu, L]$. Then Eq. (3.6) implies that

$$
\begin{equation*}
\lambda \leq g_{j}\left(\xi_{i}\right) \tag{3.7}
\end{equation*}
$$

for all $i=1, \ldots, n$ and $j=1,2$. Since each $g_{j}(\xi)$ is a piece-wise linear in $\xi$, it is not hard to see that

$$
\min _{\xi \in[\mu, L]} g_{j}(\xi)=\min \left\{g_{j}(\mu), g_{j}(L)\right\}
$$

for $j=1,2$. This proves the formula for $\lambda$. When $A=\frac{\mu+L}{2+(\mu+L) \varepsilon \gamma}$, we have $g_{1}(\mu)=g_{1}(L)$, and

$$
\mu(1-\varepsilon \gamma A)-A=-L(1-\varepsilon \gamma A)+A
$$

Further, requiring $g_{1}(\mu)=g_{2}(\mu)$ yields $\varepsilon=\mu \gamma$. And we obtain

$$
\lambda=\mu \gamma A-\frac{1}{2}|A-\mu(1-\varepsilon \gamma \boldsymbol{A})|
$$

$$
\begin{align*}
& =\mu \gamma A-\frac{1}{2}(A-\mu(1-\varepsilon \gamma A)) \\
& =\frac{\mu}{2}+A\left(\gamma \mu-\frac{1}{2} \gamma^{2} \mu^{2}-\frac{1}{2}\right) \\
& =\frac{\mu}{2}-\frac{A}{2}(\gamma \mu-1)^{2} . \tag{3.8}
\end{align*}
$$

We note that $\lambda$ is maximized when taking $\gamma=\mu^{-1}$. We obtain $\lambda=\frac{\mu}{2}$.

### 3.2. Discrete time Lyapunov analysis

In this subsection, we study the convergence criterion for the discretized linearized PDHG flow given by Eq. (1.2) and Algorithm 1. From now on, we assume that $f$ is a $\mathcal{C}^{4}$ strongly convex function. We can rewrite the iterations as

$$
\begin{align*}
\boldsymbol{p}^{n+1} & =\frac{1}{1+\sigma \varepsilon A} \boldsymbol{p}^{n}+\frac{\sigma A}{1+\sigma \varepsilon A} \nabla f\left(\boldsymbol{x}^{n}\right)  \tag{3.9a}\\
\boldsymbol{x}^{n+1} & =\boldsymbol{x}^{n}-\tau \boldsymbol{B}\left(\boldsymbol{x}^{n}\right) \nabla^{2} f\left(\boldsymbol{x}^{n}\right)\left(\frac{1-\varepsilon \gamma A}{1+\sigma \varepsilon A} \boldsymbol{p}^{n}+\frac{\sigma A+\gamma A}{1+\sigma \varepsilon A} \nabla f\left(\boldsymbol{x}^{n}\right)\right) \tag{3.9b}
\end{align*}
$$

where $\gamma=\sigma \omega$. We define the following notations which will be used later.

$$
\begin{aligned}
& \boldsymbol{N}\left(\boldsymbol{x}^{n}\right) \\
& =\frac{1}{1+\sigma \varepsilon A}\left(\begin{array}{cc}
\boldsymbol{B}\left(\boldsymbol{x}^{n}\right) \nabla^{2} f\left(\boldsymbol{x}^{n}\right)(\sigma A+\gamma A) & \boldsymbol{B}\left(\boldsymbol{x}^{n}\right) \nabla^{2} f\left(\boldsymbol{x}^{n}\right)(1-\varepsilon \gamma A) \\
-\frac{\sigma}{\tau} A & \frac{\sigma}{\tau} \varepsilon A
\end{array}\right) .
\end{aligned}
$$

And

$$
\boldsymbol{H}\left(\boldsymbol{x}^{n}\right)=\operatorname{sym}\left(\left(\begin{array}{cc}
\nabla^{2} f\left(\boldsymbol{x}^{n}\right) & 0 \\
0 & \mathbb{I}
\end{array}\right) \cdot \boldsymbol{N}\left(\boldsymbol{x}^{n}\right)\right) .
$$

Remark 3.4. The matrix $\boldsymbol{N}\left(\boldsymbol{x}^{n}\right)$ and $\boldsymbol{H}\left(\boldsymbol{x}^{n}\right)$ also depends on the $\tau, \sigma, A, \varepsilon$ and $\omega$.

Define the Lyapunov functional in discrete time as

$$
\mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)=\frac{1}{2}\left\|\nabla f\left(\boldsymbol{x}^{n}\right)\right\|^{2}+\frac{1}{2}\left\|\boldsymbol{p}^{n}\right\|^{2} .
$$

Theorem 3.5. Suppose that there exists positive constants $\lambda, M_{1} \in \mathbb{R}_{+}$, such that

$$
\boldsymbol{H}(\boldsymbol{x}) \succeq \lambda \mathbb{I}
$$

$$
\boldsymbol{N}(\boldsymbol{x})^{T} \nabla^{2} \mathcal{I}(\tilde{\boldsymbol{x}}, \tilde{\boldsymbol{p}}) \boldsymbol{N}(\boldsymbol{x}) \preceq M_{1} \mathbb{I},
$$

for all $\boldsymbol{x}, \tilde{\boldsymbol{x}} \in \mathbb{R}^{n}$. If $\tau=a \frac{\lambda}{M}$ for some $a \in(0,2)$, then the functional $\mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)$ decreases geometrically, i.e.

$$
\mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right) \leq \mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right)\left(1+\left(a^{2}-2 a\right) \frac{\lambda^{2}}{2 M_{1}}\right)^{n} .
$$

Proof. It follows from our definition of $\boldsymbol{N}\left(\boldsymbol{x}^{n}\right)$ that

$$
\begin{equation*}
\binom{\boldsymbol{x}^{n+1}-\boldsymbol{x}^{n}}{\boldsymbol{p}^{n+1}-\boldsymbol{p}^{n}}=-\tau \boldsymbol{N}\left(\boldsymbol{x}^{n}\right)\binom{\nabla f\left(\boldsymbol{x}^{n}\right)}{\boldsymbol{p}^{n}}, \tag{3.10}
\end{equation*}
$$

By the mean-value theorem, we obtain

$$
\begin{aligned}
\mathcal{I}\left(\boldsymbol{x}^{n+1}, \boldsymbol{p}^{n+1}\right)-\mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right) & =\binom{\nabla_{\boldsymbol{x}} \mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)}{\nabla_{\boldsymbol{p}} \mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)}^{T}\binom{\boldsymbol{x}^{n+1}-\boldsymbol{x}^{n}}{\boldsymbol{p}^{n+1}-\boldsymbol{p}^{n}} \\
& +\frac{1}{2}\binom{\boldsymbol{x}^{n+1}-\boldsymbol{x}^{n}}{\boldsymbol{p}^{n+1}-\boldsymbol{p}^{n}}^{T} \nabla^{2} \mathcal{I}(\tilde{\boldsymbol{x}}, \tilde{\boldsymbol{p}})\binom{\boldsymbol{x}^{n+1}-\boldsymbol{x}^{n}}{\boldsymbol{p}^{n+1}-\boldsymbol{p}^{n}},
\end{aligned}
$$

where $(\tilde{\boldsymbol{x}}, \tilde{\boldsymbol{p}})$ is in between $\left(\boldsymbol{x}^{n+1}, \boldsymbol{p}^{n+1}\right)$ and $\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)$. And

$$
\begin{aligned}
& \nabla_{\boldsymbol{x}} \mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)=\nabla^{2} f\left(\boldsymbol{x}^{n}\right) \nabla f\left(\boldsymbol{x}^{n}\right), \\
& \nabla_{\boldsymbol{p}} \mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)=\boldsymbol{p}^{n}, \\
& \nabla^{2} \mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)=\left(\begin{array}{cc}
\nabla^{3} f\left(\boldsymbol{x}^{n}\right) \nabla f\left(\boldsymbol{x}^{n}\right)+\nabla^{2} f\left(\boldsymbol{x}^{n}\right) \nabla^{2} f\left(\boldsymbol{x}^{n}\right) & 0 \\
0 & \mathbb{I}
\end{array}\right) .
\end{aligned}
$$

Then using Eq. (3.10) and definition of $\boldsymbol{H}\left(\boldsymbol{x}^{n}\right)$, we obtain

$$
\begin{align*}
& \mathcal{I}\left(\boldsymbol{x}^{n+1}, \boldsymbol{p}^{n+1}\right)-\mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right) \\
= & -\tau\binom{\nabla f\left(\boldsymbol{x}^{n}\right)}{\boldsymbol{p}^{n}}^{T}\left(\begin{array}{cc}
\nabla^{2} f\left(\boldsymbol{x}^{n}\right) & 0 \\
0 & \mathbb{I}
\end{array}\right) \cdot \boldsymbol{N}\left(\boldsymbol{x}^{n}\right)\binom{\nabla f\left(\boldsymbol{x}^{n}\right)}{\boldsymbol{p}^{n}} \\
& +\frac{\tau^{2}}{2}\binom{\nabla f\left(\boldsymbol{x}^{n}\right)}{\boldsymbol{p}^{n}}^{T} \boldsymbol{N}\left(\boldsymbol{x}^{n}\right)^{T} \nabla^{2} \mathcal{I}(\tilde{\boldsymbol{x}}, \tilde{\boldsymbol{p}}) \boldsymbol{N}\left(\boldsymbol{x}^{n}\right)\binom{\nabla f\left(\boldsymbol{x}^{n}\right)}{\boldsymbol{p}^{n}} \\
= & -\tau\binom{\nabla f\left(\boldsymbol{x}^{n}\right)}{\boldsymbol{p}^{n}}^{T} \boldsymbol{H}\left(\boldsymbol{x}^{n}\right)\binom{\nabla f\left(\boldsymbol{x}^{n}\right)}{\boldsymbol{p}^{n}} \\
& +\frac{\tau^{2}}{2}\binom{\nabla f\left(\boldsymbol{x}^{n}\right)}{\boldsymbol{p}^{n}}^{T} \boldsymbol{N}\left(\boldsymbol{x}^{n}\right)^{T} \nabla^{2} \mathcal{I}(\tilde{\boldsymbol{x}}, \tilde{\boldsymbol{p}}) \boldsymbol{N}\left(\boldsymbol{x}^{n}\right)\binom{\nabla f\left(\boldsymbol{x}^{n}\right)}{\boldsymbol{p}^{n}}, \tag{3.11}
\end{align*}
$$

From Eq. (3.11) and our assumption on $\boldsymbol{N}(\boldsymbol{x})$ and $\boldsymbol{H}(\boldsymbol{x})$, we obtain

$$
\begin{align*}
\mathcal{I}\left(\boldsymbol{x}^{n+1}, \boldsymbol{p}^{n+1}\right)-\mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right) & \leq\left(-\tau \lambda+\frac{\tau^{2} M_{1}}{2}\right) \mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right) \\
& =\frac{M_{1}}{2}\left(\left(\tau-\frac{\lambda}{M_{1}}\right)^{2}-\frac{\lambda^{2}}{M_{1}^{2}}\right) \mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right) \\
& =\left(a^{2}-2 a\right) \frac{\lambda^{2}}{2 M_{1}} \mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right), \tag{3.12}
\end{align*}
$$

where we used $\tau=a \frac{\lambda}{M_{1}}$. Hence,

$$
\begin{aligned}
\mathcal{I}\left(\boldsymbol{x}^{n+1}, \boldsymbol{p}^{n+1}\right) & \leq \mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)\left(1+\left(a^{2}-2 a\right) \frac{\lambda^{2}}{2 M_{1}}\right) \\
& \leq \mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right)\left(1+\left(a^{2}-2 a\right) \frac{\lambda^{2}}{2 M_{1}}\right)^{n+1}
\end{aligned}
$$

When $0<a<2$, we have $a^{2}-2 a<0$. Thus we obtain the desired convergence result.
Theorem 3.6. Let $f: \mathbb{R}^{d} \rightarrow \mathbb{R}$ be a $\mathcal{C}^{4}$ strongly convex function. Suppose $\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right)$ satisfies

$$
\begin{equation*}
\mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right)^{1 / 2} \leq \frac{\delta}{\tau D_{0}\|\boldsymbol{N}(\boldsymbol{x})\|_{2}^{3}}, \tag{3.13}
\end{equation*}
$$

for some $\delta>0$ and all $\boldsymbol{x}$. Here

$$
D_{0}=\sup _{\boldsymbol{x}, \boldsymbol{p}, \boldsymbol{x}^{\prime}, \boldsymbol{p}^{\prime}} \frac{\binom{\boldsymbol{x}^{\prime}}{\boldsymbol{p}^{\prime}}^{T}\left(\nabla^{3} \mathcal{I}(\boldsymbol{x}, \boldsymbol{p})\binom{\boldsymbol{x}^{\prime}}{\boldsymbol{p}^{\prime}}\right)\binom{\boldsymbol{x}^{\prime}}{\boldsymbol{p}^{\prime}}}{\left\|\binom{\boldsymbol{x}^{\prime}}{\boldsymbol{p}^{\prime}}\right\|_{2}^{3}}
$$

Suppose further that there exists positive constants $\lambda, M_{2} \in \mathbb{R}_{+}$such that

$$
\begin{aligned}
\boldsymbol{H}(\boldsymbol{x}) & \succeq \lambda \mathbb{I}, \\
\boldsymbol{N}(\boldsymbol{x})^{T} \nabla^{2} \mathcal{I}(\boldsymbol{x}, \boldsymbol{p}) \boldsymbol{N}(\boldsymbol{x}) & \preceq M_{2} \mathbb{I}
\end{aligned}
$$

for all $\boldsymbol{x} \in \mathbb{R}^{n}$. If $\tau=a \frac{\lambda}{M_{2}+\delta}$ for some $a \in(0,2)$, then the functional $\mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)$ decreases geometrically, i.e.

$$
\mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right) \leq \mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right)\left(1+\frac{a^{2}-2 a}{2} \frac{\lambda^{2}}{M_{2}+\delta}\right)^{n}
$$

Remark 3.7. Note that the constant $M_{2}$ in Theorem 3.6 can be better than the constant $M_{1}$ in Theorem 3.5 because $\boldsymbol{N}$ and $\nabla^{2} \mathcal{I}$ are evaluated at the same $\boldsymbol{x}$ in Theorem 3.6.

Proof. We will prove it by induction. Using the mean-value theorem, we have

$$
\begin{align*}
& \mathcal{I}\left(\boldsymbol{x}^{n+1}, \boldsymbol{p}^{n+1}\right)-\mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)=\binom{\nabla_{\boldsymbol{x}} \mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)}{\nabla_{\boldsymbol{p}} \mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)}^{T}\binom{\boldsymbol{x}^{n+1}-\boldsymbol{x}^{n}}{\boldsymbol{p}^{n+1}-\boldsymbol{p}^{n}} \\
& \quad+\frac{1}{2}\binom{\boldsymbol{x}^{n+1}-\boldsymbol{x}^{n}}{\boldsymbol{p}^{n+1}-\boldsymbol{p}^{n}}^{T} \nabla^{2} \mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)\binom{\boldsymbol{x}^{n+1}-\boldsymbol{x}^{n}}{\boldsymbol{p}^{n+1}-\boldsymbol{p}^{n}} \\
& (3.14) \quad+\frac{1}{6}\binom{\boldsymbol{x}^{n+1}-\boldsymbol{x}^{n}}{\boldsymbol{p}^{n+1}-\boldsymbol{p}^{n}}^{T}\left(\nabla^{3} \mathcal{I}\left(\tilde{\boldsymbol{x}}^{n}, \tilde{\boldsymbol{p}}^{n}\right)\binom{\boldsymbol{x}^{n+1}-\boldsymbol{x}^{n}}{\boldsymbol{p}^{n+1}-\boldsymbol{p}^{n}}\right)\binom{\boldsymbol{x}^{n+1}-\boldsymbol{x}^{n}}{\boldsymbol{p}^{n+1}-\boldsymbol{p}^{n}}, \tag{3.14}
\end{align*}
$$

where ( $\tilde{\boldsymbol{x}}^{n}, \tilde{\boldsymbol{p}}^{n}$ ) is in between $\left(\boldsymbol{x}^{n+1}, \boldsymbol{p}^{n+1}\right)$ and ( $\boldsymbol{x}^{n}, \boldsymbol{p}^{n}$ ). By Eq. (3.14) and Eq. (3.10), we can bound

$$
\begin{aligned}
& \mathcal{I}\left(\boldsymbol{x}^{1}, \boldsymbol{p}^{1}\right)-\mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right) \\
&=-\tau\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}}^{T} \boldsymbol{H}\left(\boldsymbol{x}^{0}\right)\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}} \\
&+\frac{\tau^{2}}{2}\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}}^{T} \boldsymbol{N}\left(\boldsymbol{x}^{0}\right)^{T} \nabla^{2} \mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right) \boldsymbol{N}\left(\boldsymbol{x}^{0}\right)\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}} \\
&-\frac{\tau^{3}}{6}\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}}^{T} \boldsymbol{N}\left(\boldsymbol{x}^{0}\right)^{T}\left(\nabla^{3} \mathcal{I}\left(\tilde{\boldsymbol{x}}^{0}, \tilde{\boldsymbol{p}}^{0}\right) \boldsymbol{N}\left(\boldsymbol{x}^{0}\right)\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}}\right) \\
& \cdot \boldsymbol{N}\left(\boldsymbol{x}^{0}\right)\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}} \\
& \leq-\tau\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}}^{T} \boldsymbol{H}\left(\boldsymbol{x}^{0}\right)\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}} \\
&+\frac{\tau^{2}}{2}\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}}^{T} \boldsymbol{N}\left(\boldsymbol{x}^{0}\right)^{T} \nabla^{2} \mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right) \boldsymbol{N}\left(\boldsymbol{x}^{0}\right)\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}} \\
&=-\tau\binom{\nabla\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}}^{T}\left(\begin{array}{c}
D_{0}\left\|\boldsymbol{N}\left(\boldsymbol{x}^{0}\right)\right\|_{2}^{3}
\end{array} \|\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}}\right)^{T} \boldsymbol{H}\left(\boldsymbol{x}^{0}\right)\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}}\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}} \\
&+\frac{\tau^{2}\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}}^{T} \boldsymbol{N}\left(\boldsymbol{x}^{0}\right)^{T} \nabla^{2} \mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right) \boldsymbol{N}\left(\boldsymbol{x}^{0}\right)\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}}}{}
\end{aligned}
$$

$$
\begin{align*}
& +\frac{\tau^{3}}{6}\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}}^{T} D_{0}\left\|\boldsymbol{N}\left(\boldsymbol{x}^{0}\right)\right\|_{2}^{3} \mathcal{I}\left(\boldsymbol{x}_{0}, \boldsymbol{p}_{0}\right)^{1 / 2}\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}} \\
\leq & -\tau\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}}^{T} \boldsymbol{H}\left(\boldsymbol{x}^{0}\right)\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}} \\
& +\frac{\tau^{2}}{2}\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}}^{T} \boldsymbol{N}\left(\boldsymbol{x}^{0}\right)^{T} \nabla^{2} \mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right) \boldsymbol{N}\left(\boldsymbol{x}^{0}\right)\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}} \\
& +\frac{\tau^{2} \delta}{6}\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}}^{T}\binom{\nabla f\left(\boldsymbol{x}^{0}\right)}{\boldsymbol{p}^{0}}, \tag{3.15}
\end{align*}
$$

where the last inequality is by our assumption on $\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right)$. Using our assumptions on the lower bound of $\boldsymbol{H}$ and the upper bound of $\boldsymbol{N}^{T} \cdot \nabla^{2} \mathcal{I} \cdot \boldsymbol{N}$, we obtain

$$
\begin{align*}
\mathcal{I}\left(\boldsymbol{x}^{1}, \boldsymbol{p}^{1}\right)-\mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right) & \leq\left(-\tau \lambda+\frac{\tau^{2} \delta}{6}+\frac{\tau^{2} M_{2}}{2}\right) \mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right) \\
& \leq\left(-\tau \lambda+\frac{\tau^{2}\left(\delta+M_{2}\right)}{2}\right) \mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right) \\
& =\frac{1}{2}\left(a^{2}-2 a\right) \frac{\lambda^{2}}{M_{2}+\delta} \mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right), \tag{3.16}
\end{align*}
$$

where we used $\tau=a \frac{\lambda}{M_{2}+\delta}$ for some $a \in(0,2)$. Hence,

$$
\mathcal{I}\left(\boldsymbol{x}^{1}, \boldsymbol{p}^{1}\right) \leq \mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right)\left(1+\frac{a^{2}-2 a}{2} \frac{\lambda^{2}}{M_{2}+\delta}\right) .
$$

This proves the base case. Now suppose it holds that

$$
\mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right) \leq \mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right)\left(1+\frac{a^{2}-2 a}{2} \frac{\lambda^{2}}{M_{2}+\delta}\right)^{n},
$$

for some $n \geq 1$. In particular, this implies that

$$
\mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)<\mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right)
$$

which yields

$$
\tau D_{0}\|\boldsymbol{N}(\boldsymbol{x})\|_{2}^{3} \mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)^{1 / 2}<\tau D_{0}\|\boldsymbol{N}(\boldsymbol{x})\|_{2}^{3} \mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right)^{1 / 2} \leq \delta .
$$

Then, repeating the derivation of Eq. (3.15) and Eq. (3.16) yields

$$
\mathcal{I}\left(\boldsymbol{x}^{n+1}, \boldsymbol{p}^{n+1}\right) \leq \mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right)\left(1+\frac{a^{2}-2 a}{2} \frac{\lambda^{2}}{M_{2}+\delta}\right)
$$

Combining with our induction hypothesis, we conclude that

$$
\mathcal{I}\left(\boldsymbol{x}^{n+1}, \boldsymbol{p}^{n+1}\right) \leq \mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right)\left(1+\frac{a^{2}-2 a}{2} \frac{\lambda^{2}}{M_{2}+\delta}\right)^{n+1} .
$$

The proof is complete by induction.
Corollary 3.8. Suppose Assumption 3.1 and Assumption 3.2 hold. When $\sigma=\tau, \gamma=\frac{1-\sigma \mu}{\mu}, \varepsilon=1, A=\frac{\mu+L}{2+(\mu+L) \varepsilon \gamma}$, we have

$$
\boldsymbol{H}(\boldsymbol{x}) \succeq \frac{\mu}{4} \mathbb{I} .
$$

Proof. By definition of $\boldsymbol{H}$, we can compute
$(1+\sigma \varepsilon A) \cdot \boldsymbol{H}(\boldsymbol{x})=\left(\begin{array}{cc}\boldsymbol{C}_{0}(\boldsymbol{x})(\sigma \boldsymbol{A}+\gamma \boldsymbol{A}) & \frac{1}{2} \boldsymbol{C}_{0}(\boldsymbol{x})(1-\varepsilon \gamma A)-\frac{1}{2} \eta \boldsymbol{A} \\ \frac{1}{2} \boldsymbol{C}_{0}(\boldsymbol{x})(1-\varepsilon \gamma A)-\frac{1}{2} \eta \boldsymbol{A} & \eta \varepsilon \boldsymbol{A}\end{array}\right)$,
where $\eta=\sigma / \tau=1, \boldsymbol{C}_{0}(\boldsymbol{x})=\nabla^{2} f(\boldsymbol{x}) \boldsymbol{B}(\boldsymbol{x}) \nabla^{2} f(\boldsymbol{x})$. We want to find some constant $\lambda>0$, such that

$$
\binom{\boldsymbol{z}}{\boldsymbol{w}}^{T} \boldsymbol{H}(\boldsymbol{x})\binom{\boldsymbol{z}}{\boldsymbol{w}} \geq \lambda\left(\|\boldsymbol{z}\|^{2}+\|\boldsymbol{w}\|^{2}\right) .
$$

Observe that

$$
\begin{align*}
& \binom{\boldsymbol{z}}{\boldsymbol{w}}^{T} \boldsymbol{H}(\boldsymbol{x})\binom{\boldsymbol{z}}{\boldsymbol{w}}-\lambda\left(\|\boldsymbol{z}\|^{2}+\|\boldsymbol{w}\|^{2}\right) \\
& =\boldsymbol{z}^{T}\left(\boldsymbol{C}_{0}(\gamma+\sigma) A /(1+\sigma \varepsilon A)-\lambda \mathbb{I}\right) \boldsymbol{z}+\boldsymbol{w}^{T}(\varepsilon A /(1+\sigma \varepsilon A)-\lambda \mathbb{I}) \boldsymbol{w} \\
& \quad+\boldsymbol{z}^{T}\left(-A+\boldsymbol{C}_{0}(\mathbb{I}-\varepsilon \gamma A)\right) \boldsymbol{w} /(1+\sigma \varepsilon A), \tag{3.17}
\end{align*}
$$

which is almost the same as Eq. (3.5). Thus, following a similar procedure in Theorem 3.3 with the provided parameters, we obtain that

$$
\lambda \geq \frac{\mu}{2} \frac{1+\frac{A \sigma}{2}}{1+\sigma A} \geq \frac{\mu}{4} .
$$

This implies

$$
\begin{equation*}
\boldsymbol{H}(\boldsymbol{x}) \succeq \frac{\mu}{4} \mathbb{I} . \tag{3.18}
\end{equation*}
$$

Corollary 3.9. Let $f: \mathbb{R}^{d} \rightarrow \mathbb{R}$ be a $\mathcal{C}^{4}$ strongly convex function. Suppose Assumption 3.1 and Assumption 3.2 hold. If $\sigma=\tau, \gamma=\frac{1-\sigma \mu}{\mu}, \varepsilon=1, A=$ $\frac{\mu+L}{2+(\mu+L) \varepsilon \gamma}$, we have

$$
\begin{equation*}
\|\boldsymbol{N}(\boldsymbol{x})\|_{2} \leq \frac{\max \{L, 1\} \cdot(A(\sigma+2 \gamma+2)+1)}{(1+\sigma A)} \tag{1}
\end{equation*}
$$

$$
\begin{equation*}
\boldsymbol{N}(\boldsymbol{x})^{T} \nabla^{2} \mathcal{I}(\boldsymbol{x}, \boldsymbol{p}) \boldsymbol{N}(\boldsymbol{x}) \preceq \frac{(3+\sigma A+2 A)^{2}}{(1+\sigma A)^{2}} \cdot \max \left\{L^{\prime}, 1\right\} \cdot \mathbb{I} \tag{2}
\end{equation*}
$$

Proof. We can decompose

$$
(1+\sigma A) \cdot \boldsymbol{N}(\boldsymbol{x})=\left(\begin{array}{cc}
\boldsymbol{B}(\boldsymbol{x}) \nabla^{2} f(\boldsymbol{x}) & 0 \\
0 & \mathbb{I}
\end{array}\right)\left(\begin{array}{cc}
(\sigma+\gamma) \boldsymbol{A} & (\mathbb{I}-\gamma \boldsymbol{A}) \\
-\boldsymbol{A} & \boldsymbol{A}
\end{array}\right) .
$$

Observe that

$$
\left(\begin{array}{cc}
(\sigma+\gamma) \boldsymbol{A} & (\mathbb{I}-\gamma \boldsymbol{A}) \\
-\boldsymbol{A} & \boldsymbol{A}
\end{array}\right)=\left(\begin{array}{cc}
(\sigma+\gamma) \mathbb{I} & -\gamma \mathbb{I} \\
-\mathbb{I} & \mathbb{I}
\end{array}\right) \cdot\left(\begin{array}{cc}
\boldsymbol{A} & 0 \\
0 & \boldsymbol{A}
\end{array}\right)+\left(\begin{array}{ll}
0 & \mathbb{I} \\
0 & 0
\end{array}\right) .
$$

Therefore,

$$
\begin{align*}
\left\|\left(\begin{array}{cc}
(\sigma+\gamma) \boldsymbol{A} & (\mathbb{I}-\gamma \boldsymbol{A}) \\
-\boldsymbol{A} & \boldsymbol{A}
\end{array}\right)\right\|_{2} & \leq A\left\|\left(\begin{array}{cc}
(\sigma+\gamma) \mathbb{I} & -\gamma \mathbb{I} \\
-\mathbb{I} & \mathbb{I}
\end{array}\right)\right\|_{2}+\left\|\left(\begin{array}{ll}
0 & \mathbb{I} \\
0 & 0
\end{array}\right)\right\|_{2} \\
& \leq A\left\|\left(\begin{array}{cc}
(\sigma+\gamma) \mathbb{I} & -\gamma \mathbb{I} \\
-\mathbb{I} & \mathbb{I}
\end{array}\right)\right\|_{2}+1 \\
& \leq A(\sigma+2 \gamma+2)+1 . \tag{3.19}
\end{align*}
$$

To get the last inequality, we consider $(\boldsymbol{z}, \boldsymbol{w})$ such that $\|\boldsymbol{z}\|^{2}+\|\boldsymbol{w}\|^{2}=1$. Thus

$$
\begin{aligned}
\left\|\left(\begin{array}{cc}
(\sigma+\gamma) \mathbb{I} & -\gamma \mathbb{I} \\
-\mathbb{I} & \mathbb{I}
\end{array}\right)\binom{\boldsymbol{z}}{\boldsymbol{w}}\right\|_{2} & =\left\|\binom{(\sigma+\gamma) \boldsymbol{z}-\gamma \boldsymbol{w}}{-\boldsymbol{z}+\boldsymbol{w}}\right\| \\
& \leq \sigma\left\|\binom{\boldsymbol{z}}{\boldsymbol{w}}\right\|+\gamma\|\boldsymbol{z}-\boldsymbol{w}\|+\|\boldsymbol{z}-\boldsymbol{w}\| \\
& \leq \sigma+2 \gamma+2 .
\end{aligned}
$$

We now have

$$
(1+\sigma A)\|\boldsymbol{N}(\boldsymbol{x})\|_{2} \leq\left\|\left(\begin{array}{cc}
\boldsymbol{B}(\boldsymbol{x}) \nabla^{2} f(\boldsymbol{x}) & 0 \\
0 & \mathbb{I}
\end{array}\right)\right\|_{2}\left\|\left(\begin{array}{cc}
(\sigma+\gamma) \boldsymbol{A} & (\mathbb{I}-\gamma \boldsymbol{A}) \\
-\boldsymbol{A} & \boldsymbol{A}
\end{array}\right)\right\|_{2}
$$

$$
\leq \max \{L, 1\} \cdot(A(\sigma+2 \gamma+2)+1) .
$$

This proves part (1) of our Corollary. It follows that

$$
\begin{align*}
& \boldsymbol{N}(\boldsymbol{x})^{T} \nabla^{2} \mathcal{I}(\boldsymbol{x}, \boldsymbol{p}) \boldsymbol{N}(\boldsymbol{x})= \\
& =\frac{1}{(1+\sigma A)^{2}}\left(\begin{array}{cc}
(\sigma+\gamma) \boldsymbol{A} & -\boldsymbol{A} \\
(\mathbb{I}-\gamma \boldsymbol{A}) & \boldsymbol{A}
\end{array}\right) \\
&  \tag{3.20}\\
& \cdot\left(\begin{array}{cc}
\boldsymbol{C}(\boldsymbol{x})^{T}\left(\nabla^{3} f(\boldsymbol{x}) \nabla f(\boldsymbol{x})+\left(\nabla^{2} f(\boldsymbol{x})\right)^{2}\right) \boldsymbol{C}(\boldsymbol{x}) & 0 \\
0
\end{array}\right) \\
& \\
& 3.20)
\end{align*}
$$

where we recall $\boldsymbol{C}(\boldsymbol{x})=\boldsymbol{B}(\boldsymbol{x}) \nabla^{2} f(\boldsymbol{x})$.
By assumption, there exists $L^{\prime}>0$, such that

$$
\boldsymbol{C}(\boldsymbol{x})^{T}\left(\nabla^{3} f(\boldsymbol{x}) \nabla f(\boldsymbol{x})+\left(\nabla^{2} f(\boldsymbol{x})\right)^{2}\right) \boldsymbol{C}(\boldsymbol{x}) \preceq L^{\prime} \mathbb{I},
$$

for all $\boldsymbol{x}$. Then, combining Eq. (3.20) and Eq. (3.19), we obtain that

$$
\begin{align*}
\left\|\boldsymbol{N}(\boldsymbol{x})^{T} \nabla^{2} \mathcal{I}(\boldsymbol{x}, \boldsymbol{p}) \boldsymbol{N}(\boldsymbol{x})\right\|_{2} & \leq \frac{(A(\sigma+2 \gamma+2)+1)^{2}}{(1+\sigma A)^{2}} \cdot \max \left\{L^{\prime}, 1\right\} \\
& \leq \frac{(3+\sigma A+2 A)^{2}}{(1+\sigma A)^{2}} \cdot \max \left\{L^{\prime}, 1\right\} \tag{3.21}
\end{align*}
$$

where we have used $\gamma A<1$ to derive the last inequality.
Theorem 3.10 (Restatement of Theorem 1.1). Suppose Assumption 3.1 and Assumption 3.2 hold. Let $\sigma=\tau, \gamma=\frac{1-\sigma \mu}{\mu}, \varepsilon=1, A=\frac{\mu+L}{2+(\mu+L) \varepsilon \gamma}$. And suppose (3.13) holds for some $\delta>0$ and all $\boldsymbol{x}$. If $\tau=\frac{1}{4} \frac{\mu}{\delta+36 \max \left\{L^{\prime}, 1\right\}}$, then

$$
\mathcal{I}\left(\boldsymbol{x}^{n}, \boldsymbol{p}^{n}\right) \leq \mathcal{I}\left(\boldsymbol{x}^{0}, \boldsymbol{p}^{0}\right)\left(1-\frac{\mu^{2} / 32}{\delta+36 \max \left\{L^{\prime}, 1\right\}}\right)^{n}
$$

Proof. By Assumption 3.1 and Assumption 3.2, we have $\mu \leq L \leq L^{\prime}$. Thus $\mu / L^{\prime} \leq 1$ and $\sigma=\tau<1 / 36$. Moreover,

$$
\gamma=\frac{1}{\mu}-\sigma \geq 1-\frac{1}{36}=\frac{35}{36} .
$$

And

$$
A=\frac{\mu+L}{2+(\mu+L) \varepsilon \gamma}<\frac{1}{\gamma} \leq \frac{36}{35} .
$$

Then it follows

$$
\frac{3+\sigma A+2 A}{1+\sigma A} \leq 3+\sigma A+2 A<6
$$

By Corollary 3.9, we have

$$
\left\|\boldsymbol{N}(\boldsymbol{x})^{T} \nabla^{2} \mathcal{I}(\boldsymbol{x}, \boldsymbol{p}) \boldsymbol{N}(\boldsymbol{x})\right\|_{2} \leq 36 \max \left\{L^{\prime}, 1\right\}
$$

Combining this with Theorem 3.6 and Corollary 3.8, we finish the proof.

Remark 3.11. The choice of parameters in Theorem 3.10 may not be optimal. The main purpose of Theorem 3.10 is to show the existence of geometric convergence in Algorithm 1.

## 4. Numerical experiment

We test our PDD algorithm using several convex and non-convex functions and compare the results with other commonly used optimizers, such as gradient descent, Nesterov's accelerated gradient (NAG), IGAHD (inertial gradient algorithm with Hessian damping) [2], and IGAHD-SC (inertial gradient algorithm with Hessian damping for strongly convex functions) [2]. In all of our numerical examples, we set tol $=0$ for better demonstration of the convergence comparison. In practice, one can choose the size of tol.

### 4.1. Summary of algorithms

For reference, we list down the iterations of gradient descent, NAG, IGAHD, and IGAHD-SC for better comparison.

Gradient descent:

$$
\boldsymbol{x}^{n+1}=\boldsymbol{x}^{n}-\tau_{\mathrm{gd}} \nabla f\left(\boldsymbol{x}^{n}\right),
$$

where $\tau_{\mathrm{gd}}>0$ is a stepsize.
NAG:

$$
\begin{aligned}
& \boldsymbol{y}^{n+1}=\boldsymbol{x}^{n}-\tau_{\mathrm{nag}} \nabla f\left(\boldsymbol{x}^{n}\right) \\
& \boldsymbol{x}^{n+1}=\boldsymbol{y}^{n+1}+\beta_{\mathrm{nag}}\left(\boldsymbol{y}^{n}-\boldsymbol{y}^{n-1}\right)
\end{aligned}
$$

where $\tau_{\text {nag }}>0$ is a stepsize, and $\beta_{\text {nag }}>0$ is a parameter.
IGAHD: Suppose $\nabla f$ is $L_{1}$-Lipschitz.

$$
\begin{aligned}
& \boldsymbol{y}^{n}= \boldsymbol{x}^{n}+\alpha_{n}\left(\boldsymbol{x}^{n}-\boldsymbol{x}^{n-1}\right)-\beta^{(1)} \sqrt{\tau_{\mathrm{att}}}\left(\nabla f\left(\boldsymbol{x}^{n}\right)-\nabla f\left(\boldsymbol{x}^{n-1}\right)\right) \\
& \quad-\frac{\beta^{(1)} \sqrt{\tau_{\mathrm{att}}} \nabla f\left(\boldsymbol{x}^{n-1}\right),}{n}, \\
& \boldsymbol{x}^{n+1}=\boldsymbol{y}^{n}-\tau_{\mathrm{att}} \nabla f\left(\boldsymbol{y}^{n}\right) .
\end{aligned}
$$

Here $\alpha_{n}=1-\frac{\alpha}{n}$ for some $\alpha \geq 3$. $\beta^{(1)}$ needs to satisfy

$$
0 \leq \beta^{(1)} \leq 2 \sqrt{\tau_{\mathrm{att}}} .
$$

And $\tau_{\text {att }}>0$ is a stepsize, which needs to satisfy

$$
\tau_{\mathrm{att}} \leq \frac{1}{L_{1}} .
$$

Remark 4.1. As mentioned earlier, in each iteration of IGAHD, $\nabla f(\cdot)$ is evaluated twice: at $\boldsymbol{x}^{n}$ and at $\boldsymbol{y}^{n}$. This differs from one gradient evaluation in gradient descent, NAG, and our method Eq. (1.2). Chen and Luo [8] proposed a slightly different algorithm from IGAHD that only requires one gradient evaluation in each iteration.

IGHD-SC: Suppose $f$ is $m_{1}$-strongly convex and $\nabla f$ is $L_{1}$-Lipschitz.

$$
\begin{aligned}
& \boldsymbol{x}^{n+1}=\boldsymbol{x}^{n}+\frac{1-\sqrt{m_{1} \tau_{\mathrm{att}}}}{1+\sqrt{m_{1} \tau_{\mathrm{att}}}}\left(\boldsymbol{x}^{n}-\boldsymbol{x}^{n-1}\right)-\frac{\beta^{(2)} \sqrt{\tau_{\mathrm{att}}}}{1+\sqrt{m_{1} \tau_{\mathrm{att}}}}\left(\nabla f\left(\boldsymbol{x}^{n}\right)\right. \\
&\left.-\nabla f\left(\boldsymbol{x}^{n-1}\right)\right)-\frac{\tau_{\mathrm{att}}}{1+\sqrt{m_{1} \tau_{\mathrm{att}}}} \nabla f\left(\boldsymbol{x}^{n}\right) .
\end{aligned}
$$

Here $\beta^{(2)}$ and $L_{1}$ need to satisfy

$$
\begin{equation*}
\beta^{(2)} \leq \frac{1}{\sqrt{m_{1}}}, \quad L_{1} \leq \min \left\{\frac{\sqrt{m_{1}}}{8 \beta^{(2)}}, \frac{\frac{\sqrt{m_{1}}}{2 \tau_{\mathrm{att}}}+\frac{m_{1}}{\sqrt{\tau_{\mathrm{tatt}}}}}{2 \beta^{(2)} m_{1}+\frac{1}{\sqrt{\tau_{\mathrm{attt}}}}+\frac{\sqrt{m_{1}}}{2}}\right\} . \tag{4.1}
\end{equation*}
$$

## 4.2. regularized log-sum-exp

Consider the regularized log-sum-exp function

$$
f(\boldsymbol{x})=\log \left(\sum_{i=1}^{n} \exp \left(\boldsymbol{q}_{i}^{T} \boldsymbol{x}\right)\right)+\frac{1}{2} \boldsymbol{x}^{T} \boldsymbol{Q} \boldsymbol{x}
$$

where $n=100, \boldsymbol{Q}=\boldsymbol{Q}^{T} \succ 0$ and $\boldsymbol{q}_{i}^{T}$ is the ith row of $\boldsymbol{Q} . \boldsymbol{Q}$ is chosen to be diagonally dominant, i.e. $Q_{i, i}>\sum_{j \neq i}\left|Q_{i, j}\right|$. In this case, we may choose the diagonal preconditioner $\boldsymbol{C}(\boldsymbol{x})=(\operatorname{diag}(\boldsymbol{Q}))^{-1}$. We compare the performance of gradient descent, preconditioned gradient descent, PDD with $\boldsymbol{C}(\boldsymbol{x})=\mathbb{I}$, PDD with diagonal preconditioner, NAG, and IGAHD-SC (inertial gradient algorithm with Hessian damping for strongly convex functions) by Attouch et al. [2] methods for minimizing $f$. The stepsize of gradient descent is determined by $\tau_{\mathrm{gd}}=\frac{2}{\lambda_{1} * 3+\lambda_{n}}$, where $\lambda_{1}$ and $\lambda_{n}$ are the maximum and minimum eigenvalues of $\boldsymbol{Q}$, respectively. For a pure quadratic objective function, $\boldsymbol{x}^{T} \boldsymbol{Q} \boldsymbol{x}$, the optimal stepsize of gradient descent is $\frac{2}{\lambda_{1}+\lambda_{n}}$. However, since our objective function also contains a log-sum-exp term, we slightly change the stepsize, with the fastest guaranteed convergence speed. Otherwise, gradient descent will not converge. Similarly, the optimal parameters of NAG for a purely quadratic function is $\tau=\frac{4}{3 \lambda_{1}+\lambda_{n}}$, and $\beta=\frac{\sqrt{3 \kappa+1}-2}{\sqrt{3 \kappa+1}+2}$, where $\kappa=\lambda_{1} / \lambda_{n}$. However, NAG with this choice of $\tau$ and $\kappa$ will not converge since $f$ is not purely quadratic. We thus choose $\tau_{\text {nag }}=\frac{4}{3 \lambda_{1}^{\prime}+\lambda_{n}}$ and $\beta_{\text {nag }}=\frac{\sqrt{3 \kappa^{\prime}+1}-2}{\sqrt{3 \kappa^{\prime}+1}+2}$, where $\kappa^{\prime}=\lambda_{1}^{\prime} / \lambda_{n}$ and $\lambda_{1}^{\prime}=a_{\text {nag }} \lambda_{1}$. In our experiment, $a_{\text {nag }}=5$ gives the fastest convergence for NAG. For PDD with $\boldsymbol{C}(\boldsymbol{x})=\mathbb{I}$, we choose $\tau_{\text {pdd }}=\sigma_{\mathrm{pdd}}=\frac{2}{a_{\mathrm{pdd}} \lambda_{1}+\lambda_{n}}, \varepsilon=1, A=10, \omega=1$. In our experiment $a_{\mathrm{pdd}}=0.5$ gives the fastest convergence. For PDD with diagonal preconditioner $\boldsymbol{C}(\boldsymbol{x})=(\operatorname{diag}(\boldsymbol{Q}))^{-1}$, we choose $\tau_{\text {pdd,prec }}=\sigma_{\text {pdd,prec }}=0.5, \varepsilon=1$, $A=1, \omega=1$. We use the same $\boldsymbol{C}(\boldsymbol{x})=(\operatorname{diag}(\boldsymbol{Q}))^{-1}$ as a preconditioner for gradient descent. The stepsize for preconditioned gradient descent is chosen to be the same as $\tau_{\text {pdd,prec }}=0.5$. For IGAHD-SC ('att'), we need $m_{1}$ as the smallest eigenvalue of $\nabla^{2} f(\boldsymbol{x})$. In this example, we may estimate $m_{1}$ as the smallest eigenvalue of $\boldsymbol{Q}$. And $\tau_{\text {att }}=0.0016$ via grid search. $\beta^{(2)}$ in IGAHD-SC is found by solving (see Theorem 11 Eq. (26) of Attouch et al. [2])

$$
\frac{\sqrt{m_{1}}}{8 \beta^{(2)}}=\frac{\frac{\sqrt{m_{1}}}{2 \tau_{\mathrm{att}}}+\frac{\sqrt{m_{1}}}{\sqrt{\tau_{\mathrm{att}}}}}{2 \beta^{(2)} m_{1}+\frac{1}{\sqrt{\tau_{\mathrm{att}}}}+\frac{\sqrt{m_{1}}}{2}}
$$

which gives

$$
\begin{equation*}
\beta^{(2)}=\frac{\sqrt{\tau_{\mathrm{att}}}+\tau_{\mathrm{att}} \sqrt{m_{1}} / 2}{4+8 \sqrt{m_{1}} \sqrt{\tau_{\mathrm{att}}}-2 m_{1} \tau_{\mathrm{att}}} . \tag{4.2}
\end{equation*}
$$

The initial condition is $\boldsymbol{x}^{0}=\operatorname{np}$.ones $(n) * 0.1$. The result is presented in Fig. 1a.


Figure 1: Comparison of gradient descent, NAG, PDD, and IGAHD-SC (we use 'att' as a shorthand for this method) on minimizing (a) the regularized log-sum-exp function and (b) the quadratic minus cosine function. The $y$ axis represents the 2-norm of the gradient of the objective function on a logarithmic scale. The x -axis represents the number of iterations on a logarithmic scale.

### 4.3. Quadratic minus cosine function

Consider the function

$$
f(\boldsymbol{x})=\|\boldsymbol{x}\|^{2}-\cos \left(\boldsymbol{c}^{T} \boldsymbol{x}\right),
$$

where $\boldsymbol{c}$ is a vector in $\mathbb{R}^{100}$ with $\|\boldsymbol{c}\|^{2}=1.9$. Then a direct calculation shows that $0.1 \mathbb{I} \preceq \nabla^{2} f(\boldsymbol{x}) \preceq 3.9 \mathbb{I}$ for any $\boldsymbol{x}$. This allows us to choose the optimal stepsize for gradient descent and NAG. When minimizing $f$ using gradient descent, we can choose $\tau_{\mathrm{gd}}=\frac{2}{0.1+3.9}=0.5$. Meanwhile, for NAG, we may choose $\tau_{\text {nag }}=\frac{4}{3 * 3.9+0.1}$, and $\beta=\frac{\sqrt{3 \kappa+1}-2}{\sqrt{3 \kappa+1}+2}$, where $\kappa=3.9 / 0.1$. For PDD with $\boldsymbol{C}(\boldsymbol{x})=\mathbb{I}$, we choose $\tau_{\mathrm{pdd}}=\sigma_{\mathrm{pdd}}=0.5, \varepsilon=1, A=1, \omega=1$. For IGAHD-SC ('att'), we choose $m_{1}=0.1, \tau_{\text {att }}=0.55$ via grid search and $\beta^{(2)}$ is given by Eq. (4.2). The initial condition is $\boldsymbol{x}^{0}=\operatorname{np}$.ones $(n) * 5$. The result is presented in Fig. 1b.

### 4.4. Rosenbrock function

4.4.1. 2-dimension. The 2-dimensional Rosenbrock function is defined as

$$
f(x, y)=(a-x)^{2}+b\left(y-x^{2}\right)^{2},
$$



Figure 2: Minimizing the Rosenbrock function with gradient descent, NAG, PDD with $\boldsymbol{C}(\boldsymbol{x})=\mathbb{I}$ and IGAHD ('att'). The left panel shows the convergence speed of each method. The right panel shows the optimization trajectories of each method.
where a common choice of parameters is $a=1, b=100$. This is a non-convex function with a global minimum of $\left(x^{*}, y^{*}\right)=\left(a, a^{2}\right)$. The global minimum is inside a long, narrow, parabolic-shaped flat valley. To find the valley is trivial. To converge to the global minimum, however, is difficult. We compare the performance of gradient descent, NAG, PDD with $\boldsymbol{C}(\boldsymbol{x})=\mathbb{I}$ and IGAHD (inertia gradient algorithm with Hessian damping) by Attouch et al. [2] when minimizing the Rosenbrock function starting from $(-3,-4)$. The stepsize of gradient descent is $\tau_{\mathrm{gd}}=0.0002$. The stepsize of NAG is $\tau_{\mathrm{nag}}=0.0002$, $\beta_{\text {nag }}=0.9$. The parameters of PDD are $\tau_{\text {pdd }}=\sigma_{\text {pdd }}=0.005, \varepsilon=1, \omega=1$, $A=5$. The stepsize of the PDD method is larger than $\tau_{\mathrm{gd}}$ and $\tau_{\text {nag }}$ because gradient descent and NAG do not allow larger stepsizes for convergence. For IGAHD ('att'), we choose $\tau_{\mathrm{att}}=0.00045, \alpha=3, \beta^{(1)}=\sqrt{\tau_{\mathrm{att}}} / 14$. The convergence result and the optimization trajectories are shown in Fig. 2.
4.4.2. N-dimension. The $N$-dimensional coupled Rosenbrock function is defined as

$$
f(\boldsymbol{x})=\sum_{i=1}^{N-1}\left(\left(a-x_{i}\right)^{2}+b\left(x_{i+1}-x_{i}^{2}\right)^{2}\right),
$$

where we choose $a=1$ and $b=100$ as in the 2 -dimensioal case and we set $N=100$. The global minimum is at $\boldsymbol{x}^{*}=(1,1, \ldots, 1)$. Using the same stepsizes as in the 2-dimensional case, we compare the performance of the three algorithms starting from $\boldsymbol{x}_{0}=(0, \ldots, 0)$. The stepsize of gradient


Figure 3: Comparison of gradient descent, NAG, PDD with $\boldsymbol{C}(\boldsymbol{x})=\mathbb{I}$ and IGAHD ('att') on minimizing the 100-dimensional coupled Rosenbrock function. The $y$-axis represents the distance between the current iterate and the global minimum on a logarithmic scale. The x -axis represents the number of iterations on a logarithmic scale.
descent is $\tau_{\mathrm{gd}}=0.001$. The stepsize of NAG is $\tau_{\text {nag }}=0.0008, \beta=0.95$. The parameters of PDD are $\tau_{\mathrm{pdd}}=\sigma_{\mathrm{pdd}}=0.01, \varepsilon=0.5, \omega=1, A=5$. The stepsize of the PDD method is larger than $\tau_{\mathrm{gd}}$ and $\tau_{\text {nag }}$ because gradient descent and NAG do not allow larger stepsizes for convergence. For IGAHD ('att'), we choose $\tau_{\mathrm{att}}=0.0002, \alpha=3, \beta^{(1)}=2 * \sqrt{\tau_{\mathrm{att}}}$. The result is summarized in Fig. 3.

### 4.5. Ackley function

We consider minimizing the two-dimensional Ackley function given by

$$
\begin{aligned}
f(x, y)=- & 20 \exp \left(-0.2 \sqrt{0.5\left(x^{2}+y^{2}\right)}\right) \\
& -\exp [0.5(\cos (2 \pi x)+\cos (2 \pi y))]+\mathrm{e}+20,
\end{aligned}
$$

which has many local minima. The unique global minimum is located at $\left(x^{*}, y^{*}\right)=(0,0)$. We compare the performance of gradient descent, NAG, PDD, and IGAHD ('att') for minimizing the two-dimensional Ackley function starting from $\left(x_{0}, y_{0}\right)=(2.5,4)$. The stepsize of gradient descent is $\tau_{\mathrm{gd}}=0.002$. The stepsize of NAG is $\tau_{\text {nag }}=0.002, \beta_{\text {nag }}=0.9$. The parameters of PDD are $\tau_{\text {pdd }}=\sigma_{\text {pdd }}=0.002, \varepsilon=1, \omega=1, A=1$. For IGAHD ('att'), we choose $\tau_{\text {att }}=0.01, \alpha=3, \beta^{(1)}=2 * \sqrt{\tau_{\text {att }}}$. The results are summarized in Fig. 4.
Remark 4.2. We remark that our algorithm has no stochasticity. It will not always converge to the global minimum for non-convex functions in


Figure 4: Minimizing the Ackley function with gradient descent, NAG, PDD and IGAHD ('att'). The left panel shows the convergence speed of each method. The right panel shows the optimization trajectories of each method. The reason why only two colors appear on the left panel is because gradient descent, NAG and IGAHD do not converge and their error curves overlap with each other.
general. For example, it will not converge for the Griewank, Drop-Wave, and Rastrigin functions.
Remark 4.3. We believe that one of the reasons why Algorithm 1 could converge for some of the non-convex functions is that we choose $\boldsymbol{p}^{0}$ to point in the direction of the global minimizer. Suppose we start at a local minimum/maximum, i.e., $\nabla f\left(\boldsymbol{x}^{0}\right)=0$. Then the gradient descent algorithm will always be stuck at $\boldsymbol{x}^{0}$. However, in Algorithm 1, even if we have $\nabla f\left(\boldsymbol{x}^{0}\right)=0$, we can still choose $\boldsymbol{p}^{0}$ to be pointing in the direction of the global minimizer if we have some information about the global minimizer. Then Eq. (3.9b) for $n=1$ reads

$$
\boldsymbol{x}^{1}=\boldsymbol{x}^{0}-\tau \boldsymbol{B}\left(\boldsymbol{x}^{0}\right) \nabla^{2} f\left(\boldsymbol{x}^{0}\right) \frac{1-\varepsilon \gamma A}{1+\sigma \varepsilon A} \boldsymbol{p}^{0},
$$

which is a non-trivial update compared with gradient descent when $\nabla f\left(\boldsymbol{x}^{0}\right)=0$. In all of our numerical examples, we have chosen $\boldsymbol{p}^{0}=\boldsymbol{x}^{0}$. This choice of $\boldsymbol{p}^{0}$ works well for the non-convex functions we choose when the global minimizer is near the origin. We emphasize that this remark only serves as an intuition on the applicability of Algorithm 1 to some non-convex examples. We leave rigorous convergence analysis for non-convex functions to future works.

| Algorithm | SGD | NAG | PDD | Adam | Att |
| :--- | :--- | :--- | :--- | :--- | :--- |
| train loss | $2.223 \pm 0.034$ | $0.964 \pm 0.244$ | $\mathbf{0 . 4 3 3} \pm \mathbf{0 . 2 7 0}$ | $0.589 \pm 0.282$ | $0.591 \pm 0.288$ |
| test acc | $29.3 \pm 8.3 \%$ | $71.2 \pm 9.4 \%$ | $\mathbf{8 5 . 4} \pm \mathbf{1 0 . 4} \%$ | $79.1 \pm 11.3 \%$ | $80.8 \pm 11.4 \%$ |

Table 1: Average training loss and test accuracy of different algorithms for MNIST handwritten digit recognition over 60 random seeds.


Figure 5: Training a two-layer neural network with the MNIST data set using gradient descent, NAG, PDD, Adam, and IGAHD ('att'). The left panel shows the convergence speed of training loss. The right panel shows the test accuracy of each method. The x-axis represents the number of iterations in terms of mini-batches.

### 4.6. Neural Networks training

4.6.1. MNIST with Two-layer neural network. We consider the classification problem using the MNIST handwritten digit data set with a twolayer neural network. The neural network has an input layer of size $784=$ $28 \times 28$, a hidden layer of size 32 followed by another hidden layer of size 32 , and an output layer of size 10 . We use ReLU activation function across the layers, and the loss is evaluated using the cross-entropy loss. We use a batch size of 200 for all the algorithms. The stepsize of gradient descent is $\tau_{\mathrm{gd}}=0.001$. The stepsize of NAG is $\tau_{\text {nag }}=0.001$, momentum $=0.9$. The parameters of PDD are $\tau_{\mathrm{pdd}}=0.001, \sigma_{\mathrm{pdd}}=5, \varepsilon=0.005, \omega=1, A=1$. For IGAHD ('att'), we choose $\tau_{\mathrm{att}}=0.001, \alpha=3, \beta^{(1)}=0.01$. For Adam, we choose $\tau_{\text {adam }}=0.001, \beta_{1}=0.9, \beta_{2}=0.999$.
4.6.2. CIFAR10 with CNN. We train a convolutional neural network using the CIFAR10 datasets with SGD, Nesterov, PDD, Adam, and IGAHD


Figure 6: Training loss (left panel) and test accuracy (right panel) of a convolutional neural network on the CIFAR10 data set. The x -axis represents the number of iterations in terms of mini-batches.
('Att'). The architecture of the network is described as follows. The network consists of two convolutional layers. The first convolutional layer has 32 output channels, and the filter size is $3 \times 3$. The second convolutional layer has 64 output channels, and the filter size is $4 \times 4$. Each convolutional layer is followed by a ReLU activation and then a $2 \times 2$ max-pooling layer. Lastly, we have 3 fully connected layers of size $(64 \cdot 4 \cdot 4,120),(120,84)$, and $(84,10)$. The loss is evaluated using the cross-entropy loss. The stepsize of gradient descent is $\tau_{\mathrm{gd}}=0.01$. The stepsize of NAG is $\tau_{\text {nag }}=0.005$, momentum $=0.9$. The parameters of PDD are $\tau_{\mathrm{pdd}}=0.005, \sigma_{\mathrm{pdd}}=5, \varepsilon=0.005, \omega=1, A=1$. For IGAHD ('att'), we choose $\tau_{\mathrm{att}}=0.005, \alpha=3, \beta^{(1)}=0.01$. For Adam, we choose $\tau_{\text {adam }}=0.005, \beta_{1}=0.9, \beta_{2}=0.999$.

| Algorithm | SGD | NAG | PDD | Adam | Att |
| :--- | :--- | :--- | :--- | :--- | :--- |
| train loss | $2.038 \pm 0.070$ | $1.347 \pm 0.100$ | $\mathbf{0 . 6 9 7} \pm \mathbf{0 . 0 7 7}$ | $0.927 \pm 0.128$ | $0.879 \pm 0.092$ |
| test acc | $27.5 \pm 2.2 \%$ | $51.2 \pm 0.7 \%$ | $\mathbf{7 0 . 3} \pm \mathbf{0 . 5} \%$ | $64.4 \pm 2.7 \%$ | $66.8 \pm 0.6 \%$ |

Table 2: Average training loss and test accuracy of different algorithms for CIFAR10 data set over 60 random seeds.

## 5. Discussion

This paper presents primal-dual hybrid gradient algorithms for solving unconstrained optimization problems. We reformulate the optimality condition of the optimization problem as a saddle-point problem and then compute
the proposed saddle-point problem by a preconditioned PDHG method. We present the geometric convergence analysis for the strongly convex objective functions. In numerical experiments, we demonstrate that the proposed method works efficiently in non-convex optimization problems, at least in some examples, such as Rosenbrock and Ackley functions. In particular, it could efficiently train two-layer and convolution neural networks in supervised learning problems.

So far, our convergence study is limited to strongly convex objective functions, not convex ones. Meanwhile, the choice of preconditioners and stepsizes are independent of time. We also have not discussed the optimal choices of parameters or general proximal operators in the updates of algorithms. These generalized choices of functions, parameters, and their convergence properties have been intensively studied in Nesterov accelerated gradient methods and Hessian-driven damping methods [1, 2, 3]. In future work, we shall explore the convergence property of PDHG methods for convex functions with time-dependent parameters. We also investigate the convergence of similar algorithms in scientific computing problems of implicit time updates of partial differential equations [12, 13, 14].

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## References

[1] H. Attouch, Z. Chbani, and H. Riahi. Fast proximal methods via time scaling of damped inertial dynamics. SIAM Journal on Optimization, 29(3):2227-2256, 2019.
[2] H. Attouch, Z. Chbani, J. Fadili, and H. Riahi. First-order optimization algorithms via inertial systems with Hessian driven damping. Mathematical Programming, pages 1-43, 2020.
[3] H. Attouch, Z. Chbani, J. Fadili, and H. Riahi. Convergence of iterates for first-order optimization algorithms with inertia and Hessian driven damping. Optimization, pages 1-40, 2021.
[4] S. Boyd and L. Vandenberghe. Convex optimization. Cambridge university press, 2004.
[5] S. Boyd, N. Parikh, E. Chu, B. Peleato, J. Eckstein, et al. Distributed optimization and statistical learning via the alternating direction method of multipliers. Foundations and Trends $\circledR^{\Omega}$ in Machine learning, 3(1):1-122, 2011.
[6] A. Chambolle and T. Pock. A first-order primal-dual algorithm for convex problems with applications to imaging. Journal of mathematical imaging and vision, 40:120-145, 2011.
[7] L. Chen and H. Luo. First order optimization methods based on Hessian-driven Nesterov accelerated gradient flow. arXiv preprint arXiv:1912.09276, 2019.
[8] L. Chen and H. Luo. A unified convergence analysis of first order convex optimization methods via strong Lyapunov functions. arXiv preprint arXiv:2108.00132, 2021.
[9] D. Gabay and B. Mercier. A dual algorithm for the solution of nonlinear variational problems via finite element approximation. Computers $\mathcal{E}$ Mathematics with Applications, 2(1):17-40, 1976. ISSN 08981221. doi: https://doi.org/10.1016/0898-1221(76)90003-1. URL https: //www.sciencedirect.com/science/article/pii/0898122176900031.
[10] M. Jacobs, F. Léger, W. Li, and S. Osher. Solving large-scale optimization problems with a convergence rate independent of grid size. SIAM Journal on Numerical Analysis, 57(3):1100-1123, 2019.
[11] D. P. Kingma and J. Ba. Adam: A method for stochastic optimization. arXiv preprint arXiv:1412.6980, 2014.
[12] W. Li, S. Liu, and S. Osher. Controlling conservation laws II: Compressible Navier-Stokes equations. Journal of Computational Physics, 463:111264, 2022.
[13] W. Li, S. Liu, and S. Osher. Controlling conservation laws I: Entropyentropy flux. Journal of Computational Physics, 480:112019, 2023.
[14] S. Liu, S. Osher, W. Li, and C.-W. Shu. A primal-dual approach for solving conservation laws with implicit in time approximations. Journal of Computational Physics, 472:111654, 2023.
[15] Y. Liu, Y. Xu, and W. Yin. Acceleration of primal-dual methods by preconditioning and simple subproblem procedures. Journal of Scientific Computing, 86(2):21, 2021.
[16] Y. E. Nesterov. A method of solving a convex programming problem with convergence rate $\mathcal{O}\left(\frac{1}{k^{2}}\right)$. In Doklady Akademii Nauk, volume 269, pages 543-547. Russian Academy of Sciences, 1983.
[17] S. Osher, B. Wang, P. Yin, X. Luo, F. Barekat, M. Pham, and A. Lin. Laplacian smoothing gradient descent. Research in the Mathematical Sciences, 9(3):55, 2022.
[18] Y. Ouyang, Y. Chen, G. Lan, and E. Pasiliao Jr. An accelerated linearized alternating direction method of multipliers. SIAM Journal on Imaging Sciences, 8(1):644-681, 2015.
[19] J.-H. Park, A. J. Salgado, and S. M. Wise. Preconditioned accelerated gradient descent methods for locally lipschitz smooth objectives with applications to the solution of nonlinear PDEs. Journal of Scientific Computing, 89(1):17, 2021.
[20] R. Pascanu and Y. Bengio. Revisiting natural gradient for deep networks. arXiv preprint arXiv:1301.3584, 2013.
[21] T. Pock and A. Chambolle. Diagonal preconditioning for first order primal-dual algorithms in convex optimization. In 2011 International Conference on Computer Vision, pages 1762-1769, 2011. doi: 10.1109/ ICCV.2011.6126441.
[22] B. T. Polyak. Some methods of speeding up the convergence of iteration methods. USSR computational mathematics and mathematical physics, 4(5):1-17, 1964.
[23] T. Rees. Preconditioning iterative methods for PDE constrained optimization. PhD thesis, University of Oxford Oxford, UK, 2010.
[24] J. W. Siegel. Accelerated first-order methods: Differential equations and Lyapunov functions. arXiv preprint arXiv:1903.05671, 2019.
[25] W. Su, S. Boyd, and E. J. Candes. A differential equation for modeling Nesterov's accelerated gradient method: theory and insights. arXiv preprint arXiv:1503.01243, 2015.
[26] L. N. Trefethen and D. Bau. Numerical linear algebra, volume 181. SIAM, 2022.
[27] T. Valkonen. A primal-dual hybrid gradient method for nonlinear operators with applications to MRI. Inverse Problems, 30(5):055012, 2014.
[28] A. Wibisono, A. C. Wilson, and M. I. Jordan. A variational perspective
on accelerated methods in optimization. Proceedings of the National Academy of Sciences, 113(47):E7351-E7358, 2016.
[29] M. Zhu and T. Chan. An efficient primal-dual hybrid gradient algorithm for total variation image restoration. Ucla Cam Report, 34:8-34, 2008.

## Appendix A Matrix lemma

Lemma A.1. Let $\boldsymbol{A}, \boldsymbol{B}, \boldsymbol{C} \in \mathbb{R}^{n}$ be real symmetric matrices that are simultaneously diagonalizable. Then for any $\boldsymbol{x}, \boldsymbol{y} \in \mathbb{R}^{n}$, if

$$
\begin{aligned}
& \lambda_{A, i}+\frac{\left|\lambda_{C, i}\right|}{2} \leq 0 \\
& \lambda_{B, i}+\frac{\left|\lambda_{C, i}\right|}{2} \leq 0
\end{aligned}
$$

for all $i$, where $\lambda_{A, i}, \lambda_{B, i}, \lambda_{C, i}$ are the ith eigenvalues of $\boldsymbol{A}, \boldsymbol{B}, \boldsymbol{C}$ respectively in the same basis. Then

$$
\boldsymbol{x}^{T} \boldsymbol{A} \boldsymbol{x}+\boldsymbol{y}^{T} \boldsymbol{B} \boldsymbol{y}+\boldsymbol{x}^{T} \boldsymbol{C} \boldsymbol{y} \leq 0,
$$

for all $\boldsymbol{x}, \boldsymbol{y} \in \mathbb{R}^{n}$.
Proof. Let $\boldsymbol{x}, \boldsymbol{y} \in \mathbb{R}^{n}$. By our assumption, there exists $\boldsymbol{Q}$ unitary such that $\boldsymbol{A}, \boldsymbol{B}, \boldsymbol{C}$ are simultaneously diagonalizable by $\boldsymbol{Q}$. Set $\tilde{\boldsymbol{x}}=\boldsymbol{Q} \boldsymbol{x}$ and $\tilde{\boldsymbol{y}}=\boldsymbol{Q} \boldsymbol{y}$. Then we can compute

$$
\begin{aligned}
\boldsymbol{x}^{T} \boldsymbol{A} \boldsymbol{x}+\boldsymbol{y}^{T} \boldsymbol{B} \boldsymbol{y}+\boldsymbol{x}^{T} \boldsymbol{C} \boldsymbol{y} & =\sum_{i=1}^{n} \tilde{x}_{i}^{2} \lambda_{A, i}+\tilde{y}_{i}^{2} \lambda_{B, i}+\tilde{x}_{i} \lambda_{C, i} \tilde{y}_{i} \\
& \leq \sum_{i=1}^{n} \tilde{x}_{i}^{2}\left(\lambda_{A, i}+\frac{\left|\lambda_{C, i}\right|}{2}\right)+\tilde{y}_{i}^{2}\left(\lambda_{B, i}+\frac{\left|\lambda_{C, i}\right|}{2}\right) \\
& \leq 0,
\end{aligned}
$$

where the first inequality follows from $\alpha x y \leq\left(x^{2}+y^{2}\right)|\alpha| / 2$ for any $\alpha, x, y \in$ $\mathbb{R}$.

## Appendix B Proof of Theorem 2.6

## B. 1 Part (a)

We have the following system of ODE:

$$
\binom{\dot{\boldsymbol{x}}}{\dot{\boldsymbol{p}}}=\left(\begin{array}{cc}
-\gamma \boldsymbol{B} \boldsymbol{Q} \boldsymbol{A} \boldsymbol{Q} & -\boldsymbol{B} \boldsymbol{Q}(\mathbb{I}-\gamma \varepsilon \boldsymbol{A})  \tag{B.1}\\
\boldsymbol{A} \boldsymbol{Q} & -\varepsilon \boldsymbol{A}
\end{array}\right)\binom{\boldsymbol{x}}{\boldsymbol{p}} .
$$

Let us compute the eigenvalues of the above system. Let $\alpha$ be an eigenvalue, then $\alpha$ satisfies

$$
\begin{aligned}
& 0=\operatorname{det}\left(\begin{array}{cc}
-\gamma \boldsymbol{B} \boldsymbol{Q} \boldsymbol{A} \boldsymbol{Q}-\alpha \mathbb{I} & -\boldsymbol{B} \boldsymbol{Q}(\mathbb{I}-\gamma \varepsilon \boldsymbol{A}) \\
\boldsymbol{A} \boldsymbol{Q} & -\varepsilon \boldsymbol{A}-\alpha \mathbb{I}
\end{array}\right), \\
& 0=\operatorname{det}((-\gamma \boldsymbol{B} \boldsymbol{Q} \boldsymbol{A} \boldsymbol{Q}-\alpha \mathbb{I})(-\varepsilon \boldsymbol{A}-\alpha \mathbb{I})+\boldsymbol{B} \boldsymbol{Q}(\mathbb{I}-\gamma \varepsilon \boldsymbol{A}) \boldsymbol{A} \boldsymbol{Q}), \\
& 0=\operatorname{det}\left(\alpha^{2} \mathbb{I}+\alpha(\varepsilon \boldsymbol{A}+\gamma \boldsymbol{B} \boldsymbol{Q} \boldsymbol{A} \boldsymbol{Q})+\boldsymbol{B} \boldsymbol{Q} \boldsymbol{A} \boldsymbol{Q}\right) .
\end{aligned}
$$

The last equality is because $\boldsymbol{A}$ commutes with $\boldsymbol{Q}$. We assume that $\boldsymbol{A}$ and $B Q A Q$ are simultaneously diagonalizable. Thus,

$$
\begin{aligned}
0 & =\alpha^{2}+\alpha\left(\varepsilon a_{i}+\gamma \mu_{i}\right)+\mu_{i} \\
\alpha & =\frac{-\varepsilon a_{i}-\gamma \mu_{i} \pm \sqrt{\left(\varepsilon+\gamma \mu_{i}\right)^{2}-4 \mu_{i}}}{2}
\end{aligned}
$$

If $\gamma>0$ and $\varepsilon \geq 0$, then the real part of the eigenvalues are negative, and the system will converge. The convergence rate depends on the largest real part of the eigenvalues, which is

$$
\max _{i} \frac{1}{2}\left[-\gamma \mu_{i}-\varepsilon a_{i}+\Re\left(\sqrt{\left(\gamma \mu_{i}+\varepsilon\right)^{2}-4 \mu_{i}}\right)\right] .
$$

## B. 2 Part (c)

When $\gamma=\varepsilon=0$, we see that $\alpha$ is purely imaginary. Thus solutions to Eq. (B.1) will be oscillatory and will not converge.

## B. 3 Part (b)

Let us define

$$
g(\gamma)=\max _{i}\left\{\frac{\mu_{i}\left(-\gamma+\Re\left(\sqrt{\gamma^{2}-4 / \mu_{i}}\right)\right)}{2}\right\} .
$$

Essentially, we would like to find $\gamma^{*}=\arg \min _{\gamma} g(\gamma)$. We then define

$$
\begin{aligned}
\gamma(\mu) & :=\underset{\gamma}{\arg \min } \frac{\mu\left(-\gamma+\Re\left(\sqrt{\gamma^{2}-4 / \mu}\right)\right)}{2} \\
& =\frac{2}{\sqrt{\mu}} .
\end{aligned}
$$

Observe that if $\gamma \geq 2 / \sqrt{\mu_{n}}$, then $\gamma^{2}-4 / \mu_{i} \geq 0$ for all $i$. Thus

$$
g(\gamma)=\max _{i}\left\{\frac{\mu_{i}\left(-\gamma+\sqrt{\gamma^{2}-4 / \mu_{i}}\right)}{2}\right\} .
$$

For $\mu \in\left[\mu_{n}, \mu_{1}\right]$ and $\gamma \geq 2 / \sqrt{\mu_{n}}$, one can check that the function $\mu(-$ $\left.\gamma+\sqrt{\gamma^{2}-4 / \mu}\right)$ is increasing in $\mu$ by computing the partial derivative with respect to $\mu$. Then we get

$$
\begin{aligned}
g(\gamma) & =\frac{\mu_{1}\left(-\gamma+\sqrt{\gamma^{2}-4 / \mu_{1}}\right)}{2} \\
& \geq g\left(2 / \sqrt{\mu_{n}}\right)=\sqrt{\mu_{1}}(\sqrt{\kappa-1}-\sqrt{\kappa}) \approx-\sqrt{\mu_{n}} / 2,
\end{aligned}
$$

where $\kappa=\mu_{1} / \mu_{n}>1$. The last approximation is valid for $\mu_{1} / \mu_{n} \gg 1$. This shows that $\gamma^{*} \leq 2 / \sqrt{\mu_{n}}$. Similarly, if $\gamma \leq 2 / \sqrt{\mu_{1}}$, then $\gamma^{2}-4 / \mu_{i} \leq 0$ for all $i$. Thus

$$
\begin{aligned}
g(\gamma) & =\max _{i}\left\{\frac{-\mu_{i} \gamma}{2}\right\} \\
& =\frac{-\mu_{n} \gamma}{2} \\
& \geq-\frac{\mu_{n}}{\sqrt{\mu_{1}}}=g\left(2 / \sqrt{\mu_{1}}\right) .
\end{aligned}
$$

This shows that $\gamma^{*} \geq 2 / \sqrt{\mu_{1}}$. Combining with our previous observation, we get $\gamma^{*} \in\left[2 / \sqrt{\mu_{1}}, 2 / \sqrt{\mu_{n}}\right]$. Now let us fix some $\gamma^{\prime} \in\left[2 / \sqrt{\mu_{1}}, 2 / \sqrt{\mu_{n}}\right]$. Let $j=\inf \left\{i: 1 \leq i \leq n, \gamma^{\prime 2}-4 / \mu_{i} \leq 0\right\}$. By our assumption on $\gamma^{\prime}$, we know that $1<j<n$. Now for $1 \leq i \leq j-1$, we have

$$
\begin{aligned}
\frac{\mu_{i}\left(-\gamma^{\prime}+\Re\left(\sqrt{\gamma^{\prime 2}-4 / \mu_{i}}\right)\right)}{2} & =\frac{\mu_{i}\left(-\gamma^{\prime}+\sqrt{\gamma^{\prime 2}-4 / \mu_{i}}\right)}{2} \\
& \leq \frac{\mu_{1}\left(-\gamma^{\prime}+\sqrt{\gamma^{\prime 2}-4 / \mu_{1}}\right)}{2} .
\end{aligned}
$$

And for $j \leq k \leq n$, we have

$$
\frac{\mu_{k}\left(-\gamma^{\prime}+\Re\left(\sqrt{\gamma^{\prime 2}-4 / \mu_{k}}\right)\right)}{2}=\frac{-\mu_{k} \gamma^{\prime}}{2} \leq \frac{-\mu_{n} \gamma^{\prime}}{2} .
$$

It is thus clear that for $\gamma^{\prime} \in\left[2 / \sqrt{\mu_{1}}, 2 / \sqrt{\mu_{n}}\right]$,

$$
g\left(\gamma^{\prime}\right)=\max \left\{\frac{\mu_{1}\left(-\gamma^{\prime}+\sqrt{\gamma^{\prime 2}-4 / \mu_{1}}\right)}{2}, \frac{-\mu_{n} \gamma^{\prime}}{2}\right\} .
$$

It is straightforward to calculate that for $\gamma \in\left[\frac{2}{\sqrt{\mu_{1}}}, \frac{2 \sqrt{\mu_{1}}}{\sqrt{\mu_{n}\left(2 \mu_{1}-\mu_{n}\right)}}\right]$, we have

$$
\frac{-\mu_{n} \gamma}{2} \geq \frac{\mu_{1}\left(-\gamma+\sqrt{\gamma^{2}-4 / \mu_{1}}\right)}{2} .
$$

So

$$
g(\gamma)=\frac{-\mu_{n} \gamma}{2} \geq g\left(\frac{2 \sqrt{\mu_{1}}}{\sqrt{\mu_{n}\left(2 \mu_{1}-\mu_{n}\right)}}\right)=\frac{-\sqrt{\mu_{n}}}{\sqrt{2-\frac{1}{\kappa}}} .
$$

And for $\gamma \in\left[\frac{2 \sqrt{\mu_{1}}}{\sqrt{\mu_{n}\left(2 \mu_{1}-\mu_{n}\right)}}, 2 / \sqrt{\mu_{n}}\right]$ we have

$$
\frac{-\mu_{n} \gamma}{2} \leq \frac{\mu_{1}\left(-\gamma+\sqrt{\gamma^{2}-4 / \mu_{1}}\right)}{2} .
$$

This implies

$$
g(\gamma)=\frac{\mu_{1}\left(-\gamma+\sqrt{\gamma^{2}-4 / \mu_{1}}\right)}{2} \geq g\left(\frac{2 \sqrt{\mu_{1}}}{\sqrt{\mu_{n}\left(2 \mu_{1}-\mu_{n}\right)}}\right)=\frac{-\sqrt{\mu_{n}}}{\sqrt{2-\frac{1}{\kappa}}}
$$

This shows $\gamma^{*}=\frac{2 \sqrt{\mu_{1}}}{\sqrt{\mu_{n}\left(2 \mu_{1}-\mu_{n}\right)}}$.

## B. 4 Part (d)

Define $\Delta_{\gamma}(\mu, \varepsilon)=(\gamma \mu+\varepsilon)^{2}-4 \mu$. Also define $g_{\gamma}(\mu)=2 \sqrt{\mu}-\gamma \mu$. Then for $\mu \geq 0$, we have $\Delta_{\gamma}(\mu, \varepsilon) \leq 0$ if and only if $\varepsilon \leq g_{\gamma}(\mu)$. Note that $g_{\gamma}^{\prime}(\mu)=$ $\frac{1}{\sqrt{\mu}}-\gamma \geq 0$ for $\mu \leq \mu_{1}$ if $\gamma \leq \frac{1}{\sqrt{\mu_{1}}}$. Then $\Delta_{\gamma}(\mu, \varepsilon) \leq 0$ for all $\mu \leq \mu_{1}$ if $\gamma \leq \frac{1}{\sqrt{\mu_{1}}}$ and $\varepsilon \leq g_{\gamma}\left(\mu_{n}\right)$. In particular, $\Delta_{\gamma}(\mu, \varepsilon) \leq 0$ for all $\mu \leq \mu_{1}$ if $\varepsilon=g_{\gamma}\left(\mu^{\prime}\right)$ for some $\mu^{\prime} \leq \mu_{n}$. We have

$$
\begin{aligned}
\alpha & =\max _{i} \frac{1}{2}\left[-\gamma \mu_{i}-\varepsilon+\Re\left(\sqrt{\left(\gamma \mu_{i}+\varepsilon\right)^{2}-4 \mu_{i}}\right)\right] \\
& =\max _{i} \frac{1}{2}\left[-\gamma \mu_{i}-\varepsilon\right] \\
& =\max _{i} \frac{1}{2}\left[-\gamma \mu_{i}-2 \sqrt{\mu^{\prime}}+\gamma \mu^{\prime}\right] \\
& =-\sqrt{\mu^{\prime}}-\frac{\gamma\left(\mu_{n}-\mu^{\prime}\right)}{2} .
\end{aligned}
$$

## Appendix C Proof of Proposition 2.4

We directly compute

$$
\begin{aligned}
\ddot{\boldsymbol{x}}= & -\boldsymbol{C}\left((\mathbb{I}-\gamma \varepsilon \boldsymbol{A}) \dot{\boldsymbol{p}}+\gamma \boldsymbol{A} \nabla^{2} f(\boldsymbol{x}) \dot{\boldsymbol{x}}\right)-\dot{\boldsymbol{C}}((\mathbb{I}-\gamma \varepsilon \boldsymbol{A}) \boldsymbol{p}+\gamma \boldsymbol{A} \nabla f(\boldsymbol{x})) \\
= & -\boldsymbol{C}\left((\mathbb{I}-\gamma \varepsilon \boldsymbol{A})(\boldsymbol{A} \nabla f(\boldsymbol{x})-\varepsilon \boldsymbol{A} \boldsymbol{p})+\gamma \boldsymbol{A} \nabla^{2} f(\boldsymbol{x}) \dot{\boldsymbol{x}}\right) \\
& -\dot{\boldsymbol{C}}((\mathbb{I}-\gamma \varepsilon \boldsymbol{A}) \boldsymbol{p}+\gamma \boldsymbol{A} \nabla f(\boldsymbol{x})) \\
= & -\boldsymbol{C}\left[(\mathbb{I}-\gamma \varepsilon \boldsymbol{A}) \boldsymbol{A} \nabla f(\boldsymbol{x})+\varepsilon \boldsymbol{A}\left(\boldsymbol{C}^{-1} \dot{\boldsymbol{x}}+\gamma \boldsymbol{A} \nabla f(\boldsymbol{x})\right)+\gamma \boldsymbol{A} \nabla^{2} f(\boldsymbol{x}) \dot{\boldsymbol{x}}\right] \\
& +\dot{\boldsymbol{C}} \boldsymbol{C}^{-1} \dot{\boldsymbol{x}} \\
= & -\boldsymbol{C}\left[\boldsymbol{A} \nabla f(\boldsymbol{x})+\varepsilon \boldsymbol{A} \boldsymbol{C}^{-1} \dot{\boldsymbol{x}}+\gamma \boldsymbol{A} \nabla^{2} f(\boldsymbol{x}) \dot{\boldsymbol{x}}\right]+\dot{\boldsymbol{C}} \boldsymbol{C}^{-1} \dot{\boldsymbol{x}} .
\end{aligned}
$$

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