# A first-order primal-dual method with adaptivity to local smoothness 

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#### Abstract

We consider the problem of finding a saddle point for the convex-concave objective $\min _{x} \max _{y} f(x)+\langle A x, y\rangle-g^{*}(y)$, where $f$ is a convex function with locally Lipschitz gradient and $g$ is convex and possibly non-smooth. We propose an adaptive version of the Condat-Vũ algorithm, which alternates between primal gradient steps and dual proximal steps. The method achieves stepsize adaptivity through a simple rule involving $\|A\|$ and the norm of recently computed gradients of $f$. Under standard assumptions, we prove an $\mathcal{O}\left(k^{-1}\right)$ ergodic convergence rate. Furthermore, when $f$ is also locally strongly convex and $A$ has full row rank we show that our method converges with a linear rate. Numerical experiments are provided for illustrating the practical performance of the algorithm.


## 1 Introduction

In this paper we study a particular instance of the composite minimization problem

$$
\begin{equation*}
\min _{x \in \mathcal{X}} f(x)+g(A x) \tag{1}
\end{equation*}
$$

where $f$ and $g$ are convex, proper and lower-semicontinuous (l.s.c.), and $A$ is a linear operator.
Problems of the form (1) have been studied in the literature under various assumptions on $f$ and $g$. For the particular instances where $g \circ A$ is proximal-friendly and $f$ is $L$-smooth, the objective is suitable for applying forward-backward splitting algorithms like the Proximal Gradient algorithm and its accelerated counterpart [Nesterov, 2013, Beck and Teboulle, 2009]. In general, however, the proximal operator of $g \circ A$ is not easily computable and in such cases a popular approach is to decouple $A$ and $g$ by reformulating problem (1) as a convex-concave saddle-point problem:

$$
\begin{equation*}
\min _{x \in \mathcal{X}} \max _{y \in \mathcal{Y}}\langle A x, y\rangle+f(x)-g^{*}(y) \tag{2}
\end{equation*}
$$

where $g^{*}$ denotes the Fenchel conjugate of $g$. Objective (2) is typically addressed by primal-dual splitting algorithms which, under strong duality, can recover the solution to the original problem (1). In the particular case when $f$ and $g$ are proximal-friendly and possibly non-smooth, a very popular method is the Primal-Dual Hybrid Gradient proposed in [Chambolle and Pock, 2011], which was further extended to handle an additional $L$-smooth component with the Condat-Vũ algorithm [Condat, 2013, Vũ, 2013]. Convergence rates for the latter are studied in [Chambolle and Pock, 2016a].
Together, these classes of algorithms cover a broad range of problems in diverse fields such as signal processing, machine learning, inverse problems, telecommunications and many others. As a result, a great amount of research effort has gone into addressing practical concerns such as robustness
to inexact oracles, acceleration and automation of stepsize selection. For a comprehensive list of examples and theoretical details we refer the reader to review papers [Combettes and Pesquet, 2011, Parikh and Boyd, 2014, Komodakis and Pesquet, 2015, Chambolle and Pock, 2016b]. In this work, we focus on the line of investigation studying stepsize regime automation for primal-dual algorithms targeting problem (2).
In their basic form, primal-dual methods require as input stepsize parameters belonging to a designated interval of stability, which depends on problem specific constants like the global smoothness parameter $L$ and $\|A\|$. Dependence on such constants is undesirable because they may be costly to compute and oftentimes one can only access upper-bound estimates, thus leading to overly-conservative stepsizes and slower convergence. Moreover, the need to know $L$ for setting the stepsizes prevents these methods from being applied to functions which are not globally smooth.

Consequently, recent efforts have gone towards devising methods with adaptive stepsizes [Goldstein et al., 2013, 2015, Malitsky and Pock, 2018, Pedregosa and Gidel, 2018]. These approaches resort to linesearch for finding good stepsizes at every iteration, and exhibit improved practical performance. It thus appears that better convergence comes at the cost of an indeterminate number of extra steps spent in subprocedures aimed at finding appropriate stepsizes.
In this work, we study problem (2) under the assumption that $\nabla f$ is locally Lipschitz continuous and $g$ is proximal-friendly. To illustrate the motivation of our framework, we take a prototypical example in image processing:

$$
\min _{x} \frac{1}{2}\|K x-b\|^{2}+\lambda\|D x\|_{2,1}, \quad K \in \mathbb{R}^{m \times d}, D \in \mathbb{R}^{2 d \times d}
$$

where $x$ is an image, $K$ is a problem-specific measurement operator, $b$ is the vector of (possibly noisy) observations and $D$ is the discrete gradient operator and the regularization term represents the isotropic TV norm. In order to apply any of the aforementioned primal-dual algorithms, one needs to first choose how to decouple the linear operators. There are three options: decoupling with respect to $K$ leaves us with having to compute the proximal operator of the TV norm for the primal step, which is an iterative procedure [Chambolle, 2004]. Decoupling $D$ implies performing gradient steps on $f$, since in general its proximal operator is not efficient. Finally, decoupling with respect to both implies increasing the dimensionality of the dual variable to $m+2 d$, which is problematic for large $d$ and $m$. The sensible choice is the second one (i.e., decoupling $D$ ), and the question we seek to answer with this work is:

## Does there exist a method for solving (2) that adapts to the local problem geometry without resorting to linesearch?

Our contribution is to propose a first-order primal-dual scheme that answers this question in the affirmative and is accompanied by theoretical convergence guarantees. Using standard analysis techniques we show an ergodic convergence of $\mathcal{O}\left(k^{-1}\right)$ when $\nabla f$ is locally Lipschitz and $g$ is proximal-friendly, and a linear convergence rate for the case when $f$ is in addition locally strongly convex and $A$ has full row rank. We provide numerical experiments for sparse logistic regression and image inpainting, as well as use our method as a heuristic for TV-regularized nonconvex phase retrieval.

The rest of the paper is structured as follows: Section 2 provides details about related work; Section 3 introduces notation, along with technical preliminaries and assumptions to be used in our analysis; Section 4 reports the main theoretical results alongside partial proofs; finally, partial numerical results are provided in Section 5 with the rest being deferred to the appendix due to lack of space.

## 2 Related Work

Adaptive Gradient Descent (GD) methods. Arguably the most widespread of optimization methods, GD presents similar shortcomings for setting the stepsize as those described in the previous section. In particular, much research effort has gone in devising variants of the algorithm that remove the need to estimate the global smoothness constant $L$. In a recent work, Malitsky and Mishchenko [2020] propose an extremely simple and effective alternative for setting the stepsize $\tau_{k}$ adaptively at every iteration, as follows:

$$
\begin{equation*}
\tau_{k}=\min \left\{\tau_{k-1} \sqrt{1+\frac{\tau_{k-1}}{\tau_{k-2}}}, \frac{\left\|x_{k}-x_{k-1}\right\|}{2\left\|\nabla f\left(x_{k}\right)-\nabla f\left(x_{k-1}\right)\right\|}\right\} \tag{3}
\end{equation*}
$$

Adaptivity essentially comes 'for free' in (3), as it involves solely quantities which have already been computed. Moreover, the method requires only the weaker assumption of local smoothness, thus extending the reach of provably-convergent GD to a wider class of differentiable functions while maintaining the standard $\mathcal{O}\left(k^{-1}\right)$ convergence rate.

In this work we show that the above technique can be extended to the analysis of primal-dual methods, where it gives rise to an algorithm whose stepsizes adapt to the local geometry of the objective's (locally) smooth component $f$.

Adaptive monotone variational inequality (VI) methods. Malitsky [2020] proposes an algorithm for solving monotone VIs with a stepsize that adapts to local smoothness similarly to (3). This method solves the very general formulation of finding $u^{*}$ such that $\left\langle F\left(u^{*}\right), u-u^{*}\right\rangle+h(u)-h\left(u^{*}\right) \geq 0, \forall u$ for a given monotone operator $F$ which is locally Lipschitz continuous. Our template (2) can be recovered from theirs by setting $u=(x, y)$, with

$$
F(u)=F(x, y)=\left[\begin{array}{c}
\nabla f(x)+A^{T} y \\
-A x
\end{array}\right]
$$

and $h(u)=g^{*}(y)$. The advantages of this approach are the relaxed requirement of local Lipschitz continuity for $F$ and the fact that knowledge of $\|A\|$ is not required. However, since the VI framework is very general and does not take advantage of the problem structure (e.g. the fact that $\langle A x, y\rangle$ is a bilinear term), the method comes with worse convergence bounds than algorithms specifically designed to solve (2). In addition, the algorithm requires as input an upper bound on the stepsizes, despite them being set in accordance to the estimated local smoothness.

First order primal dual algorithms and adaptive versions. A popular method for solving (2) when $f$ is $L$-smooth is the Condat-Vũ algorithm (CVA) [Condat, 2013, Vũ, 2013]. The method's convergence is subject to a global stepsize validity condition given by $\left(\frac{1}{\tau}-L\right) \frac{1}{\sigma} \geq\|A\|^{2}$, where $\tau$ and $\sigma$ are the primal and dual stepsizes, respectively.
Another approach to solving problem (2) is via the Primal-Dual Fixed-Point algorithm based on the Proximity Operator $\left(\mathrm{PDFP}^{2} \mathrm{O}\right)$ or the Proximal Alternating Predictor-Corrector (PAPC) methods [Loris and Verhoeven, 2011, Chen et al., 2013, Drori et al., 2015]. This approach comes with less restrictive stepsize conditions than CVA owing to a different iteration style, but which nevertheless depend on the global smoothness constant $L$ and $\|A\|$ and have to be carefully chosen.

In order to alleviate the burden of choosing the stepsize parameters in CVA, Malitsky and Pock [2018] propose a linesearch procedure involving only dual variable updates and which, for certain problems such as regularized least squares, does not require any additional matrix-vector multiplications. A characteristic of this algorithm is that it maintains a constant ratio between primal and dual stepsizes through a hyperparameter $\beta$ - a setup which we also use in this work.

## 3 Preliminaries

Consider problem (2) and let $\mathcal{X}, \mathcal{Y}$ be finite dimensional real vector spaces equipped with the standard inner product $\langle\cdot, \cdot\rangle$ and the associated Euclidean norm $\|\cdot\|=\sqrt{\langle\cdot, \cdot\rangle}$. We denote by $g^{*}$ the Fenchel conjugate of $g$ in (1) defined as $g^{*}(y)=\sup _{x}\{\langle x, y\rangle-g(x)\}$. In order to not overload the * notation, we use $A^{T}$ to denote the adjoint operator of $A$.
One can easily see that (2) is a primal-dual formulation of the following primal and dual optimization problems, of which the former is the same as (1):

$$
\min _{x \in \mathcal{X}} f(x)+g(A x), \quad \max _{y \in \mathcal{Y}}-\left(f^{*}\left(-A^{T} y\right)+g^{*}(y)\right)
$$

A saddle-point $\left(x^{*}, y^{*}\right) \in \mathcal{X} \times \mathcal{Y}$ of problem (2) satisfies the following optimality conditions:

$$
\begin{equation*}
-A^{T} y^{*}=\nabla f\left(x^{*}\right), \quad A x^{*} \in \partial g^{*}\left(y^{*}\right) \tag{4}
\end{equation*}
$$

For $\left(x^{\prime}, y^{\prime}\right) \in \mathcal{X} \times \mathcal{Y}$ we define the following quantities:

$$
P_{x^{\prime}, y^{\prime}}(x):=f(x)-f\left(x^{\prime}\right)+\left\langle x-x^{\prime}, A^{T} y^{\prime}\right\rangle
$$

$$
\begin{aligned}
& D_{x^{\prime}, y^{\prime}}(y):=g^{*}(y)-g^{*}\left(y^{\prime}\right)-\left\langle A x^{\prime}, y-y^{\prime}\right\rangle \\
& \mathcal{G}_{x^{\prime}, y^{\prime}}(x, y):=P_{x^{\prime}, y^{\prime}}(x)+D_{x^{\prime}, y^{\prime}}(y)
\end{aligned}
$$

These functions are convex for fixed $\left(x^{\prime}, y^{\prime}\right)$ and whenever $\left(x^{\prime}, y^{\prime}\right)=\left(x^{*}, y^{*}\right)$, it holds that $P_{x^{*}, y^{*}}(x) \geq 0, D_{x^{*}, y^{*}}(x) \geq 0$ and $\mathcal{G}_{x^{*}, y^{*}}(x, y) \geq 0$, with the latter quantity representing the primal-dual gap. We also define the gap restricted to a bounded subset $B_{1} \times B_{2} \subset \mathcal{X} \times \mathcal{Y}$ as:

$$
\mathcal{G}_{B_{1} \times B_{2}}(x, y):=\sup _{\left(x^{\prime}, y^{\prime}\right) \in B_{1} \times B_{2}} P_{x^{\prime}, y^{\prime}}(x)+D_{x^{\prime}, y^{\prime}}(y),
$$

and note that it is non-negative whenever $B_{1} \times B_{2}$ contains a saddle-point.
Given a function $f: \mathcal{X} \rightarrow \mathbb{R}$ and $L>0$, we say that $f$ is $L$-smooth if its gradient $\nabla f$ is Lipschitz continuous: $\|\nabla f(x)-\nabla f(y)\| \leq L\|x-y\|, \forall x, y$. Furthermore, $f$ is locally smooth if $\nabla f$ is Lipschitz continuous on any compact subset $\mathcal{C}: \forall \mathcal{C} \subset \mathcal{X}, \exists L_{\mathcal{C}}>0$ such that $\|\nabla f(x)-\nabla f(y)\| \leq$ $L_{\mathcal{C}}\|x-y\|, \forall x, y \in \mathcal{C}$.
We also say that $f$ is $\mu$-strongly convex if $f(y) \geq f(x)+\langle\nabla f(x), y-x\rangle+\frac{\mu}{2}\|x-y\|^{2}, \forall x, y$. Similarly, $f$ is locally strongly convex if it is strongly convex on any compact subset $\mathcal{C}: \forall \mathcal{C} \subset$ $\mathcal{X}, \exists \mu_{\mathcal{C}}>0$ such that $f(y) \geq f(x)+\langle\nabla f(x), y-x\rangle+\frac{\mu_{\mathcal{C}}}{2}\|x-y\|^{2}, \forall x, y \in \mathcal{C}$.
We define the proximal operator of a convex function $g: \mathcal{X} \rightarrow \mathbb{R} \cup\{\infty\}$ as $\operatorname{prox}_{g}(x)=$ $\arg \min _{z}\left\{g(z)+\frac{1}{2}\|x-z\|^{2}\right\}$, and say that $g$ is 'proximal-friendly' if $\operatorname{prox}_{g}(x)$ has a closed form solution or can be efficiently computed to high accuracy.
Finally, the following two blanket assumptions will hold throughout the paper:
Assumption 3.1. Function $f$ is convex and locally smooth, while $g$ convex, l.s.c., and proximalfriendly.
Assumption 3.2. A saddle-point exists for problem (2) and thus strong duality holds.
We note that Assumption 3.2 is standard in the literature (see e.g., [Chambolle and Pock, 2011]). Assumption 3.1, on the other hand, is weaker than the usual global $L$-smoothness premise and thus enlarges the category of admissible functions $f$ with instances such as $x \mapsto \exp (x)$. To illustrate, consider the aforementioned function defined on the reals: the global smoothness assumption clearly does not hold, however for any fixed interval $[a, b] \subset \mathbb{R}$ the smoothness constant can be chosen as $\exp (b)$.
For showing linear convergence of our method, we will add the following assumption:
Assumption 3.3. Function $f$ is locally strongly convex and operator $A$ has full row-rank.

## 4 Algorithm and convergence

The primal-dual method proposed for solving problem (2) under assumptions 3.1 and 3.2, is provided in Algorithm 1 under the abbreviation APDA, which we use from here onwards. APDA follows the same structure as the basic CVA [Chambolle and Pock, 2016a] for the given assumptions. Notice that if we restrict Assumption 3.1 to $L$-smooth functions $f$, we can in fact recover CVA by setting $\theta_{k}=\theta=1$ and $\tau_{k}=\tau, \sigma_{k}=\sigma$ fixed such that $\left(\frac{1}{\tau}-L\right) \frac{1}{\sigma} \geq\|A\|^{2}$.

```
Algorithm 1 Adaptive Primal Dual Algorithm (APDA)
    Input: \(x_{0} \in \mathcal{X}, y_{0} \in \mathcal{Y}, \tau_{\text {init }}>0, \tau_{0}=\infty, \theta_{0}=1, \beta>0, c \in(0,1)\)
    \(x_{1}=x_{0}-\tau_{\text {init }}\left(\nabla f\left(x_{0}\right)+A^{T} y_{0}\right)\)
    for \(k=1,2, \ldots\) do
        Set \(\tau_{k}=\min \left\{\frac{1}{2 \sqrt{L_{k}^{2}+(\beta /(1-c))\|A\|^{2}}}, \tau_{k-1} \sqrt{1+\theta_{k-1}}\right\}, \sigma_{k}=\beta \tau_{k}, \theta_{k}=\frac{\tau_{k}}{\tau_{k-1}}\)
        \(\tilde{x}_{k}=x_{k}+\theta_{k}\left(x_{k}-x_{k-1}\right)\)
        \(y_{k+1}=\operatorname{prox}_{\sigma_{k} g^{*}}\left(y_{k}+\sigma_{k} A \tilde{x}_{k}\right)\)
        \(x_{k+1}=x_{k}-\tau_{k}\left(\nabla f\left(x_{k}\right)+A^{T} y_{k+1}\right)\)
    end for
```


### 4.1 High level ideas

We can rephrase the global stepsize condition of CVA by introducing a free parameter $\beta>0$ which represents the ratio between the fixed dual and primal stepsizes: $\beta=\frac{\sigma}{\tau}$. With this change of variables, the stepsize validity condition becomes $\tau \in\left(0, \frac{2}{L+\sqrt{L^{2}+4 \beta\|A\|^{2}}}\right)$.
Our algorithm disposes of CVA's global condition and relies instead on a very similar but local criterion given by $\tau_{k} \in\left(0, \frac{1}{L_{k}+\sqrt{L_{k}^{2}+2 \beta\|A\|^{2}}}\right)$, where $L_{k}:=\frac{\left\|\nabla f\left(x_{k}\right)-\nabla f\left(x_{k-1}\right)\right\|}{\left\|x_{k}-x_{k-1}\right\|}$ provides an estimate of the local smoothness constant and $\beta=\frac{\sigma_{k}}{\tau_{k}}$. In particular, this requirement is satisfied by the first part of the expression defining $\tau_{k}$ in APDA:

$$
\begin{equation*}
\tau_{k}=\min \left\{\frac{1}{2 \sqrt{L_{k}^{2}+(\beta /(1-c))\|A\|^{2}}}, \tau_{k-1} \sqrt{1+\theta_{k-1}}\right\} \tag{5}
\end{equation*}
$$

where $c \in(0,1)$. Intuitively, this rule demands that $\tau_{k}$ does not overstep a constant related to the local curvature, thus allowing for larger stepsizes in flatter regions and correspondingly smaller ones otherwise.
By itself, the first term of (5) does not ensure convergence, since overly-aggressive and possibly destabilizing stepsizes might occur in near-linear regions. This issue is addressed by the second part of the expression (5) which, informally, prevents the stepsize from increasing 'too fast' in consecutive iterations. Specifically, the increase factor is at most $\sqrt{1+\theta_{k-1}}$, where $\theta_{k}=\frac{\tau_{k-1}}{\tau_{k-2}}$.

Under these two local stepsize conditions we are able to show APDA's convergence using the weaker assumption of local smoothness of $f$, thus conveniently removing the need of estimating a global smoothness constant $L$.
Remark 4.1. While $\tau_{k}$ does not adapt to $\|A\|$, for many practical problems this fact is not a big hindrance. Function $f$ typically represents the data fidelity term, whose smoothness constant $L$ (should it exist) can far exceed $\|A\|$ - the linear operator enforcing structured regularization on $x$. A specific example are TV-regularized imaging problems, where $A$ is the discrete gradient operator whose norm is bounded by $\sqrt{8}$ [Chambolle, 2004], while the data fidelity term may involve a very large number of measurements and a larger norm, consequently.
Remark 4.2. APDA takes an additional primal step prior to the for-loop, which is controlled by $\tau_{\text {init }}$ given as input. This is needed for estimating $L_{1}$ in the first iteration. In practice we set $\tau_{\text {init }}=1 e-9$, a sufficiently small value to ensure that $x_{1}$ does not depart too far from $x_{0}$ and yield a good estimate of $L_{1}$. Furthermore, the setting of $\tau_{0}=\infty$ simply ensures that in the first step, $\tau_{1}=\frac{1}{2 \sqrt{L_{1}^{2}+(\beta /(1-c))\|A\|^{2}}}$ and has no impact on further steps. Finally, in our experiments we set $c=1 e-15-$ this is a parameter introduced for theoretical purposes as explained in the following section.

### 4.2 Analysis - the base case

In short, the main steps of our analysis are: first, we establish the inequality that characterizes the dynamics of APDA given in Lemma 4.1 below. Based on it, we are able to prove the boundedness of sequences $\left\{x_{k}\right\}$ and $\left\{y_{k}\right\}$ in Theorem 4.1. In turn, sequence boundedness alongside the local smoothness property of $f$ allows us to conclude that there exists a constant $L>0$ such that $f$ is $L$-smooth on the compact set $\overline{\operatorname{Conv}}\left(\left\{x^{*}, x_{0}, x_{1}, \ldots\right\}\right)$ - the closed convex hull generated by $\left\{x^{*}, x_{0}, x_{1}, \ldots\right\}$. Finally, we leverage this information to show that $\left(x_{k}, y_{k}\right)$ converges to a saddle point of (2) and derive the associated ergodic convergence rates presented in Theorem 4.1.
Lemma 4.1. Consider APDA along with Assumptions 3.1 and 3.2 and $(x, y) \in \mathcal{X} \times \mathcal{Y}$. Then, for all $k$ and $\eta_{k} \in\left(\frac{\beta \tau_{k}\|A\|}{1-c}, \frac{1-2 \tau_{k} L_{k}}{2 \tau_{k}\|A\|}\right)$,

$$
\begin{aligned}
\left\|x_{k+1}-x\right\|^{2} & +\frac{1}{\beta}\left\|y_{k+1}-y\right\|^{2}+\left(1-\eta_{k} \tau_{k}\|A\|-\tau_{k} L_{k}\right)\left\|x_{k+1}-x_{k}\right\|^{2} \\
& +\frac{\eta_{k}-\tau_{k} \beta\|A\|}{\beta \eta_{k}}\left\|y_{k+1}-y_{k}\right\|^{2}+2 \tau_{k}\left(1+\theta_{k}\right) P_{x, y}\left(x_{k}\right)+2 \tau_{k} D_{x, y}\left(y_{k+1}\right)
\end{aligned}
$$

$$
\leq\left\|x_{k}-x\right\|^{2}+\frac{1}{\beta}\left\|y_{k}-y\right\|^{2}+\tau_{k} L_{k}\left\|x_{k}-x_{k-1}\right\|^{2}+2 \tau_{k} \theta_{k} P_{x, y}\left(x_{k-1}\right)
$$

Moreover, it holds that:

1) $\tau_{k} L_{k}<\frac{1}{2}<1-\eta_{k} \tau_{k}\|A\|-\tau_{k} L_{k}$,
2) $\frac{1}{\beta}-\frac{\tau_{k}\|A\|}{\eta_{k}}>\frac{c}{\beta}>0$.

Proof sketch. The full proof is deferred to the appendix. We use algebraic manipulations, APDA's update rules, the Cauchy-Schwarz and Young inequalities and properties of the prox operator to get the recurrence:

$$
\begin{align*}
&\left\|x_{k+1}-x\right\|^{2}+\frac{1}{\beta}\left\|y_{k+1}-y\right\|^{2}+\left(1-\tau_{k}\|A\| \eta_{k}-\tau_{k} L_{k}\right)\left\|x_{k+1}-x_{k}\right\|^{2} \\
&+\left(\frac{1}{\beta}-\frac{\tau_{k}\|A\|}{\eta_{k}}\right)\left\|y_{k+1}-y_{k}\right\|^{2}+2 \tau_{k}\left(1+\theta_{k}\right) P_{x, y}\left(x_{k}\right)+2 \tau_{k} D_{x, y}\left(y_{k+1}\right) \\
& \leq\left\|x_{k}-x\right\|^{2}+\frac{1}{\beta}\left\|y_{k}-y\right\|^{2}+\tau_{k} L_{k}\left\|x_{k}-x_{k-1}\right\|^{2}+2 \tau_{k} \theta_{k} P_{x, y}\left(x_{k-1}\right) \tag{6}
\end{align*}
$$

where $\eta_{k}>0$ is a free iteration-dependent constant involved in Young's inequality.
In order to obtain anything worthwhile we would like to set $\eta_{k}$ such that, when unrolling (6) over the iterations, the terms containing $\left\|x_{k+1}-x_{k}\right\|^{2}$ and $\left\|y_{k+1}-y_{k}\right\|^{2}$ accumulate on the LHS with positive coefficients. More precisely, we ask that:

$$
\left\{\begin{array}{l}
\frac{1}{\beta}-\frac{\tau_{k}\|A\|}{\eta_{k}}>\frac{c}{\beta}  \tag{7}\\
1-\tau_{k}\|A\| \eta_{k}-\tau_{k} L_{k}>\frac{1}{2}
\end{array}\right.
$$

where $c \in(0,1)$. We note that the RHS of the first inequality could have been chosen as 0 , however, we made it strictly positive due to technical reasons related to controlling the sequence $\left\|y_{k+1}-y_{k}\right\|^{2}$. In practice, we choose $c$ to be as small as possible.
A similar remark holds for the second inequality, where it would have been sufficient to set its RHS to $\tau_{k+1} L_{k+1}$. Since this would considerably complicate the analysis, we make the observation that $\tau_{k} L_{k}<\frac{1}{2}, \forall k$ and use this simpler uniform upper-bound instead.

The inequalities (7) are equivalent to asking that $\eta_{k} \in\left(\frac{\tau_{k} \beta\|A\|}{1-c}, \frac{1-2 \tau_{k} L_{k}}{2 \tau_{k}\|A\|}\right)$ and what is left to show is that this is a valid interval i.e., that the left endpoint is strictly smaller than its right counterpart. This condition amounts to solving a quadratic inequality in $\tau_{k}$, whose solutions lie in the interval $\left(0, \frac{1}{L_{k}+\sqrt{L_{k}^{2}+2(\beta /(1-c))\|A\|^{2}}}\right)$. The proof is concluded by showing that our choice of $\tau_{k}$ indeed satisfies this constraint.

We are now ready to state the main convergence result in Theorem 4.1 below, whose full proof is given in the appendix.
Theorem 4.1. Consider APDA along with Assumptions 3.1 and 3.2, and let $\left(x^{*}, y^{*}\right) \in \mathcal{X} \times \mathcal{Y}$ be a saddle point of problem (2). Then, for all $k$

1) Boundedness. The sequence $\left\{\left(x_{k}, y_{k}\right)\right\}$ is bounded. Specifically, for all $k$,

$$
\left\|x_{k}-x^{*}\right\|^{2}+\left\|y_{k}-y^{*}\right\|^{2} \leq M
$$

where $M:=\left\|x_{1}-x^{*}\right\|^{2}+\frac{1}{\beta}\left\|y_{1}-y^{*}\right\|^{2}+\frac{1}{2}\left\|x_{1}-x_{0}\right\|^{2}<\infty$.
2) Convergence to a saddle point. The sequence $\left\{\left(x_{k}, y_{k}\right)\right\}$ converges to a saddle point of (2).
3) Ergodic convergence. Let $S_{k}:=\sum_{i=1}^{k} \tau_{i}, \quad X_{k}:=\frac{1}{S_{k}}\left(\tau_{k}\left(1+\theta_{k}\right) x_{k}+\right.$ $\left.\sum_{i=1}^{k-1}\left(\tau_{i}\left(1+\theta_{i}\right)-\tau_{i+1} \theta_{i+1}\right) x_{i}\right)$ and $Y_{k}:=\frac{1}{S_{k}} \sum_{i=1}^{k} \tau_{i} y_{i+1}$. Then, for any bounded
$B_{1} \times B_{2} \in \mathcal{X} \times \mathcal{Y}$ and for all $k$,

$$
\mathcal{G}_{B_{1} \times B_{2}}\left(X_{k}, Y_{k}\right) \leq \frac{M\left(B_{1}, B_{2}\right) \sqrt{L^{2}+(\beta /(1-c))\|A\|^{2}}}{k}
$$

where $L$ is the Lipschitz constant of $\nabla f$ over the compact set $\overline{\operatorname{Conv}}\left(\left\{x^{*}, x_{0}, x_{1}, \ldots\right\}\right)$ and $M\left(B_{1}, B_{2}\right)=\sup _{(x, y) \in B_{1} \times B_{2}}\left\|x_{1}-x\right\|^{2}+\frac{1}{\beta}\left\|y_{1}-x\right\|^{2}+\frac{1}{2}\left\|x_{1}-x_{0}\right\|^{2}$.

The boundedness result of Theorem 4.1 point 1) implies that the closed set $\mathcal{C}=$ $\overline{\operatorname{Conv}}\left(\left\{x^{*}, x_{0}, x_{1}, \ldots\right\}\right)$ is also bounded and hence compact. The local smoothness assumption on $f$ then ensures that there exists $L>0$ such that $f$ is $L$-smooth over $\mathcal{C}$. Note that such an $L$ exists for any $x_{0}, y_{0}$ since the boundedness result itself holds for any initial conditions (though the value of such $L$ cannot be generally known, as it is path-dependent). Using this fact, we can show a uniform lower-bound on the primal stepsize: $\tau_{k} \geq \frac{1}{2}\left(L^{2}+(\beta /(1-c))\|A\|^{2}\right)^{-1 / 2}>0, \forall k$, which is instrumental in deriving the subsequent convergence results, as well as Theorem 4.2. We emphasize that the appearance of constant $L$ in the provided rates is a consequence of iterate boundedness, whose proof does not require its knowledge. Finally, we note that our rate is comparable to that of CVA in terms of constants.

### 4.3 Analysis under the additional Assumption 3.3

We now study APDA under the additional assumption of locally strongly convex $f$ and full row rank A. Before proving the result of Theorem 4.2, a few remarks are in order. First, the boundedness result of Theorem 4.1 point 1) also holds for constant $c=0$, since this constant was required only for proving convergence to a saddle point in point 2) of the theorem. Second, taking a smaller stepsize than the originally defined $\tau_{k}$ will not change the validity of Lemma 4.1 or the boundedness result of Theorem 4.1, as it remains within the required interval mentioned in section 4.1.
Consequently, for studying APDA under the additional Assumption 3.3 we can simplify the stepsize expression by taking $c=0$, because now we are able to show iterate convergence directly by using the strong convexity and full row-rank assumptions. Specifically, we consider the stepsize:

$$
\begin{equation*}
\tau_{k}=\min \left\{\frac{1}{2 \sqrt{4 L_{k}^{2}+\beta\|A\|^{2}}}, \tau_{k-1} \sqrt{1+\theta_{k-1} / 2}\right\} \tag{8}
\end{equation*}
$$

which is smaller than the one originally considered and, due to the aforementioned remarks it ensures that APDA produces a bounded sequence. It follows that, under the local smoothness and local strong convexity assumptions, there exist constants $L$ and $\mu$ such that $f$ is $L$-smooth and $\mu$-strongly convex over $\overline{\operatorname{Conv}}\left(\left\{x^{*}, x_{0}, x_{1}, \ldots\right\}\right)$.
The existence of these constants along with $A$ being full row rank, in turn, allows us to derive a strengthened version of the inequality in Lemma 4.1 for $(x, y)=\left(x^{*}, y^{*}\right)$ :

$$
\begin{aligned}
&\left\|x_{k+1}-x^{*}\right\|^{2}+\left(\frac{1}{\beta}+q_{1}\right)\left\|y_{k+1}-y^{*}\right\|^{2}+\left(\frac{1}{2}+q_{2}\right)\left\|x_{k}-x_{k+1}\right\|^{2}+q_{3}\left\|y_{k+1}-y_{k}\right\|^{2} \\
&+2 \tau_{k}\left(1+\theta_{k}\right) P_{x^{*}, y^{*}}\left(x_{k}\right)+2 \tau_{k} D_{x^{*}, y^{*}}\left(y_{k+1}\right) \\
& \leq\left(1-q_{4}\right)\left\|x_{k}-x^{*}\right\|^{2}+\frac{1}{\beta}\left\|y_{k}-y^{*}\right\|^{2}+\left(\frac{1}{2}-q_{5}\right)\left\|x_{k}-x_{k-1}\right\|^{2}+2 \tau_{k} \theta_{k} P_{x^{*}, y^{*}}\left(x_{k-1}\right)
\end{aligned}
$$

where $q_{1}, q_{2}, q_{3}, q_{4}, q_{5}>0$ are constants given in the appendix. This new inequality represents in fact a contraction, which guarantees the linear convergence rate stated in Theorem 4.2.

Theorem 4.2. Consider APDA along with Assumptions 3.1, 3.2 and 3.3. Let $\left(x^{*}, y^{*}\right) \in \mathcal{X} \times \mathcal{Y}$ be a saddle point of problem (2). Furthermore, let $\tau_{k}$ be defined by (8) and let $s:=\sqrt{4 L^{2}+\beta\|A\|^{2}}$ and $t:=\sqrt{4 \mu^{2}+\beta\|A\|^{2}}$, where $\mu, L$ are the strong convexity and smoothness constants of $f$ over the compact set $\overline{\operatorname{Conv}}\left(\left\{x^{*}, x_{0}, x_{1}, \ldots\right\}\right)$.

Then, for all $k$ :

$$
\left\|x_{k}-x^{*}\right\|^{2}+\frac{1}{\beta}\left\|y_{k}-y^{*}\right\|^{2} \leq(1-\min \{p, q, r\})^{k} M
$$

where the rate constants are given by:

$$
p=\frac{1}{2}, \quad q=\frac{\mu}{4 s}, \quad r=\frac{\beta \sigma_{\min }^{2}(A) \mu}{\beta \sigma_{\min }^{2}(A) \mu+8 s^{2} t+4 L^{2} s}
$$

and $M=\left\|x_{2}-x^{*}\right\|^{2}+\left(\frac{1}{\beta}+T\right)\left\|y_{2}-y^{*}\right\|^{2}+\frac{1}{2}\left\|x_{2}-x_{1}\right\|^{2}+2 \tau_{1} P_{x^{*}, y^{*}}\left(x_{1}\right), \quad T=$ $\frac{\sigma_{\min }^{2}(A) \mu}{8 s^{2} t+4 L^{2} s}$, with $\sigma_{\min }(A)$ representing the smallest singular value of $A$.

A few remarks are in order: first, as a sanity check, we observe that when $A=0$ we recover the contraction factor of [Malitsky and Mishchenko, 2020] which is equal to $q$.

Second, we make some notes on how our rate compares with existing ones. To our knowledge, there are no explicit results regarding the linear convergence of CVA under assumptions similar to ours (linear rates are usually shown for the 3 -component objective without assumptions on $A$ - see e.g., [Chambolle and Pock, 2016a]). However, in the case of $L$-smooth and $\mu$-strongly-convex $f$ and full row-rank $A$, Chen et al. [2013] show the linear convergence of $\mathrm{PDFP}^{2} \mathrm{O}$ with rate:

$$
\left\|x_{k}-x^{*}\right\|^{2} \leq\left(\left\|x_{1}-x_{0}\right\|^{2}+\frac{1}{\sigma_{\max (A)}}\left\|y_{1}-y_{0}\right\|^{2}\right)\left(1-\min \left\{\frac{\sigma_{\min (A)}^{2}}{\sigma_{\max (A)}^{2}}, \frac{\mu}{L}\right\}\right)^{k-1}
$$

The rate presented in Theorem 4.2 has a comparatively worse contraction factor. The reason is that our iteration is set up in the style of CVA, where we essentially have a single stepsize to compute using the rephrasing from Section 4.1. Therefore, $\tau_{k}$ needs to obey the problem structure with respect to both $L$ and $\|A\|$, resulting in the 'mixed' term appearing in the denominator.

Keeping the above in mind, the interested reader may find in the appendix that constants $q$ and $r$ come from a product between $\tau_{k}$ and other condition number-related quantities, which is tightly liked to the structure of the main inequality used in the paper. This makes the nice separation of condition numbers achieved in $\mathrm{PDFP}^{2} \mathrm{O}^{\prime}$ 's rate not possible in our case and, it seems, the analysis necessary to achieve this kind of adaptivity comes at the cost of worse constants (the same remark holds for [Malitsky and Mishchenko, 2020]).
$\mathrm{PDFP}^{2} \mathrm{O}$, on the other hand, achieves a clean bound by having a different iteration style than CVA, as well as a fundamentally different kind of analysis where the iteration is expressed in fixed-point form to show convergence. In this context the stability conditions on the stepsizes are also relaxed - specifically, $0<\lambda \leq 1 / \sigma_{\max }^{2}(A)$ and $0<\gamma<2 L$ in [Chen et al., 2013]. A drawback of this approach, however, is that the algorithm has no rate guarantees when $f$ is only smooth and not strongly convex and only asymptotic convergence is shown. Also, $\mathrm{PDFP}^{2} \mathrm{O}$ requires 3 matrix-vector multiplications per iteration whereas we only require 2 .

## 5 Experiments

We now present some numerical experiments conducted for APDA ${ }^{1}$. Additional problems and results are included the appendix. The experiments were implemented in Python 3.9 and executed on a MacBook Pro with 32 GB RAM and a 2,9 GHz 6-Core Intel Core i9 processor.

[^0]The baseline we compare against in this section as well as the appendix is CVA, for which we use Algorithm 1 in [Chambolle and Pock, 2016a] (using $g \equiv 0$ ). In the particular case of sparse logistic regression we also compare against FISTA [Beck and Teboulle, 2009]. For obtaining $x^{*}$ we ran one of the algorithms for a large number of iterations.

### 5.1 Sparse binary logistic regression



Figure 1: The first column shows algorithm convergence. The second column shows a comparison of primal stepsizes between APDA and CVA. The third column shows a comparison of dual stepsizes between APDA and CVA. Each subfigure represents a different dataset: (a) ijcnn; (b) mushrooms; (c) a9a; (d)covtype.

We consider the problem of sparse binary Logistic Regression on 4 LIBSVM datasets [Chang and Lin, 2011] and show that adaptivity provides faster convergence in 3 of these cases. The objective we
consider is:

$$
\begin{equation*}
\min _{x \in \mathbb{R}^{d}} F(x):=\underbrace{\sum_{i=1}^{m} \log \left(1+\exp \left(-b_{i}\left\langle q_{i}, x\right\rangle\right)\right)}_{f}+\underbrace{\lambda\|x\|_{1}}_{g}, \tag{9}
\end{equation*}
$$

where $\left(q_{i}, b_{i}\right) \in \mathbb{R}^{d} \times\{-1,1\}$ and $\lambda$ is the regularization parameter. APDA and CVA can be applied to this problem by setting $A=I$ in formulation (2). Primal-dual algorithms are not the typical choice for solving (9), which is usually addressed by methods such as Proximal Gradient or FISTA [Beck and Teboulle, 2009]. However, we note that the computational costs of APDA and FISTA are comparable since the matrix-vector multiplication cost of the former is removed due to a $A=I$.
We choose $\lambda=0.005\left\|Q^{T} b\right\|_{\infty}$, where $Q^{T}=\left[q_{1}^{T}, \ldots q_{m}^{T}\right]^{T}$. For APDA we perform a parameter sweep over $\beta \in[1 \mathrm{e}-3,1 \mathrm{e} 6]$ for each dataset and settle for: $\beta=2.68 \mathrm{e} 3$ for $\mathrm{ijcnn} ; \beta=5.18 \mathrm{e} 4$ for a9a; $\beta=3.16 \mathrm{e} 1$ for mushrooms; $\beta=3.73 \mathrm{e}-1$ for covtype.
For CVA we sweep $p \in[1 \mathrm{e}-3,1 \mathrm{e} 6]$ and set $\tau=\frac{1}{\|A\| / p+L}$ and $\sigma=\frac{1}{p\|A\|}$ - by construction, these stepsizes satisfy the validity condition and are as large as possible since the condition is satisfied with equality. We do an additional tuning procedure where we choose constants $\tau \in[1 \mathrm{e}-10,1 \mathrm{e} 2]$ and $\xi \in[1 \mathrm{e}-5,1 \mathrm{e} 2]$ and set $\sigma=\tau \xi$, which are subject to verifying the stepsize validity condition of CVA. Finally we select the best stepsizes across the two tuning phases to be (truncated to 3 decimals): $\tau=9.869 \mathrm{e}-4, \sigma=1.125 \mathrm{e} 1$ for $\mathrm{ijcnn} ; \tau=2.655 \mathrm{e}-4, \sigma=7.896 \mathrm{e} 1$ for a9a; $\tau=9.936 \mathrm{e}-4$, $\sigma=5.878 \mathrm{e} 0$ for mushrooms; $\tau=7.728 \mathrm{e}-06, \sigma=1 \mathrm{e}-06$ for covtype.
Note that the Hessian of $f$ is given by $\nabla^{2} f(x)=Q^{T} D(x) Q$, where $D(x)$ is a diagonal matrix such that $D_{i, i}(x)=\sigma_{i}(x)\left(1-\sigma_{i}(x)\right)$, where $\sigma_{i}(x)=\frac{1}{1+\exp \left(-b_{i}\left\langle q_{i}, x\right\rangle\right)} \in(0,1)$. Clearly, over any compact set in $\mathcal{C} \subset \mathcal{X}$ there exist $D_{\text {min }}:=\min _{i, x \in \mathcal{C}} D_{i, i}(x) \in(0,1)$ such that $D_{\min } Q^{T} Q \preceq Q^{T} D(x) Q$. As a result, a sufficient condition for local strong convexity is that the minimum eigenvalue of $Q^{T} Q$ be greater than 0 .
The convergence results along with stepsize comparison plots are presented in Figure 1. For dataset ij cnn we run APDA with the modified $\tau_{k}$ used in Theorem 4.2, since $\lambda_{\min }\left(Q^{T} Q\right)=75.13$ and $A$ has full rank. In the latter case, the legend identifier is APDA-strcnv. For the remaining datasets we use only the basic setting for $\tau_{k}$, as $\lambda_{\min }\left(Q^{T} Q\right) \leq 1 \mathrm{e}-13$.
While APDA outperforms FISTA and CVA on ijcnn, a9a and mushrooms, it shows a relatively poor performance on covtype. We hypothesize that this is related to the condition number of $Q^{T} Q$, which is almost three orders of magnitude larger in the latter case: 9.2 e 22 versus $5.3 \mathrm{e} 1,2 \mathrm{e} 20$ and 2 e 17 for ijcnn, mushrooms and a9a, respectively. A similar behavior is seen in Figure 1.(c) of [Malitsky and Mishchenko, 2020].
Finally, the adaptive property of APDA's stepsizes is visible in the stepsize comparison plots where they are shown to oscillate within at least one order of magnitude throughout the optimization process.

## 6 Limitations of APDA

The experiments presented in this paper (Section 5 and Appendix A) have the common trait of not imposing hard constraints on the primal variables. As a consequence, we are able to take plain gradient steps in the primal domain. However, for instances such as Poisson linear inverse problems [Bertero et al., 2009], the iterates $x_{k}$ need to reside in $\mathbb{R}_{+}^{n}$ because the primal objective contains log functions. APDA cannot handle such cases, as any constraints imposed on the primal variables will only be satisfied asymptotically. We consider such scenarios as a future research direction.

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## Appendix

## A Additional experiments

## A. 1 Nonconvex phase retrieval



Figure 2: (a) Convergence rate. (b) Primal stepsize comparison. (c) Primal stepsize comparison. (d) APDA reconstruction, $\mathrm{PSNR}=21.34, \mathrm{SSIM}=0.76$. (e) CVA reconstruction, $\mathrm{PSNR}=20.56, \mathrm{SSIM}=$ 0.70.

In this section we provide the results for applying our algorithm, heuristically, on the nonconvex least squares formulation of the phase retrieval (PR) problem. The phase-retrieval problem has attracted intense interest recently, due to its application is domains such as optical imaging [Walther, 1963], astronomy [Fienup and Dainty, 1987] and many others. Here, we consider the real counterpart of the original complex PR formulation for square images, where given $\left\{\left(A_{i}, b_{i}\right) \in \mathbb{R}^{n \times n} \times \mathbb{R}\right\}$ we want to recover $X \in \mathbb{R}^{n \times n}$ up to its sign, such that $b_{i}=\operatorname{Tr}\left(A_{i}^{T} x\right)^{2}$. To this end, we consider the following optimization objective:

$$
\begin{equation*}
\min _{X \in \mathbb{R}^{n \times n}} F(X):=\underbrace{\frac{1}{4 m} \sum_{i=1}^{m}\left(b_{i}-\operatorname{Tr}\left(A_{i}^{T} X\right)^{2}\right)^{2}}_{f(X)}+\underbrace{\lambda\|D X\|_{2,1}}_{g(X) \equiv\|\cdot\|_{\mathrm{Tv}}} . \tag{10}
\end{equation*}
$$

We note a few things: first, objective (10) is nonconvex with $f$ being only locally smooth. Secondly, Sun et al. [2018] have recently shown that given $m$ i.i.d Gaussian measurements, the global geometry of $F(X)$ is 'benign' for $m>C d \log (d)^{3}$, where $d$ is the problem dimension. By benign, the authors specifically mean '(1) there are no spurious local minimizers, and all global minimizers are equal to the target signal $x$ up to a global phase; and (2) the objective function has a negative directional curvature around each saddle point'. It is posed that in such cases iterative algorithms should, with high probability, find the minimizer without requiring special initialization as is needed for current state of the art solvers.

For our experiments we use $84 \times 84$-sized images and choose a smaller number of measurements than suggested above: $m=d \log (d) \approx 27,155$. We generate $m$ sparse matrices $A_{i} \in R^{n \times n}$ with $30 \%$ non-zero entries sampled i.i.d from the standard normal distribution, and corrupt a random subset containing $10 \%$ of elements in $b_{i}$ by setting them to 0 . We perform parameter sweep for $\lambda \in[1 \mathrm{e}-4,1 \mathrm{e} 4], \beta \in[1 \mathrm{e}-3,1 \mathrm{e} 4]$ and settle for $\lambda=1 \mathrm{e} 2$ and $\beta=2.78 \mathrm{e} 2$. Without guidelines for setting $\tau, \sigma$ for CVA since $f$ is not $L$-smooth, we search for the best $\tau \in[1 \mathrm{e}-4,1 \mathrm{e} 4]$ and
$p \in[1 \mathrm{e}-2,1 \mathrm{e} 2]$ such that $\sigma=\frac{1}{p \tau\|A\|}$ and settle for $\tau=1 \mathrm{e}-4, p=1.02$. We note that CVA diverged for $32 / 40$ grid points, whereas our method converged for all instances. Finally, the initial points $x_{0}$ and $y_{0}$ are sampled from the standard normal distribution.
The results are depicted in Figure 2, which contains the reconstructions and convergence plots. For each reconstruction we report the Peak Signal to Noise Ratio (PSNR) and the Structural Similarity Index Measure (SSIM). We tested several random seeds and obtained similar results. We also tried running CVA with the stepsize values used by APDA in its last iteration (notice how in Figure 2 (d) $\tau_{k}$ essentially stabilizes in a very narrow band just above $1 \mathrm{e}-3$ after the first 250 iterations) however, CVA diverged in this setting as well.

## A. 2 Image inpainting



Figure 3: (a) Original image downloaded from http://www.cs.tut.fi/~foi/GCF-BM3D/. (b) APDA reconstruction, $\mathrm{PSNR}=25.63, \mathrm{SSIM}=0.91$. (c) CVA reconstruction, $\mathrm{PSNR}=25.63, \mathrm{SSIM}=$ 0.91. (d) Convergence rate. (e) Primal stepsize comparison. (f) Dual stepsize comparison.

Image inpainting consists in reconstructing the missing parts of a subsampled image $b=P_{\Omega}\left(X^{\natural}\right)$, where $P_{\Omega}: \mathbb{R}^{m \times n} \rightarrow \mathbb{R}^{m \times n}$ is an operator that selects a subset of $q$ pixels from the original image $X^{\natural} \in \mathbb{R}^{m \times n}$, where $q \ll m n$. This problem can be formulated as a regularized optimization objective:

$$
\begin{equation*}
\min _{X \in \mathbb{R}^{m \times n}} F(X):=\underbrace{\frac{1}{2}\left\|b-P_{\Omega}(X)\right\|_{2}^{2}}_{f(X)}+\underbrace{\lambda\|D X\|_{2,1}}_{g(X) \equiv\|\cdot\|_{\mathrm{TV}}}, \tag{11}
\end{equation*}
$$

where $D: \mathbb{R}^{m \times n} \rightarrow \mathbb{R}^{m \times n \times 2}$ represents the discrete gradient operator, and $\|D X\|_{2,1}=$ $\sum_{i, j=1}^{m, n} \sqrt{(D X)_{i, j, 1}^{2}+(D X)_{i, j, 2}^{2}}$. The regularization term represents the isotropic TV norm, which is known to help in recovering sharp images by preserving discontinuities and reducing noise [Chambolle et al., 2010, Condat, 2017].

For our experiments, we vectorize the images of size $256 \times 256$ and transform $D$ accordingly. We represent $P_{\Omega}$ as a matrix built by removing rows uniformly at random from $I$ and which removes $60 \%$ of pixels from the original image (sampling ratio 0.4 ). We perform parameter sweep for $\lambda \in[1 \mathrm{e}-4,1 \mathrm{e}]$, and settle for $\lambda=1 \mathrm{e}-2$. We also sweep $\beta \in[1 \mathrm{e}-5,1 \mathrm{e}]$ and settle for $\beta=$
1.291e-2. Finally, we perform a similar two-phase tuning for CVA as that described in Section 5 with $p \in[1 \mathrm{e}-5,1 \mathrm{e} 3]$ for the first phase and $\tau \in[1 \mathrm{e}-5,1 \mathrm{e} 2], \xi \in[1 \mathrm{e}-5,1 \mathrm{e} 1]$ for the second phase. We settle for stepsizes $\tau=8.722 \mathrm{e}-1$ and $\sigma=1.831 \mathrm{e}-2$.

Experiment results are presented in Figure 3, where we show the reconstructions, alongside the convergence plot and a comparison of the fixed stepsizes of CVA with those of APDA. The two algorithms are comparable both in terms of reconstruction quality and convergence speed, with APDA being marginally better for the latter criterion. The convergence plot also shows an instance of CVA whose stepsizes were set to the values of those used by APDA in the final iteration of these experiments. Finally, subfigures (e) and (f) show APDA's stepsizes oscillating within close range of CVA's.

## B Missing proofs

## B. 1 Proof of Lemma 4.1

Lemma 4.1. Consider APDA along with Assumptions 3.1 and 3.2 and $(x, y) \in \mathcal{X} \times \mathcal{Y}$. Then, for all $k$ and $\eta_{k} \in\left(\frac{\beta \tau_{k}\|A\|}{1-c}, \frac{1-2 \tau_{k} L_{k}}{2 \tau_{k}\|A\|}\right)$,

$$
\begin{aligned}
& \left\|x_{k+1}-x\right\|^{2}+\frac{1}{\beta}\left\|y_{k+1}-y\right\|^{2}+\left(1-\eta_{k} \tau_{k}\|A\|-\tau_{k} L_{k}\right)\left\|x_{k+1}-x_{k}\right\|^{2} \\
& +\frac{\eta_{k}-\tau_{k} \beta\|A\|}{\beta \eta_{k}}\left\|y_{k+1}-y_{k}\right\|^{2}+2 \tau_{k}\left(1+\theta_{k}\right) P_{x, y}\left(x_{k}\right)+2 \tau_{k} D_{x, y}\left(y_{k+1}\right) \\
& \quad \leq\left\|x_{k}-x\right\|^{2}+\frac{1}{\beta}\left\|y_{k}-y\right\|^{2}+\tau_{k} L_{k}\left\|x_{k}-x_{k-1}\right\|^{2}+2 \tau_{k} \theta_{k} P_{x, y}\left(x_{k-1}\right)
\end{aligned}
$$

Moreover, it holds that:

$$
\text { 1) } \tau_{k} L_{k}<\frac{1}{2}<1-\eta_{k} \tau_{k}\|A\|-\tau_{k} L_{k}
$$

2) $\frac{1}{\beta}-\frac{\tau_{k}\|A\|}{\eta_{k}}>\frac{c}{\beta}>0$.

Proof. Using the primal update rule, we have

$$
\begin{equation*}
\left\|x_{k+1}-x\right\|^{2}=\left\|x_{k}-x\right\|^{2}+\left\|x_{k+1}-x_{k}\right\|^{2}-2 \tau_{k}\left\langle\nabla f\left(x_{k}\right)+A^{T} y_{k+1}, x_{k}-x\right\rangle \tag{12}
\end{equation*}
$$

We address each term in the RHS separately. Using the convexity of $f$ we bound the last term of (12):

$$
\begin{equation*}
-2 \tau_{k}\left\langle\nabla f\left(x_{k}\right)+A^{T} y_{k+1}, x_{k}-x\right\rangle \leq 2 \tau_{k}\left(f(x)-f\left(x_{k}\right)\right)+2 \tau_{k}\left\langle A\left(x-x_{k}\right), y_{k+1}\right\rangle \tag{13}
\end{equation*}
$$

For the second term of (12) we use an expansion similar to the analysis in [Malitsky and Mishchenko, 2020] along with the primal update rule:

$$
\begin{align*}
& \left\|x_{k+1}-x_{k}\right\|^{2}=2\left\|x_{k+1}-x_{k}\right\|^{2}-\left\|x_{k+1}-x_{k}\right\|^{2} \\
& =2 \tau_{k}\left\langle\nabla f\left(x_{k}\right)+A^{T} y_{k+1}, x_{k}-x_{k+1}\right\rangle-\left\|x_{k+1}-x_{k}\right\|^{2} \\
& = \\
& =2 \tau_{k}\left\langle\nabla f\left(x_{k}\right)-\nabla f\left(x_{k-1}\right), x_{k}-x_{k+1}\right\rangle+2 \tau_{k}\left\langle A^{T} y_{k+1}-A^{T} y_{k}, x_{k}-x_{k+1}\right\rangle  \tag{14}\\
& \\
& \quad+2 \tau_{k}\left\langle\nabla f\left(x_{k-1}\right)+A^{T} y_{k}, x_{k}-x_{k+1}\right\rangle-\left\|x_{k+1}-x_{k}\right\|^{2}
\end{align*}
$$

Notice that the first term in (14) gives us the opportunity to insert a dependence on the local Lipschitz constant $L_{k}$. Using Cauchy-Schwarz, the definition of $L_{k}$ and Young's inequality, we indeed take this opportunity and get:

$$
\begin{align*}
\left\langle\nabla f\left(x_{k}\right)-\nabla f\left(x_{k-1}\right), x_{k}-x_{k+1}\right\rangle & \leq L_{k}\left\|x_{k}-x_{k-1}\right\|\left\|x_{k+1}-x_{k}\right\| \\
& \leq \frac{L_{k}}{2}\left(\left\|x_{k}-x_{k-1}\right\|^{2}+\left\|x_{k+1}-x_{k}\right\|^{2}\right) \tag{15}
\end{align*}
$$

Similarly, we bound the second term in (14) and obtain:

$$
\begin{equation*}
\left\langle A^{T} y_{k+1}-A^{T} y_{k}, x_{k}-x_{k+1}\right\rangle \leq \frac{\|A\| \eta}{2}\left\|x_{k+1}-x_{k}\right\|^{2}+\frac{\|A\|}{2 \eta}\left\|y_{k+1}-y_{k}\right\|^{2} \tag{16}
\end{equation*}
$$

where $\eta>0$ is a free parameter coming from Young's inequality.
Finally, for the third term in (14) we use the update rule and the convexity of $f$ :

$$
\begin{align*}
\left\langle\nabla f\left(x_{k-1}\right)+A^{T} y_{k}, x_{k}-x_{k+1}\right\rangle & =\left\langle\frac{1}{\tau_{k-1}}\left(x_{k-1}-x_{k}\right), \tau_{k}\left(\nabla f\left(x_{k}\right)+A^{T} y_{k+1}\right)\right\rangle \\
& \leq \theta_{k}\left(f\left(x_{k-1}\right)-f\left(x_{k}\right)\right)+\theta_{k}\left\langle A\left(x_{k-1}-x_{k}\right), y_{k+1}\right\rangle \tag{17}
\end{align*}
$$

Replacing (15), (16) and (17) into (14), we get

$$
\begin{array}{r}
\left\|x_{k+1}-x_{k}\right\|^{2} \leq \tau_{k} L_{k}\left\|x_{k}-x_{k-1}\right\|^{2}+\left(\tau_{k}\|A\| \eta+\tau_{k} L_{k}-1\right)\left\|x_{k+1}-x_{k}\right\|^{2} \\
+\frac{\tau_{k}\|A\|}{\eta}\left\|y_{k+1}-y_{k}\right\|^{2}+2 \tau_{k} \theta_{k}\left(f\left(x_{k-1}\right)-f\left(x_{k}\right)\right) \\
+2 \tau_{k} \theta_{k}\left\langle A\left(x_{k-1}-x_{k}\right), y_{k+1}\right\rangle \tag{18}
\end{array}
$$

Finally, replacing (18) and (13) back into (12) and using the fact that $\theta_{k}\left\langle A\left(x_{k-1}-x_{k}\right), y_{k+1}\right\rangle+$ $\left\langle A\left(x-x_{k}\right), y_{k+1}\right\rangle=-\left\langle A\left(\tilde{x}_{k}-x\right), y_{k+1}\right\rangle$, we obtain the inequality for the primal iterate sequence:

$$
\begin{array}{r}
\left\|x_{k+1}-x\right\|^{2} \leq\left\|x_{k}-x\right\|^{2}+\tau_{k} L_{k}\left\|x_{k}-x_{k-1}\right\|^{2}+\left(\tau_{k}\|A\| \eta+\tau_{k} L_{k}-1\right)\left\|x_{k+1}-x_{k}\right\|^{2} \\
+\frac{\tau_{k}\|A\|}{\eta}\left\|y_{k+1}-y_{k}\right\|^{2}+2 \tau_{k} \theta_{k}\left(f\left(x_{k-1}\right)-f\left(x_{k}\right)\right)+2 \tau_{k}\left(f(x)-f\left(x_{k}\right)\right) \\
-2 \tau_{k}\left\langle A\left(\tilde{x}_{k}-x\right), y_{k+1}\right\rangle \tag{19}
\end{array}
$$

We now seek a similar result for the dual sequence. For this, we use the following characterization of the proximal operator:

$$
\begin{equation*}
u=\operatorname{prox}_{g^{*}}(x) \Longleftrightarrow\langle u-x, z-u\rangle \geq g^{*}(u)-g^{*}(z) \forall z . \tag{20}
\end{equation*}
$$

Thus, letting $u=y_{k+1}, x=y_{k}$ and $z=y$ in (20), we obtain:

$$
g^{*}(y) \geq g^{*}\left(y_{k+1}\right)+\left\langle\frac{1}{\sigma_{k}}\left(y_{k}-y_{k+1}\right), y-y_{k+1}\right\rangle+\left\langle A \tilde{x}_{k}, y-y_{k+1}\right\rangle
$$

Using the cosine rule for the second term, the fact that $\sigma_{k}=\beta \tau_{k}$ and multiplying both sides by $2 \tau_{k}>0$, we obtain:

$$
\begin{array}{r}
\frac{1}{\beta}\left\|y_{k+1}-y\right\|^{2} \leq \frac{1}{\beta}\left\|y_{k}-y\right\|^{2}-\frac{1}{\beta}\left\|y_{k+1}-y_{k}\right\|^{2}+2 \tau_{k}\left(g^{*}(y)-g^{*}\left(y_{k+1}\right)\right) \\
+2 \tau_{k}\left\langle A \tilde{x}_{k}, y_{k+1}-y\right\rangle \tag{21}
\end{array}
$$

Summing (21) with (19) we obtain the following recurrence:

$$
\begin{align*}
& \left\|x_{k+1}-x\right\|^{2}+\frac{1}{\beta}\left\|y_{k+1}-y\right\|^{2} \\
& \qquad \begin{aligned}
\leq & \left\|x_{k}-x\right\|^{2}+\frac{1}{\beta}\left\|y_{k}-y\right\|^{2}+\tau_{k} L_{k}\left\|x_{k}-x_{k-1}\right\|^{2} \\
& +\left(\tau_{k}\|A\| \eta+\tau_{k} L_{k}-1\right)\left\|x_{k+1}-x_{k}\right\|^{2}+\left(\frac{\tau_{k}\|A\|}{\eta}-\frac{1}{\beta}\right)\left\|y_{k+1}-y_{k}\right\|^{2} \\
& +2 \tau_{k}\left(\theta_{k}\left(f\left(x_{k-1}\right)-f\left(x_{k}\right)\right)+f(x)-f\left(x_{k}\right)+g^{*}(y)-g^{*}\left(y_{k+1}\right)\right. \\
& +2 \tau_{k}\left\langle A x, y_{k+1}\right\rangle-2 \tau_{k}\left\langle A \tilde{x}_{k}, y\right\rangle .
\end{aligned}
\end{align*}
$$

We further process the terms involving $f$ and $g^{*}$ on the right-hand side in order to form the $P_{x, y}(\cdot)$ and $D_{x, y}(\cdot)$ :

$$
f(x)-f\left(x_{k}\right)=-P_{x, y}\left(x_{k}\right)+\left\langle A\left(x_{k}-x\right), y\right\rangle
$$

$$
\begin{aligned}
\theta_{k}\left(f\left(x_{k-1}\right)-f\left(x_{k}\right)\right) & =\theta_{k} P_{x, y}\left(x_{k-1}\right)-\theta_{k} P_{x, y}\left(x_{k}\right)+\left\langle\theta_{k} A\left(x_{k}-x_{k-1}\right), y\right\rangle \\
g^{*}(y)-g^{*}\left(y_{k+1}\right) & =-D_{x, y}\left(y_{k+1}\right)-\left\langle A x, y_{k+1}-y\right\rangle
\end{aligned}
$$

Replacing the above expressions into (22) and noting that $\left\langle A\left(\tilde{x}_{k}-x\right), y\right\rangle-\left\langle A x, y_{k+1}-y\right\rangle+$ $\left\langle A x, y_{k+1}\right\rangle-\left\langle A \tilde{x}_{k}, y\right\rangle=0$, we obtain:

$$
\begin{align*}
& \left\|x_{k+1}-x\right\|^{2}+\frac{1}{\beta}\left\|y_{k+1}-y\right\|^{2}+\left(1-\tau_{k}\|A\| \eta-\tau_{k} L_{k}\right)\left\|x_{k+1}-x_{k}\right\|^{2} \\
& +\left(\frac{1}{\beta}-\frac{\tau_{k}\|A\|}{\eta}\right)\left\|y_{k+1}-y_{k}\right\|^{2}+2 \tau_{k}\left(1+\theta_{k}\right) P_{x, y}\left(x_{k}\right)+2 \tau_{k} D_{x, y}\left(y_{k+1}\right) \\
& \quad \leq\left\|x_{k}-x\right\|^{2}+\frac{1}{\beta}\left\|y_{k}-y\right\|^{2}+\tau_{k} L_{k}\left\|x_{k}-x_{k-1}\right\|^{2}+2 \tau_{k} \theta_{k} P_{x, y}\left(x_{k-1}\right) \tag{23}
\end{align*}
$$

What is left to do for obtaining the stated result is to choose $\eta$, possibly depending on $k$, such that the corresponding terms are positive. First, note that $\tau_{k} L_{k}\left\|x_{k}-x_{k-1}\right\|^{2}<\frac{1}{2}\left\|x_{k}-x_{k-1}\right\|^{2}$ because $z \mapsto \frac{z}{2 \sqrt{z^{2}+a}}, a>0$ is an increasing function whose limit at $\infty$ is $\frac{1}{2}$ and we have:

$$
\begin{equation*}
\tau_{k} L_{k} \leq \frac{L_{k}}{2 \sqrt{L_{k}^{2}+(\beta /(1-c))\|A\|^{2}}}<\frac{1}{2} \tag{24}
\end{equation*}
$$

Next we need to choose $\eta=\eta_{k}$ (iteration-dependent) to satisfy:

$$
\left\{\begin{array}{l}
\frac{1}{\beta}-\frac{\tau_{k}\|A\|}{\eta_{k}}>0 \\
1-\tau_{k}\|A\| \eta_{k}-\tau_{k} L_{k}>\frac{1}{2}
\end{array}\right.
$$

However, for theoretical purposes related to controlling the sequence $\left\|y_{k+1}-y_{k}\right\|^{2}$, we strengthen the first inequality to $\frac{1}{\beta}-\frac{\tau_{k}\|A\|}{\eta_{k}}>\frac{c}{\beta}, c \in(0,1)$. In practice, this constant is chosen as small as possible. The new conditions to be satisfied are:

$$
\left\{\begin{array} { l } 
{ \frac { 1 } { \beta } - \frac { \tau _ { k } \| A \| } { \eta _ { k } } > \frac { c } { \beta } , }  \tag{25}\\
{ 1 - \tau _ { k } \| A \| \eta _ { k } - \tau _ { k } L _ { k } > \frac { 1 } { 2 } , }
\end{array} \Longleftrightarrow \left\{\begin{array}{l}
\eta_{k}>\frac{\beta \tau_{k}\|A\|}{1-c} \\
\eta_{k}<\frac{1-2 \tau_{k} L_{k}}{2 \tau_{k}\|A\|}
\end{array}\right.\right.
$$

The question we need to answer therefore is: given the expression of $\tau_{k}$, is the interval always valid for choosing $\eta_{k} \in\left(\frac{\beta \tau_{k}\|A\|}{1-c}, \frac{1-2 \tau_{k} L_{k}}{2 \tau_{k}\|A\|}\right)$ ?
To answer, we form the corresponding quadratic inequality in $\tau_{k}$ :

$$
\begin{equation*}
\frac{\beta \tau_{k}\|A\|}{1-c}-\frac{1-2 \tau_{k} L_{k}}{2 \tau_{k}\|A\|}<0 \Longleftrightarrow \frac{2 \beta \tau_{k}^{2}\|A\|^{2}}{1-c}+2 \tau_{k} L_{k}-1<0 \tag{26}
\end{equation*}
$$

whose 2 real roots are given by:

$$
\begin{cases}\tau_{k, 1}=\frac{1}{L_{k}-\sqrt{L_{k}^{2}+2(\beta /(1-c))\|A\|^{2}}} & <0 \\ \tau_{k, 2}=\frac{1}{L_{k}+\sqrt{L_{k}^{2}+2(\beta /(1-c))\|A\|^{2}}}>0\end{cases}
$$

For inequality (26) to be satisfied, we need:

$$
\begin{equation*}
\tau_{k} \in\left(0, \tau_{k, 2}\right)=\left(0, \frac{1}{L_{k}+\sqrt{L_{k}^{2}+2(\beta /(1-c))\|A\|^{2}}}\right), \forall k \tag{27}
\end{equation*}
$$

The lower bound for $\tau_{k}$ trivially holds, and for the upper bound we make the following observation:

$$
\begin{aligned}
L_{k}+\sqrt{L_{k}^{2}+2(\beta /(1-c))\|A\|^{2}} & =\frac{2\left[\sqrt{L_{k}^{2}}+\sqrt{L_{k}^{2}+2(\beta /(1-c))\|A\|^{2}}\right]}{2} \\
& \stackrel{\text { Jensen }}{<} 2 \sqrt{\frac{2 L_{k}^{2}+2(\beta /(1-c))\|A\|^{2}}{2}} \\
& =2 \sqrt{L_{k}^{2}+(\beta /(1-c))\|A\|^{2}}
\end{aligned}
$$

Here Jensen's inequality holds strictly because function $\sqrt{ }$. is strictly concave and $L_{k}^{2} \neq L_{k}^{2}+$ $2\|A\|^{2} \beta$. Thus, we obtain:

$$
0<\tau_{k} \leq \frac{1}{2 \sqrt{L_{k}^{2}+\|A\|^{2} \beta}}<\frac{1}{L_{k}+\sqrt{L_{k}^{2}+2 \beta\|A\|^{2}}}=\tau_{k, 2} \quad \forall k
$$

It follows that we can find an $\eta_{k} \in\left(\frac{\beta \tau_{k}\|A\|}{1-c}, \frac{1-2 \tau_{k} L_{k}}{2 \tau_{k}\|A\|}\right), \forall k$, which implies that conditions (25) can always be satisfied. This concludes the proof.

## B. 2 Proof of Theorem 4.1

Theorem 4.1. Consider APDA along with Assumptions 3.1 and 3.2, and let $\left(x^{*}, y^{*}\right) \in \mathcal{X} \times \mathcal{Y}$ be a saddle point of problem (2). Then, for all $k$

1) Boundedness. The sequence $\left\{\left(x_{k}, y_{k}\right)\right\}$ is bounded. Specifically, for all $k$,

$$
\left\|x_{k}-x^{*}\right\|^{2}+\left\|y_{k}-y^{*}\right\|^{2} \leq M
$$

where $M:=\left\|x_{1}-x^{*}\right\|^{2}+\frac{1}{\beta}\left\|y_{1}-y^{*}\right\|^{2}+\frac{1}{2}\left\|x_{1}-x_{0}\right\|^{2}<\infty$.
2) Convergence to a saddle point. The sequence $\left\{\left(x_{k}, y_{k}\right)\right\}$ converges to a saddle point of (2).
3) Ergodic convergence. Let $S_{k}:=\sum_{i=1}^{k} \tau_{i}, \quad X_{k} \quad:=\frac{1}{S_{k}}\left(\tau_{k}\left(1+\theta_{k}\right) x_{k}+\right.$ $\left.\sum_{i=1}^{k-1}\left(\tau_{i}\left(1+\theta_{i}\right)-\tau_{i+1} \theta_{i+1}\right) x_{i}\right)$ and $Y_{k}:=\frac{1}{S_{k}} \sum_{i=1}^{k} \tau_{i} y_{i+1} . \quad$ Then, for any bounded $B_{1} \times B_{2} \in \mathcal{X} \times \mathcal{Y}$ and for all $k$,

$$
\mathcal{G}_{B_{1} \times B_{2}}\left(X_{k}, Y_{k}\right) \leq \frac{M\left(B_{1}, B_{2}\right) \sqrt{L^{2}+(\beta /(1-c))\|A\|^{2}}}{k}
$$

where $L$ is the Lipschitz constant of $\nabla f$ over the compact set $\overline{\operatorname{Conv}}\left(\left\{x^{*}, x_{0}, x_{1}, \ldots\right\}\right)$ and $M\left(B_{1}, B_{2}\right)=\sup _{(x, y) \in B_{1} \times B_{2}}\left\|x_{1}-x\right\|^{2}+\frac{1}{\beta}\left\|y_{1}-x\right\|^{2}+\frac{1}{2}\left\|x_{1}-x_{0}\right\|^{2}$.

Proof. 1) Sequence boundedness. Using the inequality of Lemma (4.1) with $(x, y)=\left(x^{*}, y^{*}\right)$ and the fact that $\tau_{k} L_{k}<\frac{1}{2}, \forall k$, unrolling it over the iterations and rearranging the terms we obtain:

$$
\begin{aligned}
& \left\|x_{k+1}-x^{*}\right\|^{2}+\frac{1}{\beta}\left\|y_{k+1}-y^{*}\right\|^{2}+\left(1-\eta_{k} \tau_{k}\|A\|-\tau_{k} L_{k}\right)\left\|x_{k+1}-x_{k}\right\|^{2} \\
& +\sum_{i=1}^{k-1}\left(\frac{1}{2}-\eta_{i} \tau_{i}\|A\|-\tau_{i} L_{i}\right)\left\|x_{i+1}-x_{i}\right\|^{2}+\frac{c}{\beta} \sum_{i=1}^{k}\left\|y_{i+1}-y_{i}\right\|^{2}+2 \tau_{k}\left(1+\theta_{k}\right) P_{x^{*}, y^{*}}\left(x_{k}\right) \\
& \quad+2 \sum_{i=2}^{k-1}\left(\tau_{i}\left(1+\theta_{i}\right)-\tau_{i+1} \theta_{i+1}\right) P_{x^{*}, y^{*}}\left(x_{i}\right)+2 \sum_{i=1}^{k} \tau_{i} D_{x^{*}, y^{*}}\left(y_{i+1}\right)
\end{aligned}
$$

$$
\begin{equation*}
\leq\left\|x_{1}-x^{*}\right\|^{2}+\frac{1}{\beta}\left\|y_{1}-y^{*}\right\|^{2}+\frac{1}{2}\left\|x_{1}-x_{0}\right\|^{2}+2 \tau_{1} \theta_{1} P_{x^{*}, y^{*}}\left(x_{0}\right) \tag{28}
\end{equation*}
$$

All the terms on the left hand-side of (28) are non-negative:

$$
\left\{\begin{array}{l}
\frac{1}{\beta}-\frac{\tau_{k}\|A\|}{\eta_{i}}>\frac{c}{\beta}>0, \forall i  \tag{byLemma4.1}\\
\frac{1}{2}-\eta_{i} \tau_{i}\|A\|-\tau_{i} L_{i}>0, \forall i \\
\tau_{i+1} \theta_{i+1} \leq \tau_{i} \sqrt{1+\theta_{i}} \theta_{i+1} \leq \tau_{i}\left(1+\theta_{i}\right) \\
P_{x^{*}, y^{*}}(x) \geq 0, D_{x^{*}, y^{*}}(y) \geq 0 \quad \forall x, y
\end{array}\right.
$$

(by stepsize update rule) (by the saddle point property)

Also, by our parameter setup we have that $\theta_{1}=0$. Consequently, it holds that:

$$
\left\|x_{k+1}-x^{*}\right\|^{2}+\frac{1}{\beta}\left\|y_{k+1}-y^{*}\right\|^{2} \leq M<\infty \quad \forall k
$$

where $M:=\left\|x_{1}-x^{*}\right\|^{2}+\frac{1}{\beta}\left\|y_{1}-y^{*}\right\|^{2}+\frac{1}{2}\left\|x_{1}-x_{0}\right\|^{2}$, which implies that the sequence is bounded.

We make the following remarks which will be useful for the remainder of the theorem's proof:

- Boundedness of $\left\{x_{k}\right\}$ together with the local Lipschitz continuity of $\nabla f$ from Assumption 3.1 implies that there exists $L>0$ such that $f$ is $L$-smooth over $\overline{\operatorname{Conv}}\left(\left\{x^{*}, x_{0}, x_{1}, \ldots\right\}\right)$. Furthermore, $L \geq L_{k} \forall k$.
- A consequence of the prior point is that $\tau_{k}$ has a uniform and positive lower-bound. By the definition of APDA it holds that:

$$
\tau_{1}=\frac{1}{2 \sqrt{L_{1}^{2}+(\beta /(1-c))\|A\|^{2}}} \geq \frac{1}{2 \sqrt{L^{2}+(\beta /(1-c))\|A\|^{2}}}
$$

and, from the definition of $\tau_{k}$ it is straightforward to see that at every iteration we either explicitly increase $\tau_{k}$ relative to $\tau_{k-1}$ or otherwise set it to an expression dictated by the local smoothness constant $L_{k}$. Thus it holds that:

$$
\begin{equation*}
\tau_{k} \geq \frac{1}{2 \sqrt{L^{2}+(\beta /(1-c))\|A\|^{2}}}, \forall k \tag{29}
\end{equation*}
$$

- Furthermore, the existence of $L$ guarantees that $\tau_{k} L_{k}$ can have a tighter upper bound than the $1 / 2$ shown before, as follows:

$$
\begin{align*}
\tau_{k} L_{k} & \leq \frac{L_{k}}{2 \sqrt{L_{k}^{2}+(\beta /(1-c))\|A\|^{2}}} \\
& \leq \frac{L}{2 \sqrt{L^{2}+(\beta /(1-c))\|A\|^{2}}}, \tag{30}
\end{align*}
$$

where we used the fact that $z \mapsto \frac{z}{2 \sqrt{z^{2}+a}}, a>0$ is an increasing function.

- Finally, due to the point above, we can uniformly lower bound the coefficients of terms $\left\|x_{k+1}-x_{k}\right\|^{2}$ on the LHS of (28), and thus obtain:

$$
\frac{1}{2}\left(1-\frac{L}{\sqrt{L^{2}+(\beta /(1-c))\|A\|^{2}}}\right) \sum_{i=1}^{k-1}\left\|x_{i}-x_{i+1}\right\|^{2}+\frac{c}{\beta} \sum_{i=1}^{k}\left\|y_{i+1}-y_{i}\right\|^{2} \leq M
$$

which conveniently ensures that:

$$
\left\{\begin{array}{l}
\lim _{k \rightarrow \infty}\left\|x_{k}-x_{k-1}\right\|^{2}=0  \tag{31}\\
\lim _{k \rightarrow \infty}\left\|y_{k}-y_{k-1}\right\|^{2}=0
\end{array}\right.
$$

2) Convergence to a saddle point. Let $(\hat{x}, \hat{y})$ be an arbitrary cluster point of the sequence $\left\{\left(x_{k}, y_{k}\right)\right\}$. Since we have shown that the sequence is bounded, then there must exist a subsequence $\left\{\left(x_{k_{i}}, y_{k_{i}}\right)\right\}$, such that $\lim _{i \rightarrow \infty}\left(x_{k_{i}}, y_{k_{i}}\right)=(\hat{x}, \hat{y})$. We wish to prove that $(\hat{x}, \hat{y})$ is a saddle point of (2).
More precisely, we wish to prove that $P_{\hat{x}, \hat{y}}(x) \geq 0$ and $D_{\hat{x}, \hat{y}}(y) \geq 0$ for $\forall x, y$, respectively. For convenience, we remind the reader the definitions of these two quantities:

$$
\begin{aligned}
& P_{\hat{x}, \hat{y}}(x)=f(x)-f(\hat{x})+\langle A(x-\hat{x}), \hat{y}\rangle \\
& D_{\hat{x}, \hat{y}}(y)=g^{*}(y)-g^{*}(\hat{y})-\langle A \hat{x}, y-\hat{y}\rangle
\end{aligned}
$$

We start with $P_{\hat{x}, \hat{y}}(x)$ :

$$
\begin{array}{rlr}
P_{\hat{x}, \hat{y}}(x) & =f(x)-f(\hat{x})+\langle A(x-\hat{x}), \hat{y}\rangle \\
& =\lim _{i \rightarrow \infty} f(x)-f\left(x_{k_{i}}\right)+\left\langle A\left(x-x_{k_{i}}\right), y_{k_{i}}\right\rangle \\
& \geq \lim _{i \rightarrow \infty}\left\langle\nabla f\left(x_{k_{i}}\right)+A^{*} y_{k_{i}+1}, x-x_{k_{i}}\right\rangle+\left\langle A^{*}\left(y_{k_{i}}-y_{k_{i}+1}\right), x-x_{k_{i}}\right\rangle \quad \text { (Convexity of } f \text { ) } \\
& =\lim _{i \rightarrow \infty}\left\langle\frac{x_{k_{i}+1}-x_{k_{i}}}{\tau_{k_{i}}}, x-x_{k_{i}}\right\rangle+\left\langle A^{*}\left(y_{k_{i}}-y_{k_{i}+1}\right), x-x_{k_{i}}\right\rangle \\
& =0 . & \text { (Primal update rule) }  \tag{29}\\
& \text { (By (29), (31)) }
\end{array}
$$

Showing the analogous result for $D_{\hat{x}, \hat{y}}(y)$ relies on similar arguments, with the additional requirement that $\theta_{k}$ is uniformly upper bounded. From the update rule of $\tau_{k}$ we have:

$$
\begin{equation*}
\theta_{k}=\frac{\tau_{k}}{\tau_{k-1}} \leq \sqrt{1+\theta_{k-1}} \Longrightarrow \theta_{k} \leq \sqrt{1+\ldots+\sqrt{1+\theta_{2}}} \leq \underbrace{\sqrt{1+\ldots+\sqrt{1+1}}}_{k-2 \text { times }} \leq 2 \tag{32}
\end{equation*}
$$

where the second to last inequality comes from the way APDA's first two iterations are set up.
Therefore, we have that $\forall y \in \mathcal{Y}$ :

$$
\left.\begin{array}{rl}
D_{\hat{x}, \hat{y}}(y) & =g^{*}(y)-g^{*}(\hat{y})-\langle A \hat{x}, y-\hat{y}\rangle \\
& \geq g^{*}(y)-\liminf _{i \rightarrow \infty} g^{*}\left(y_{k_{i}}\right)-\left\langle A \liminf _{i \rightarrow \infty} x_{k_{i}}, y-\liminf _{i \rightarrow \infty} y_{k_{i}}\right\rangle \\
& =\limsup _{i \rightarrow \infty} g^{*}(y)-g^{*}\left(y_{k_{i}}\right)-\left\langle A x_{k_{i}}, y-y_{k_{i}}\right\rangle \\
& \geq \limsup _{i \rightarrow \infty}\left\langle\frac{y_{k_{i}-1}-y_{k_{i}}}{\sigma_{k_{i}-1}}, y-y_{k_{i}}\right\rangle+\left\langle A\left(\tilde{x}_{k_{i}-1}-x_{k_{i}}\right), y-y_{k_{i}}\right\rangle \\
& =\limsup _{i \rightarrow \infty}\left\langle\frac{y_{k_{i}-1}-y_{k_{i}}}{\beta \tau_{k_{i}-1}}, y-y_{k_{i}}\right\rangle+\left\langle A\left[x_{k_{i}-1}-x_{k_{i}}+\theta_{k_{i}-1}\left(x_{k_{i}-1}-x_{k_{i}-2}\right)\right], y-y_{k_{i}}\right\rangle \\
& =0 \tag{29}
\end{array} \quad \text { (By (29), (31), (32)) }\right)
$$

3) Gap rate. Unrolling the inequality of Lemma 4.1 for some $(x, y) \in B_{1} \times B_{2}$, we obtain:

$$
\begin{align*}
& \left\|x_{k+1}-x\right\|^{2}+\frac{1}{\beta}\left\|y_{k+1}-x\right\|^{2}+\left(1-\eta_{k} \tau_{k}\|A\|-\tau_{k} L_{k}\right)\left\|x_{k+1}-x_{k}\right\|^{2} \\
& +\sum_{i=1}^{k-1}\left(\frac{1}{2}-\eta_{i} \tau_{i}\|A\|-\tau_{i} L_{i}\right)\left\|x_{i+1}-x_{i}\right\|^{2}+\frac{c}{\beta} \sum_{i=1}^{k}\left\|y_{i+1}-y_{i}\right\|^{2}+2 \tau_{k}\left(1+\theta_{k}\right) P_{x, y}\left(x_{k}\right) \\
& \quad+2 \sum_{i=2}^{k-1}\left(\tau_{i}\left(1+\theta_{i}\right)-\tau_{i+1} \theta_{i+1}\right) P_{x, y}\left(x_{i}\right)+2 \sum_{i=1}^{k} \tau_{i} D_{x, y}\left(y_{i+1}\right) \\
& \leq\left\|x_{1}-x\right\|^{2}+\frac{1}{\beta}\left\|y_{1}-y\right\|^{2}+\frac{1}{2}\left\|x_{1}-x_{0}\right\|^{2} \tag{33}
\end{align*}
$$

First, note that due to $\theta_{1}=0$, the following holds:

$$
\tau_{k}\left(1+\theta_{k}\right)+\sum_{i=1}^{k-1}\left(\tau_{i}\left(1+\theta_{i}\right)-\tau_{i+1} \theta_{i+1}\right)=\sum_{i=1}^{k} \tau_{i}=: S_{k}
$$

Second, since all the terms on the LHS of (33) except those involving $P_{x, y}(\cdot)$ and $D_{x, y}(\cdot)$ are non-negative and, for fixed $(x, y) \in \mathcal{X} \times \mathcal{Y}$ the functions $P_{x, y}(\cdot)$ and $D_{x, y}(\cdot)$ are convex, we have:

$$
\begin{align*}
& 2 S_{k}\left(P_{x, y}\left(X_{k}\right)+D_{x, y}\left(Y_{k}\right)\right) \\
& \leq 2 \tau_{k}\left(1+\theta_{k}\right) P_{x, y}\left(x_{k}\right)+2 \sum_{i=2}^{k-1}\left(\tau_{i}\left(1+\theta_{i}\right)-\tau_{i+1} \theta_{i+1}\right) P_{x, y}\left(x_{i}\right)+2 \sum_{i=1}^{k} \tau_{i} D_{x, y}\left(y_{i+1}\right) \\
& \leq\left\|x_{1}-x\right\|^{2}+\frac{1}{\beta}\left\|y_{1}-y\right\|^{2}+\frac{1}{2}\left\|x_{1}-x_{0}\right\|^{2} \tag{34}
\end{align*}
$$

Lastly, since $\tau_{k} \geq \frac{1}{2 \sqrt{L^{2}+(\beta /(1-c))\|A\|^{2}}}, \forall k$, we have that $S_{k} \geq \frac{k}{2 \sqrt{L^{2}+(\beta /(1-c))\|A\|^{2}}}$ and the rate for the restricted gap is:

$$
\begin{aligned}
& \mathcal{G}_{B_{1} \times B_{2}}\left(X_{k}, Y_{k}\right) \\
&=\sup _{(x, y) \in B_{1} \times B_{2}} P_{x, y}\left(X_{k}\right)+D_{x, y}\left(Y_{k}\right) \\
& \leq \sup _{(x, y) \in B_{1} \times B_{2}} \frac{\left(\left\|x_{1}-x\right\|^{2}+\frac{1}{\beta}\left\|y_{1}-y\right\|^{2}+\frac{1}{2}\left\|x_{1}-x_{0}\right\|^{2}\right) \sqrt{L^{2}+(\beta /(1-c))\|A\|^{2}}}{k} \\
& \quad=\frac{M\left(B_{1}, B_{2}\right) \sqrt{L^{2}+(\beta /(1-c))\|A\|^{2}}}{k},
\end{aligned}
$$

which concludes the proof of the theorem.

## B. 3 Proof of Theorem 4.2

Before proving the result of Theorem 4.2, a few remarks are in order. First, the boundedness result of Theorem 4.1 point 1 ) also holds for constant $c=0$, since this constant was required only for proving convergence to a saddle point in point 2) of the theorem. Second, taking a stepsize smaller than the originally considered $\tau_{k}$ will not change the validity of Lemma 4.1 or the boundedness result of Theorem 4.1, as it remains within the interval given in (27).

Consequently, for studying APDA under the additional Assumption 3.3 we can simplify the stepsize expression by taking $c=0$, since now we will prove convergence of the iterates directly by using the strong convexity and full row-rank assumptions. Specifically, we consider $\tau_{k}$ as defined in (8), which is smaller than the one originally considered and, due to the above remarks it ensures that APDA produces a bounded sequence. It follows that, under the local smoothness and local strong convexity assumptions, there exist constant $L$ and $\mu$ such that $f$ is $L$-smooth and $\mu$-strongly convex over $\overline{\operatorname{Conv}}\left(\left\{x^{*}, x_{0}, x_{1}, \ldots\right\}\right)$. This observation suffices to show linear convergence in Theorem 4.2.

Theorem 4.2. Consider APDA along with Assumptions 3.1, 3.2 and 3.3. Let $\left(x^{*}, y^{*}\right) \in \mathcal{X} \times \mathcal{Y}$ be a saddle point of problem (2). Furthermore, let $\tau_{k}$ be defined by (8) and let $s:=\sqrt{4 L^{2}+\beta\|A\|^{2}}$ and $t:=\sqrt{4 \mu^{2}+\beta\|A\|^{2}}$, where $\mu, L$ are the strong convexity and smoothness constants of $f$ over the compact set $\overline{\operatorname{Conv}}\left(\left\{x^{*}, x_{0}, x_{1}, \ldots\right\}\right)$.

Then, for all $k$ :

$$
\left\|x_{k}-x^{*}\right\|^{2}+\frac{1}{\beta}\left\|y_{k}-y^{*}\right\|^{2} \leq(1-\min \{p, q, r\})^{k} M
$$

where the rate constants are given by:

$$
p=\frac{1}{2}, \quad q=\frac{\mu}{4 s}, \quad r=\frac{\beta \sigma_{\min }^{2}(A) \mu}{\beta \sigma_{\min }^{2}(A) \mu+8 s^{2} t+4 L^{2} s}
$$

and $M=\left\|x_{2}-x^{*}\right\|^{2}+\left(\frac{1}{\beta}+T\right)\left\|y_{2}-y^{*}\right\|^{2}+\frac{1}{2}\left\|x_{2}-x_{1}\right\|^{2}+2 \tau_{1} P_{x^{*}, y^{*}}\left(x_{1}\right), T=$ $\frac{\sigma_{\min }^{2}(A) \mu}{8 s^{2} t+4 L^{2} s}$, with $\sigma_{\min }(A)$ representing the smallest singular value of $A$.

Proof. The outline of the proof is first arriving at a strengthened version of the inequality in Lemma 4.1, and then showing that the inequality expresses a contraction.
Since this new stepsize still ensures the boundedness result of Theorem 4.1, there exist $\mu$ and $L$ such that $f$ is $\mu$-strongly convex and $L$-Lipschitz smooth over the compact set $\overline{\operatorname{Conv}}\left(\left\{x^{*}, x_{0}, x_{1}, \ldots\right\}\right)$. From these properties it follows that, for all $k$ :

$$
\begin{array}{ll}
2 \tau_{k}\left\langle\nabla f\left(x_{k}\right), x^{*}-x_{k}\right\rangle & \leq 2 \tau_{k}\left(f\left(x^{*}\right)-f\left(x_{k}\right)\right)-\mu \tau_{k}\left\|x_{k}-x^{*}\right\|^{2} \\
2 \tau_{k}\left\langle\nabla f\left(x_{k}\right), x^{*}-x_{k}\right\rangle & \leq 2 \tau_{k}\left(f\left(x^{*}\right)-f\left(x_{k}\right)\right)-\frac{\tau_{k}}{L}\left\|\nabla f\left(x_{k}\right)-\nabla f\left(x^{*}\right)\right\|^{2}
\end{array}
$$

Summing these two inequalities and dividing by 2 , we obtain a stronger version of equation (13):

$$
\begin{align*}
-2 \tau_{k}\left\langle\nabla f\left(x_{k}\right)+A^{T} y_{k+1}, x_{k}-x^{*}\right\rangle & \leq 2 \tau_{k}\left(f\left(x^{*}\right)-f\left(x_{k}\right)\right)-\frac{\tau_{k} \mu}{2}\left\|x_{k}-x^{*}\right\|^{2} \\
& -\frac{\tau_{k}}{2 L}\left\|\nabla f\left(x_{k}\right)-\nabla f\left(x^{*}\right)\right\|^{2}+2 \tau_{k}\left\langle A\left(x^{*}-x_{k}\right), y_{k+1}\right\rangle \tag{35}
\end{align*}
$$

We further bound the term $\left\|\nabla f\left(x_{k}\right)-\nabla f\left(x^{*}\right)\right\|^{2}$ in (40):

$$
\begin{align*}
\left\|\nabla f\left(x^{*}\right)-\nabla f\left(x_{k}\right)\right\|^{2}= & \left\|A^{T}\left(y_{k+1}-y^{*}\right)-\frac{x_{k}-x_{k+1}}{\tau_{k}}\right\|^{2}  \tag{36}\\
\geq & \left\|A^{T}\left(y_{k+1}-y^{*}\right)\right\|^{2}+\frac{1}{\tau_{k}^{2}}\left\|x_{k+1}-x_{k}\right\|^{2} \\
& -\frac{2}{\tau_{k}}\left\|A^{T}\left(y_{k+1}-y^{*}\right)\right\|\left\|x_{k+1}-x_{k}\right\|  \tag{37}\\
\geq & \left\|A^{T}\left(y_{k+1}-y^{*}\right)\right\|^{2}+\frac{1}{\tau_{k}^{2}}\left\|x_{k+1}-x_{k}\right\|^{2} \\
& \quad-\left(\frac{1}{\xi+1}\left\|A^{T}\left(y_{k+1}-y^{*}\right)\right\|^{2}+\frac{\xi+1}{\tau_{k}^{2}}\left\|x_{k+1}-x_{k}\right\|^{2}\right)  \tag{38}\\
\geq & \frac{\xi \sigma_{\min }^{2}(A)}{\xi+1}\left\|y_{k+1}-y^{*}\right\|^{2}-\frac{\xi}{\tau_{k}^{2}}\left\|x_{k+1}-x_{k}\right\|^{2}, \tag{39}
\end{align*}
$$

where line (36) comes from the primal iterate update rule and the optimality condition (4); line (37) comes from developing the square and applying Cauchy-Schwarz; line (38) comes from applying Young's inequality with constant $1+\xi$, where $\xi>0$; line (39) comes from the assumption of $A$ having full-row rank, which implies that $\left\|A^{T}\left(y_{k+1}-y^{*}\right)\right\|^{2} \geq \sigma_{\min }^{2}(A)\left\|y_{k+1}-y^{*}\right\|^{2}$.

Finally, setting $\xi=2 \tau_{k}^{2} L_{k} L$ we obtain that:

$$
\begin{align*}
-2 \tau_{k}\left\langle\nabla f\left(x_{k}\right)+A^{T} y_{k+1}, x_{k}-x^{*}\right\rangle & \leq 2 \tau_{k}\left(f\left(x^{*}\right)-f\left(x_{k}\right)\right)-\frac{\tau_{k} \mu}{2}\left\|x_{k}-x^{*}\right\|^{2} \\
& -\frac{\tau_{k}^{3} L_{k} \sigma_{\min }^{2}(A)}{1+2 \tau_{k}^{2} L_{k} L}\left\|y_{k+1}-y^{*}\right\|^{2}+\tau_{k} L_{k}\left\|x_{k+1}-x_{k}\right\|^{2} \\
& +2 \tau_{k}\left\langle A\left(x^{*}-x_{k}\right), y_{k+1}\right\rangle \tag{40}
\end{align*}
$$

Replacing inequality (13) with inequality (40) in the proof of Lemma 4.1 and keeping everything else identical, we obtain the a strengthened version of Lemma's 4.1 result:

$$
\begin{aligned}
& \left\|x_{k+1}-x^{*}\right\|^{2}+\left(\frac{1}{\beta}+\frac{\tau_{k}^{3} L_{k} \sigma_{\min }^{2}(A)}{1+2 \tau_{k}^{2} L_{k} L}\right)\left\|y_{k+1}-y^{*}\right\|^{2}+\left(1-\eta_{k} \tau_{k}\|A\|-2 \tau_{k} L_{k}\right) \\
& \left\|x_{k+1}-x_{k}\right\|^{2}+\frac{\eta_{k}-\tau_{k} \beta\|A\|}{\beta \eta_{k}}\left\|y_{k+1}-y_{k}\right\|^{2}+2 \tau_{k}\left(1+\theta_{k}\right) P_{x^{*}, y^{*}}\left(x_{k}\right)+2 \tau_{k} D_{x^{*}, y^{*}}\left(y_{k+1}\right)
\end{aligned}
$$

$$
\begin{equation*}
\leq\left(1-\frac{\mu \tau_{k}}{2}\right)\left\|x_{k}-x^{*}\right\|^{2}+\frac{1}{\beta}\left\|y_{k}-y^{*}\right\|^{2}+\tau_{k} L_{k}\left\|x_{k}-x_{k-1}\right\|^{2}+2 \tau_{k} \theta_{k} P_{x^{*}, y^{*}}\left(x_{k-1}\right) \tag{41}
\end{equation*}
$$

In order to show that this is in fact a contraction, we note a few properties of the terms in (41):
a) It holds that $1-\eta_{k} \tau_{k}\|A\|-2 \tau_{k} L_{k}>1 / 2$ and $\frac{\eta_{k}-\tau_{k} \beta\|A\|}{\beta \eta_{k}}>0$ since:

$$
\begin{aligned}
\begin{cases}1-\eta_{k} \tau_{k}\|A\|-2 \tau_{k} L_{k}>1 / 2 \\
\frac{\eta_{k}-\tau_{k} \beta\|A\|}{\beta \eta_{k}}>0\end{cases} & \Longleftrightarrow\left\{\begin{array}{l}
\eta_{k}<\frac{1}{2 \tau_{k}\|A\|}-\frac{2 L_{k}}{\|A\|} \\
\eta_{k}>\tau_{k} \beta\|A\|
\end{array}\right. \\
& \Longleftrightarrow 2 \beta\|A\|^{2} \tau_{k}^{2}+4 L_{k} \tau_{k}-1<0
\end{aligned}
$$

which holds for any $\tau_{k} \in\left(0, \frac{1}{2 L_{k}+\sqrt{4 L_{k}^{2}+2 \beta\|A\|^{2}}}\right)$. Our choice of $\tau_{k}$ belongs to this interval, and therefore ensures the stated properties;
b) It holds that $\tau_{k} L_{k}<1 / 4$, by the same observation as that in (24) but with a different limit constant given by the new stepsize;
c) It holds that $\tau_{k} \theta_{k} \leq \tau_{k-1} \sqrt{1+\theta_{k-1} / 2} \theta_{k} \leq \tau_{k-1}\left(1+\theta_{k-1} / 2\right)$, by the definitions of $\tau_{k}$ and $\theta_{k}$;
d) It holds that:

$$
\begin{equation*}
\frac{1}{2 \sqrt{4 L^{2}+\beta\|A\|^{2}}} \leq \tau_{k} \leq \frac{1}{2 \sqrt{4 \mu^{2}+\beta\|A\|^{2}}} \tag{42}
\end{equation*}
$$

by the existence of $\mu$ and $L$ over $\overline{\operatorname{Conv}}\left(\left\{x^{*}, x_{0}, x_{1}, \ldots\right\}\right)$ and a similar argument to that in (29), plus the fact that under strong convexity $\|\nabla f(x)-\nabla f(y)\| \geq \mu\|x-y\|$;
e) It holds that:

$$
\begin{equation*}
\frac{\mu}{2 \sqrt{4 \mu^{2}+\beta\|A\|^{2}}} \leq \tau_{k} L_{k} \leq \frac{L}{2 \sqrt{4 L^{2}+\beta\|A\|^{2}}} \tag{43}
\end{equation*}
$$

by a similar argument to that in (30).
Using properties a), b), c) in the list above and ignoring the positive terms on the LHS that do not have a correspondent on the RHS of (41), the main inequality becomes:

$$
\begin{array}{r}
\left\|x_{k+1}-x^{*}\right\|^{2}+\left(\frac{1}{\beta}+T\right)\left\|y_{k+1}-y^{*}\right\|^{2}+\frac{1}{2}\left\|x_{k+1}-x_{k}\right\|^{2}+2 \tau_{k}\left(1+\theta_{k}\right) P_{x^{*}, y^{*}}\left(x_{k}\right) \\
\leq\left(1-\frac{\mu \tau_{k}}{2}\right)\left\|x_{k}-x^{*}\right\|^{2}+\frac{1}{\beta}\left\|y_{k}-y^{*}\right\|^{2}+\frac{1}{2}\left(1-\frac{1}{2}\right)\left\|x_{k}-x_{k-1}\right\|^{2} \\
+2 \tau_{k-1}\left(1+\theta_{k-1} / 2\right) P_{x^{*}, y^{*}}\left(x_{k-1}\right) \tag{44}
\end{array}
$$

where $T$ is given by:

$$
\begin{aligned}
T & :=\frac{\sigma_{\min }^{2}(A) \mu}{8 s^{2} t+4 L^{2} s} \\
& =\frac{\sigma_{\min }^{2}(A) \mu}{8\left(4 L^{2}+\beta\|A\|^{2}\right) \sqrt{4 \mu^{2}+\beta\|A\|^{2}}+4 L^{2} \sqrt{4 L^{2}+\beta\|A\|^{2}}} \\
& \leq \frac{\tau_{k}^{3} L_{k} \sigma_{\min }^{2}(A)}{1+2 \tau_{k}^{2} L_{k} L}
\end{aligned}
$$

where we used the definitions of $s=\sqrt{4 L^{2}+\beta\|A\|^{2}}$ and $t=\sqrt{4 \mu^{2}+\beta\|A\|^{2}}$ to simplify notations.

We thus have the following contractions in (44):

- For $\frac{1}{2}\left\|x_{k+1}-x_{k}\right\|^{2}$ it is: $1-\underbrace{\frac{1}{2}}_{=: p}$;
- For $\left(\frac{1}{\beta}+T\right)\left\|y_{k+1}-y^{*}\right\|^{2}$ it is:

$$
\begin{aligned}
1+\frac{1}{1+T \beta} & =1-\frac{T \beta}{1+T \beta} \\
& =1-\underbrace{\frac{\beta \sigma_{\min }^{2}(A) \mu}{\sigma_{\min }^{2}(A) \mu+8 s^{2} t+4 L^{2} s}}_{=: r}
\end{aligned}
$$

- For $\left\|x_{k+1}-x^{*}\right\|^{2}$ it is:

$$
1-\frac{\mu \tau_{k}}{2} \leq 1-\underbrace{\frac{\mu}{4 s}}_{=: q}
$$

- For $2 \tau_{k}\left(1+\theta_{k}\right) P_{x^{*}, y^{*}}\left(x_{k}\right)$ it is:

$$
\frac{1+\theta_{k-1} / 2}{1+\theta_{k-1}}=1-\frac{\theta_{k-1}}{2\left(1+\theta_{k-1}\right)}
$$

$$
\leq 1-\frac{t}{s} \quad\left(\text { By def. of } \theta_{k-1} \text { and property 5. }\right)
$$

Note that for the latter two contractions above, it always holds that $\mu /(4 s)<t / s$ so in the final bound we can ignore the latter. Finally, denoting the LHS of inequality (44) as $E_{k+1}$, we have that:

$$
E_{k+1} \leq(1-\min \{p, q, r\})^{k+1} M
$$

where $M=E_{2}$ and we used the fact that $\theta_{1}=0$.


[^0]:    ${ }^{1}$ See https://github.com/mvladarean/adaptive_pda.

