# Accelerated inexact composite gradient methods for nonconvex spectral optimization problems 

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#### Abstract

This paper presents two inexact composite gradient methods, one inner accelerated and another doubly accelerated, for solving a class of nonconvex spectral composite optimization problems. More specifically, the objective function for these problems is of the form $f_{1}+f_{2}+h$, where $f_{1}$ and $f_{2}$ are differentiable nonconvex matrix functions with Lipschitz continuous gradients, $h$ is a proper closed convex matrix function, and both $f_{2}$ and $h$ can be expressed as functions that operate on the singular values of their inputs. The methods essentially use an accelerated composite gradient method to solve a sequence of proximal subproblems involving the linear approximation of $f_{1}$ and the singular value functions underlying $f_{2}$ and $h$. Unlike other composite gradient-based methods, the proposed methods take advantage of both the composite and spectral structure underlying the objective function in order to efficiently generate their solutions. Numerical experiments are presented to demonstrate the practicality of these methods on a set of real-world and randomly generated spectral optimization problems.


Keywords Composite nonconvex problem • Iteration complexity • Inexact composite gradient method • First-order accelerated gradient method • Spectral optimization

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## 1 Introduction

There are numerous applications in electrical engineering, machine learning, and medical imaging that can be formulated as nonconvex spectral optimization problems of the form

$$
\begin{equation*}
\min _{U \in \mathbb{R}^{n \times n}}\{\phi(U):=f_{1}(U)+(\underbrace{f_{2}^{\mathcal{V}} \circ \sigma}_{f_{2}})(U)+(\underbrace{h^{\nu} \circ \sigma}_{h})(U)\}, \tag{1}
\end{equation*}
$$

where $\sigma$ is the function that maps a matrix to its singular value vector (in nonincreasing order of magnitude), $f_{1}$ and $f_{2}^{\mathcal{V}}$ are continuously differentiable functions with Lipschitz continuous gradients, and $h^{\mathcal{V}}$ is a proper, lower semicontinuous, convex function. For this paper, we are interested in solving instances of (1) where: (i) the resolvents of $\lambda \partial h$ and $\lambda \partial h^{\mathcal{V}}$, i.e., evaluations of the operators $(I+\lambda \partial h)^{-1}$ and $\left(I+\lambda \partial h^{\mathcal{V}}\right)^{-1}$, are easy compute for any $\lambda>0$; (ii) the resolvents of $\lambda\left(\nabla f_{2}+\partial h\right)$ and $\lambda\left(\nabla f_{2}^{\mathcal{V}}+\partial h^{\mathcal{V}}\right)$ cannot be computed exactly for any $\lambda>0$; and (iii) both $f_{2}^{\mathcal{L}}$ and $h^{\mathcal{L}}$ are absolutely symmetric in their arguments, i.e., they do not depend on the ordering or the sign of their arguments.

We now describe some practical instances of (1) that satisfy all three assumptions above. To avoid repetition, we let $\mathcal{R}=\mathcal{R}_{s}+\mathcal{R}_{n}$ and $\mathcal{P}$ be two sparsityinducing regularizers, where $\mathcal{R}_{s}$ and $\mathcal{P}$ are continuously differentiable functions with Lipschitz continuous gradients and $\mathcal{R}_{n}$ is a proper, lower semicontinuous, and convex function.

- Matrix Completion Let $A \in \mathbb{R}^{m \times n}$ be a given data matrix and let $r=\min \{m, n\}$. Moreover, let $\Omega$ denote a subset of the indices of $A$. The goal of the general matrix completion problem is to find a low rank approximation of $A$ that is close to $A$ in some sense. A nonconvex formulation (see, for example, [20]) of this problem is

$$
\min _{X \in \mathbb{R}^{m \times n}}\left\{\frac{1}{2}\left\|P_{\Omega}(X-A)\right\|_{F}^{2}+(\mathcal{R} \circ \sigma)(X)\right\}
$$

where $P_{\Omega}$ is the function that zeros out the entries of its input that are not in $\Omega$. Note that this problem is a special instance of (1) in which $f_{1}=\left\|P_{\Omega}(\cdot)-A\right\|_{F}^{2} / 2$, $f_{2}^{\mathcal{V}}=\mathcal{R}_{s}$, and $h^{\mathcal{V}}=\mathcal{R}_{n}$.

- Phase Retrieval Given a vector $x \in \mathbb{R}^{n}$, let $x[\omega]$ denote its discrete Fourier transform for some frequency $\omega$. Moreover, for some unknown noisy signal $\tilde{x} \in \mathbb{R}^{n}$ and a frequency set $\Omega \subseteq \mathbb{R}_{+}$, suppose that we are given measurements $\{|\tilde{x}[\omega]|\}_{\omega \in \Omega}$ and vectors $a_{\omega} \in \mathbb{C}^{n}$ such that $\left|\left\langle a_{\omega}, \tilde{x}\right\rangle\right|=|\tilde{x}[\omega]|$ for every $\omega \in \Omega$. The goal of the phase retrieval problem is to recover an approximation $x$ of $\tilde{x}$ such that $\left|\left\langle a_{\omega}, x\right\rangle\right|^{2} \approx\left|\left\langle a_{\omega}, \tilde{x}\right\rangle\right|^{2}$ for every $\omega \in \Omega$. A nonconvex formulation of this problem is

$$
\min _{X \in \mathbb{R}^{|\Omega| \times|\times 2|}}\left\{\frac{1}{2}\|\mathcal{A}(X)-b\|^{2}+(\mathcal{R} \circ \lambda)(X): X \geq 0\right\}
$$

where $\lambda$ denotes the function that maps matrices to their eigenvalue vector, $X \geq 0$ means that $X$ is symmetric positive semidefinite, and the quantities $\mathcal{A}: \mathbb{R}^{|\Omega| \times|\Omega|} \mapsto \mathbb{R}^{|\Omega|}$ and $b \in \mathbb{R}^{|\Omega|}$ are given by

$$
[\mathcal{A}(X)]_{\omega}=\operatorname{tr}\left(a_{\omega} a_{\omega}^{*} X\right), \quad b_{\omega}=|\tilde{x}[\omega]|^{2}, \quad \forall(X, \omega) \in \mathbb{R}^{|\Omega| \times|\Omega|} \times \Omega
$$

Note that this problem is a special instance of (1) in which $f_{1}=\|\mathcal{A}(\cdot)-b\|_{F}^{2} / 2$, $f_{2}^{\mathcal{V}}=\mathcal{R}_{s}$, and $h^{\mathcal{V}}=\mathcal{R}_{n}+\delta_{\mathbb{R}_{+}^{|\Omega|}}$ where $\delta_{\mathbb{R}_{+}^{|\Omega|}}$ is the indicator for the nonnegative orthant of $\mathbb{R}^{|\Omega|}$. It is worth mentioning that this formulation is a generalization of the one in [3] where the convex function $\operatorname{tr} X$ is replaced with the nonconvex function $\mathcal{R}$.

- Robust Principal Component Analysis Let $\widehat{M} \in \mathbb{R}^{m \times n}$ be a given data matrix and let $r=\min \{m, n\}$. The goal of the robust principal component analysis problem is to find an approximation $M+E$ of $\widehat{M}$ where $M$ is low-rank and $E$ is sparse. A nonconvex formulation of this problem is

$$
\min _{M, E \in \mathbb{R}^{m \times n}}\left\{\frac{1}{2}\|\widehat{M}-(M+E)\|_{F}^{2}+(\mathcal{R} \circ \sigma)(M)+\mathcal{P}(E)\right\} .
$$

Note that this problem is a special instance of (1) in which $f_{1}=\|\widehat{M}-[(\cdot)-E]\|_{F}^{2} / 2+\mathcal{P}, f_{2}^{\mathcal{V}}=\mathcal{R}_{s}$, and $h^{\mathcal{V}}=\mathcal{R}_{n}$. It is worth mentioning that this formulation is a instance of the one in [19] where more structure is imposed on the functions $\mathcal{R}$ and $\mathcal{P}$.
A natural approach for finding approximate stationary points of the above instances is to employ the exact composite gradient (ECG) method that, when applied to (1), exactly solves a sequence of matrix subproblems of the form

$$
\begin{equation*}
\min _{U \in \mathbb{R}^{m \times n}}\left\{\tilde{\lambda}_{k}\left[\left\langle\nabla\left(f_{1}+f_{2}\right)\left(Y_{k-1}\right), U\right\rangle+h(U)\right]+\frac{1}{2}\left\|U-Y_{k-1}\right\|_{F}^{2}\right\}, \tag{2}
\end{equation*}
$$

where $\lambda_{k}>0$ is an appropriately chosen stepsize and the point $Y_{k-1}$ is the previous iterate. Its computation primarily consists of computing a singular value decomposition (SVD) at the point $\tilde{Y}_{k}:=Y_{k-1}-\tilde{\lambda}_{k} \nabla\left(f_{1}+f_{2}\right)\left(Y_{k-1}\right)$ and an evaluation of the resolvent of $\tilde{\lambda}_{k} \partial h^{\mathcal{V}}$ at $\sigma\left(\tilde{Y}_{k}\right)$. Accelerated ECG (A-ECG) methods solve subproblems similar to (2) but with $Y_{k-1}$ selected in an accelerated manner. Notice that both of these approaches do not exploit the spectral structure in $f_{2}$.

Our goal in this paper is to develop two efficient inexact composite gradient (ICG) methods that find approximate stationary points of (1) by exploiting the spectral structure in both $f_{2}$ and $h$. Our first prototype, called the static inner accelerated ICG (IA-ICG) method, inexactly solves a sequence of matrix prox subproblems of the form

$$
\begin{equation*}
\min _{U \in \mathbb{R}^{n \times n}}\left\{\lambda_{k}\left[\left\langle\nabla f_{1}\left(Y_{k-1}\right), U\right\rangle+f_{2}(U)+h(U)\right]+\frac{1}{2}\left\|U-Y_{k-1}\right\|_{F}^{2}\right\} \tag{3}
\end{equation*}
$$

where $\lambda_{k}>0$ is an appropriately chosen stepsize and the point $Y_{k-1}$ is the previous iterate. It is shown (see Sect. 4.1) that the effort of finding the required inexact solution $Y_{k}$ of (3) consists of computing one SVD and applying an accelerated gradient (ACG) algorithm to find an approximate solution to the related vector prox subproblem

$$
\begin{equation*}
\min _{u \in \mathbb{R}^{r}}\left\{\lambda_{k}\left[f_{2}^{\mathcal{V}}(u)-\left\langle c_{k-1}, u\right\rangle+h^{\mathcal{V}}(u)\right]+\frac{1}{2}\|u\|^{2}\right\} \tag{4}
\end{equation*}
$$

where $r=\min \{m, n\}$ and $c_{k-1}=\sigma\left(Y_{k-1}-\lambda_{k} \nabla f_{1}\left(Y_{k-1}\right)\right)$. Notice that (4) is a problem over the vector space $\mathbb{R}^{r}$, and hence, has significantly fewer dimensions than (3) which is a problem over the matrix space $\mathbb{R}^{m \times n}$. The other prototype, called the static doubly accelerated ICG (DA-ICG), solves a subproblem similar to (3) but with $Y_{k-1}$ selected in an accelerated manner (and hence its qualifier "doubly accelerated"). Notice that the static IA-ICG (resp. DA-ICG) can be viewed as an inexact version of ECG (resp. A-ECG) where, instead of $h$ in (2), the function $f_{2}+h$ is viewed as the composite term, i.e., the part that is not linearized in the subproblems. Moreover, neither IA-ICG nor DA-ICG are able to solve (3) (or its accelerated version) exactly due to assumption (ii) made in the first paragraph of this section.

Motivation of our approach. For high-dimensional instances of (1) where $r=\min \{m, n\}$ is large, we have that the larger the Lipschitz constant of $\nabla f_{2}^{\mathcal{V}}$ is, the better the performance of the ICG methods is compared to the performance of their exact counterparts. This fact immediately follows from the following two claims:
(i) the ICG methods inexactly solve fewer matrix subproblems compared to their exact counterparts when the Lipschitz constant of $\nabla f_{2}^{\mathcal{V}}$ is large; and
(ii) the work of exactly solving (2) or inexactly solving (3) is comparable when $r$ is large.

The justification of claim (i) is as follows. First, recall that the larger the stepizes $\lambda_{k}$ 's (resp. $\tilde{\lambda}_{k}$ ) are, the smaller the number of generated subproblems (3) (resp. (2)) is. Second, the CG stepsizes chosen in either (2) or (3) to guarantee convergence of the underlying CG method are inversely proportional to the Lipschitz constant of the gradient of the function being linearized. Hence, since the inexact CG methods linearize $f_{1}$ only and the exact CG methods linearize both $f_{1}$ and $f_{2}$, claim (i) follows. Some specific applications where the Lipschitz constant of $\nabla f_{2}^{\mathcal{V}}$ may be large in practice can be found, for example, in [1, 18, 20]. The justification for claim (ii) is due to the following two observations: (a) all of the above CG methods require one SVD per subproblem; and (b) when $r$ is large, the computational bottleneck for solving a single subproblem is the aforementioned SVD.

Contributions and Main results To the best of our knowledge, this paper is the first to present ICG methods that exploit both the spectral and composite structure in (1).

When $f_{2}$ is convex or, more generally, a key inequality is satisfied at every iteration of ACG applied to (4), it is shown that for any given $\hat{\rho}>0$, both the static IAICG and the static DA-ICG always obtain a pair $(\hat{Y}, \hat{V})$ satisfying the approximate stationarity condition

$$
\begin{equation*}
\hat{V} \in \nabla f_{1}(\hat{Y})+\nabla f_{2}(\hat{Y})+\partial h(\hat{Y}), \quad\|\hat{V}\| \leq \hat{\rho} \tag{5}
\end{equation*}
$$

by inexactly solving $\mathcal{O}\left(\hat{\rho}^{-2}\right)$ matrix prox subproblems as in (3). If, in addition, $f_{1}$ is convex, it is shown that this bound improves to $\mathcal{O}\left(\hat{\rho}^{-2 / 3}\right)$ for the static DA-ICG method.

When $f_{2}$ is nonconvex, the static IA-ICG and the static DA-ICG may fail to obtain a pair as in (5). To remedy this, we develop dynamic IA-ICG and DA-ICG methods that repeatedly invoke their static counterparts to solve (1) with $\left(f_{1}, f_{2}\right)$ replaced by $\left(f_{1, \xi}, f_{2, \xi}\right)=\left(f_{1}-\xi\|\cdot\|^{2} / 2, f_{2}+\xi\|\cdot\|^{2} / 2\right)$ for strictly increasing values of $\xi>0$. These dynamic versions always obtain a pair as in (5) because: (i) $f_{1}+f_{2}=f_{1, \xi}+f_{2, \xi}$ for every $\xi>0$ and (ii) there always exists $\underline{\xi}>0$ such that $f_{2, \underline{\xi}}$ is convex due to the fact that $\nabla f_{2}$ is Lipschitz continuous.

Numerical experiments are also given to demonstrate the practicality of our proposed methods. More specifically, our experiments demonstrate that the dynamic methods are substantially faster (usually 10x) than other first-order methods at minimizing the primal residual $\|\hat{V}\|$ in terms of runtime.

Related works The earliest complexity analysis of an ACG method for solving nonconvex composite problems like the one in (1) is given in [6]. Building on the results in [6], many other papers [5, 7, 13] have proposed similar ACG-based methods.

Another common approach for solving problems like (1) is to employ an inexact proximal point method where each prox subproblem is constructed to be convex, and hence, solvable by an ACG variant. For example, papers [4, 9, 10, 16] present inner accelerated inexact proximal point methods whereas [12] presents a doubly accelerated inexact proximal point method.

Organization of the paper Section 1.1 gives some notation and basic definitions. Section 2 presents some necessary background material for describing the ICG methods. Section 3 is split into three subsections. The first one precisely describes the problem of interest, while the last two present the IA-ICG and DA-ICG methods. Section 4 describes an efficient way of solving problem (3) by modifying a solution of problem (4). Section 5 presents some numerical results. Section 6 establishes the iteration complexity of the ICG methods. Finally, some auxiliary results are presented in Appendices A to D.

### 1.1 Notation and basic definitions

This subsection provides some basic notation and definitions.
The set of real numbers is denoted by $\mathbb{R}$. The set of non-negative real numbers and the set of positive real numbers is denoted by $\mathbb{R}_{+}$and $\mathbb{R}_{++}$respectively. The set of natural numbers is denoted by $\mathbb{N}$. The set of complex numbers is $\mathbb{C}$. The set of unitary matrices of size $n-$ by $-n$ is $\mathcal{U}^{n}$. For $t>0$, define $\log _{1}^{+}(t):=\max \{1, \log (t)\}$.

Let $\mathbb{R}^{n}$ denote a real-valued $n$-dimensional Euclidean space with norm $\|\cdot\|$. Given a linear operator $A: \mathbb{R}^{n} \mapsto \mathbb{R}^{p}$, the operator norm of $A$ is denoted by $\|A\|:=\sup \left\{\|A z\| /\|z\|: z \in \mathbb{R}^{n}, z \neq 0\right\}$. Using the asymptotic notation $\mathcal{O}$, we denote $\mathcal{O}_{1}(\cdot) \equiv \mathcal{O}(1+\cdot)$.

Let $(m, n) \in \mathbb{N}^{2}$ and let $r=\min \{m, n\}$. Given matrices $X \in \mathbb{R}^{m \times n}$ and $Y \in \mathbb{R}^{n \times n}$, let the quantities $\sigma(X)$ and $\lambda(Y)$ denote the singular values and eigenvalues of $X$ and $Y$, respectively, in nonincreasing order. Let dg : $\mathbb{R}^{r} \mapsto \mathbb{R}^{r \times r}$ and $\operatorname{Dg}: \mathbb{R}^{m \times n} \mapsto \mathbb{R}^{r}$ be given pointwise by

$$
[\operatorname{dg} z]_{i j}=\left\{\begin{array}{l}
z_{i}, \text { if } i=j, \\
0, \text { otherwise },
\end{array} \quad[\operatorname{Dg} Z]_{i}=Z_{i i}\right.
$$

for every $z \in \mathbb{R}^{r}, Z \in \mathbb{R}^{m \times n}$, and $(i, j) \in\{1, \ldots, r\}^{2}$.
The following notation and definitions are for a general complete inner product space $\mathcal{Z}$, whose inner product and its associated induced norm are denoted by $\langle\cdot, \cdot\rangle$ and $\|\cdot\|$ respectively. Let $\psi: \mathcal{Z} \mapsto(-\infty, \infty]$ be given. The effective domain of $\psi$ is denoted by $\operatorname{dom} \psi:=\{x \in \mathcal{Z}: \psi(x)<\infty\}$ and $\psi$ is said to be proper if $\operatorname{dom} \psi \neq \emptyset$. For $\varepsilon \geq 0$, the $\varepsilon$-subdifferential of $\psi$ at $x \in \operatorname{dom} \psi$ is denoted by

$$
\partial_{\varepsilon} \psi(z):=\left\{w \in \mathbb{R}^{n}: \psi\left(z^{\prime}\right) \geq \psi(z)+\left\langle w, z^{\prime}-z\right\rangle-\varepsilon, \forall z^{\prime} \in \mathcal{Z}\right\},
$$

and we denote $\partial \psi \equiv \partial_{0} \psi$. The set of proper, lower semi-continuous, convex functions is denoted by $\overline{\operatorname{Conv}} \mathcal{Z}$. The convex conjugate $\psi$ is denoted by $\psi^{*}$. The linear approximation of $\psi$ at a point $z_{0} \in \operatorname{dom} \psi$ is denoted by $\ell_{\psi}\left(\cdot ; z_{0}\right):=\psi\left(z_{0}\right)+\left\langle\nabla \psi\left(z_{0}\right), \cdot-z_{0}\right\rangle$. The indicator of a closed convex set $C \subseteq \mathcal{Z}$ at a point $z \in \mathcal{Z}$ is denoted by $\delta_{C}(z)$, which is 1 if $z \in C$ and $\infty$ otherwise. The local Lipschitz constant of $\nabla \psi$ at two points $u, z \in \mathcal{Z}$ is denoted by

$$
L_{\psi}(x, y)=\left\{\begin{array}{ll}
\frac{\|\nabla \psi(x)-\nabla \psi(y)\|}{\|x-y\|}, & x \neq y,  \tag{6}\\
0, & x=y,
\end{array} \quad \forall x, y \in \operatorname{dom} \psi\right.
$$

## 2 Background material

Recall from Sect. 1 that our interest is in solving (1) by repeatedly solving a sequence of prox subproblems as in (3). This section presents some background material regarding (3).

This section considers the nonconvex composite optimization (NCO) problem

$$
\begin{equation*}
\min _{u \in \mathcal{Z}}\left\{\psi(u):=\psi_{s}(u)+\psi_{n}(u)\right\}, \tag{7}
\end{equation*}
$$

where $\mathcal{Z}$ is a finite dimensional inner product space and the functions $\psi_{s}$ and $\psi_{n}$ are assumed to satisfy the following assumptions:
(B1) $\psi_{n} \in \overline{\operatorname{Conv}} \mathcal{Z}$;
(B2) $\psi_{s}$ is continuously differentiable on $\mathcal{Z}$ and satisfies $\psi_{s}(u)-\ell_{\psi_{s}}(u ; y)$ $\leq M\|u-y\|^{2} / 2$ for some $M \geq 0$ and every $u, y \in \mathcal{Z}$.

Clearly, problems (1) and (3) are special cases of (7), and hence any definition or result that is stated in the context of (7) applies to (1) and/or (3).

An important notion of an approximate solution of (7) is as follows: given $\hat{\rho}>0$, a pair $\left(y_{r}, v_{r}\right)$ is said to be a $\hat{\rho}$-approximate solution of (7) if

$$
\begin{equation*}
v_{r} \in \nabla \psi_{s}\left(y_{r}\right)+\partial \psi_{n}\left(y_{r}\right), \quad\left\|v_{r}\right\| \leq \hat{\rho} . \tag{8}
\end{equation*}
$$

In Sect. 3, we develop prox-type methods for finding $\hat{\rho}$-approximate solutions of (1) that repeatedly solve (3) inexactly by taking advantage of its spectral decomposition.

We now discuss the inexactness criterion under which the subproblems (3) are solved. Again, the criterion is described in the context of (7) as follows.

Problem $\mathcal{A}$ : Given $(\mu, \theta) \in \mathbb{R}_{++}^{2}$ and $z_{0} \in \mathcal{Z}$, find $(y, v, \varepsilon) \in \operatorname{dom} \psi \times \mathcal{Z} \times \mathbb{R}_{+}$ such that

$$
\begin{equation*}
v \in \partial_{\varepsilon}\left(\psi-\frac{\mu}{2}\|\cdot-y\|^{2}\right)(y), \quad\|v\|^{2}+2 \varepsilon \leq \theta^{2}\left\|y-z_{0}\right\|^{2} . \tag{9}
\end{equation*}
$$

We begin by making three remarks about the above problem. First, if $(y, v, \varepsilon)$ solves Problem $\mathcal{A}$ with $\theta=0$, then $(v, \varepsilon)=(0,0)$, and $z$ is an exact solution of (7). Hence, the output $(y, v, \varepsilon)$ of Problem $\mathcal{A}$ can be viewed as an inexact solution of (7) when $\theta \in \mathbb{R}_{++}$. Second, the input $z_{0}$ is arbitrary for the purpose of this section. However, the two methods described in Sect. 3 for solving (1) repeatedly solve (3) according to Problem $\mathcal{A}$ with the input $z_{0}$ at the $k^{\text {th }}$ iteration determined by the iterates generated at the $(k-1)^{\text {th }}$ iteration. Third, defining the function

$$
\begin{equation*}
\Delta_{\mu}(u ; y, v):=\psi(y)-\psi(u)-\langle v, y-u\rangle+\frac{\mu}{2}\|u-y\|^{2} \quad \forall u \in \operatorname{dom} \psi \tag{10}
\end{equation*}
$$

another way to express the inclusion in (9) is $\Delta_{\mu}(u ; y, v) \leq \varepsilon$ for every $u \in \operatorname{dom} \psi$. Finally, the relaxed ACG (R-ACG) algorithm presented later in this subsection will be shown to solve Problem $\mathcal{A}$ when $\psi_{s}$ is convex. Moreover, it solves a weaker version of Problem $\mathcal{A}$ involving $\Delta_{\mu}$ (see Problem $\mathcal{B}$ later on) whenever $\psi_{s}$ is not convex and as long as some key inequalities are satisfied during its execution.

A technical issue in our analysis in this paper lies in the ability of refining the output of Problem $\mathcal{A}$ to an approximate solution ( $y_{r}, v_{r}$ ) of (7), i.e., one satisfying the inclusion in (8), in which $\left\|v_{r}\right\|$ is nicely bounded. We now present a refinement procedure that addresses this issue.

## Refinement Procedure

```
Input: a triple \(\left(M, \psi_{s}, \psi_{n}\right)\) satisfying (B1)-(B2) and a pair \((y, v) \in \operatorname{dom} \psi_{n} \times \mathcal{Z}\);
Output: a pair ( \(y_{r}, v_{r}\) ) satisfying the inclusion in (8);
```

1. set the quantities

$$
\begin{align*}
& y_{r}=\underset{u \in \mathcal{Z}}{\operatorname{argmin}}\left\{\left\langle\nabla \psi_{s}(y)-v, u\right\rangle+\frac{M}{2}\|u-y\|^{2}+\psi_{n}(u)\right\},  \tag{11}\\
& v_{r}=v+M\left(y-y_{r}\right)+\nabla \psi_{s}\left(y_{r}\right)-\nabla \psi_{s}(y), \tag{12}
\end{align*}
$$

and output $\left(y_{r}, v_{r}\right)$.

The result below presents the key properties of the above procedure. For the sake of brevity, we write $\left(y_{r}, v_{r}\right)=R P(y, v)$ to indicate that the pair $\left(y_{r}, v_{r}\right)$ is the output of the above procedure with inputs $\left(M, \psi_{s}, \psi_{n}\right)$ and $(y, v)$.

Proposition 1 Let $\left(M, \psi_{s}, \psi_{n}\right)$ satisfying assumptions (B1)-(B2) and a triple $(y, v, \varepsilon) \in \operatorname{dom} \psi_{n} \times \mathcal{Z} \times \mathbb{R}_{+}$be given. Moreover, let $\left(y_{r}, v_{r}\right)=R P(y, v)$, denote $L_{\psi_{s}}(\cdot, \cdot)$ simply by $L(\cdot, \cdot)$ where $L_{\psi_{s}}(\cdot, \cdot)$ is as in (6), and let $\Delta_{\mu}$ be as in (10). Then, the following statements hold:
(a) $v_{r} \in \nabla \psi_{s}\left(y_{r}\right)+\partial \psi_{n}\left(y_{r}\right)$;
(b) $\Delta_{\mu}\left(y_{r} ; y, v\right) \geq M\left\|y_{r}-y\right\|^{2} / 2$;
(c) if $\Delta_{\mu}\left(y_{r} ; y, v\right) \leq \varepsilon$ and $(y, v, \varepsilon)$ satisfies the inequality in (9), then

$$
\begin{equation*}
\left\|v_{r}\right\| \leq \theta\left[1+\frac{M+L\left(y, y_{r}\right)}{\sqrt{M}}\right]\left\|y-z_{0}\right\| \tag{13}
\end{equation*}
$$

(d) if $(y, v, \varepsilon)$ solves Problem $\mathcal{A}$, then $\Delta_{\mu}(u ; y, v) \leq \varepsilon$ for every $u \in \operatorname{dom} \psi_{n}$, and, as a consequence, bound (13) holds.

Proof (a) Using the definition of $v_{r}$ and the optimality of $y_{r}$, we have that

$$
v_{r}=v+M\left(y-y_{r}\right)+\nabla \psi_{s}\left(y_{r}\right)-\nabla \psi_{s}(y) \in \nabla \psi_{s}\left(y_{r}\right)+\partial \psi_{n}\left(y_{r}\right) .
$$

(b) Let $(y, v) \in \operatorname{dom} \psi_{n} \times \mathcal{Z}$ be fixed, and define $\widetilde{\psi}_{s}:=\psi_{s}-\langle v, \cdot\rangle$. Using Proposition 19 with $(g, h, L)=\left(\widetilde{\psi}_{s}, \psi_{n}, M\right)$ and $(z, \hat{z})=\left(y, y_{r}\right)$, and the definition of $\Delta_{\mu}$ in (10), we have

$$
\begin{aligned}
\frac{M}{2}\left\|y-y_{r}\right\|^{2} & \leq\left(\widetilde{\psi}_{s}+\psi_{n}\right)(y)-\left(\widetilde{\psi}_{s}+\psi_{n}\right)\left(y_{r}\right) \\
& =\psi(y)-\psi\left(y_{r}\right)-\left\langle v, y-y_{r}\right\rangle \leq \Delta_{\mu}\left(y_{r} ; y, v\right) .
\end{aligned}
$$

(c) Using the assumption that $\Delta_{\mu}\left(y_{r} ; y, v\right) \leq \varepsilon$, part (b), and the inequality in (9), we have that

$$
\begin{equation*}
\left\|y-y_{r}\right\| \leq \sqrt{\frac{2 \Delta_{\mu}\left(y_{r} ; y, v\right)}{M}} \leq \sqrt{\frac{2 \varepsilon}{M}} \leq \frac{\theta}{\sqrt{M}}\left\|y-z_{0}\right\| . \tag{14}
\end{equation*}
$$

Using the triangle inequality, the definition of $L(\cdot, \cdot),(14)$ and the inequality in (9) again, we conclude that

$$
\left\|v_{r}\right\| \leq\|v\|+\left[M+L\left(y, y_{r}\right)\right] \cdot\left\|y-y_{r}\right\| \leq \theta\left[1+\frac{M+L\left(y, y_{r}\right)}{\sqrt{M}}\right]\left\|y-z_{0}\right\|
$$

(d) The fact that $\Delta_{\mu}(u ; y, v) \leq \varepsilon$ for every $u \in \operatorname{dom} \psi_{n}$ follows immediately from the inclusion in (9) and the definition of $\Delta_{\mu}$ in (10). The fact that (13) holds now follows from part (c).

We make a few remarks about Proposition 1. First, it follows from (a) that ( $y_{r}, v_{r}$ ) satisfies the inclusion in (8). Second, it follows from (a) and (c) that if $\theta=0$, then $\left(\varepsilon, v_{r}\right)=(0,0)$, and hence $y_{r}$ is an exact stationary point of (7). In general, (13) implies that the residual $\left\|v_{r}\right\|$ is directly proportional to $\|y-w\|$, and hence, becomes smaller as this quantity approaches zero.

Inequality (13) plays an important technical role in the complexity analysis of the two prox-type methods of Sect. 3. Sufficient conditions for its validity are provided in (c) and (d), with (c) being the weaker one, in view of (d). When $\psi_{s}$ is convex, it is shown that every iterate of the R-ACG algorithm presented below always satisfies the inclusion in (9), and hence, verifying the the validity of the sufficient condition in (c) amounts to simply checking whether the inequality in (9) holds. When $\psi_{s}$ is not convex, verification of the inclusion in (9), and hence the sufficient condition in (d), is generally not possible, while the one in (c) is. This is a major advantage of the sufficient condition in (c), which is exploited in this paper towards the development of adaptive prox-type methods which attempt to approximately solve (7) when $\psi_{s}$ is not convex.

For the sake of future reference, we now state the following problem for finding a triple ( $y, v, \varepsilon$ ) satisfying the sufficient condition in Proposition 1(c). Its statement relies on the refinement procedure preceding Proposition 1.

Problem $\mathcal{B}:$ Given the same inputs as in Problem $\mathcal{A}$, find $(y, v, \varepsilon) \in \operatorname{dom} \psi \times \mathcal{Z} \times \mathbb{R}_{+}$satisfying the inequality in (9) and

$$
\begin{equation*}
\Delta_{\mu}\left(y_{r} ; y, v\right) \leq \varepsilon, \tag{15}
\end{equation*}
$$

where $\Delta_{\mu}(\cdot ; \cdot, \cdot)$ is as in (10) and $y_{r}$ is the first component of the refined pair $\left(y_{r}, v_{r}\right)=R P(y, v)$.

We now state the aforementioned R-ACG algorithm which solves Problem $\mathcal{A}$ when $\psi_{s}$ is convex and solves Problem $\mathcal{B}$ whenever $\psi_{s}$ is not convex and two key inequalities are satisfied, one at every iteration (i.e., (16)) and one at the end of its execution.

## R-ACG Algorithm

Input: a quadruple $\left(\mu, M, \psi_{s}, \psi_{n}\right)$ satisfying (B1)-(B2) and a pair $\left(\theta, z_{0}\right)$;
Output: a triple $(y, v, \varepsilon)$ that solves Problem $\mathcal{B}$ or a failure status;
0 . define $\psi:=\psi_{s}+\psi_{n}$ and set $z_{0}^{c}=z_{0}, B_{0}=0, \Gamma_{0} \equiv 0$, and $j=1$;

1. compute the iterates

$$
\begin{aligned}
\xi_{j-1} & =\frac{1+\mu B_{j-1}}{M-\mu}, \quad b_{j-1}=\frac{\xi_{j-1}+\sqrt{\xi_{j-1}^{2}+4 \xi_{j-1} B_{j-1}}}{2} \\
B_{j} & =B_{j-1}+b_{j-1}, \quad \tilde{z}_{j-1}=\frac{B_{j-1}}{B_{j}} z_{j-1}+\frac{b_{j-1}}{B_{j}} z_{j-1}^{c} \\
z_{j} & =\underset{u \in \mathcal{Z}}{\operatorname{argmin}}\left\{l_{\psi_{s}}\left(u ; \tilde{z}_{j-1}\right)+\psi_{n}(u)+\frac{M}{2}\left\|u-\tilde{z}_{j-1}\right\|^{2}\right\} \\
z_{j}^{c} & =\frac{1}{1+\mu B_{j}}\left[z_{j-1}^{c}-b_{j-1}(M-\mu)\left(\tilde{z}_{j-1}-z_{j}\right)+\mu\left(B_{j-1} z_{j-1}^{c}+b_{j-1} z_{j}\right)\right]
\end{aligned}
$$

2. compute the quantities

$$
\begin{aligned}
\tilde{\gamma}_{j} & =l_{\psi_{s}}\left(\cdot ; \tilde{z}_{j-1}\right)+\psi_{n}+\frac{\mu}{2}\left\|\cdot-\tilde{z}_{j-1}\right\|^{2} \\
\gamma_{j} & =\tilde{\gamma}_{j}\left(z_{j}\right)+(M-\mu)\left\langle\tilde{z}_{j-1}-z_{j}, \cdot-z_{j}\right\rangle+\frac{\mu}{2}\left\|\cdot-z_{j}\right\|^{2}, \\
\Gamma_{j} & =\frac{B_{j-1}}{B_{j}} \Gamma_{j-1}+\frac{b_{j-1}}{B_{j}} \gamma_{j}, \quad r_{j}=\frac{z_{0}^{c}-z_{j}^{c}}{B_{j}}+\mu\left(z_{j}^{c}-z_{j}\right) \\
\eta_{j} & =\max \left\{0, \psi\left(z_{j}\right)-\Gamma_{j}\left(z_{j}^{c}\right)-\left\langle r_{j}, z_{j}-z_{j}^{c}\right\rangle+\frac{\mu}{2}\left\|z_{j}-z_{j}^{c}\right\|^{2}\right\}
\end{aligned}
$$

3. if the inequality

$$
\begin{equation*}
\left(\frac{1}{1+\mu B_{j}}\right)\left\|B_{j} r_{j}+z_{j}-z_{0}\right\|^{2}+2 B_{j} \eta_{j} \leq\left\|z_{j}-z_{0}\right\|^{2} \tag{16}
\end{equation*}
$$

holds, then go to step 4; otherwise, stop with a failure status;
4. if the inequality

$$
\begin{equation*}
\left\|r_{j}\right\|^{2}+2 \eta_{j} \leq \theta^{2}\left\|z_{j}-z_{0}\right\|^{2} \tag{17}
\end{equation*}
$$

holds, then go to step 5; otherwise, go to step 1;
5. set $(y, v, \varepsilon)=\left(z_{j}, r_{j}, \eta_{j}\right)$ and compute $\left(y_{r}, v_{r}\right)=R P\left(z_{j}, r_{j}\right)$; if the condition

$$
\Delta_{\mu}\left(y_{r} ; y, v\right) \leq \varepsilon
$$

holds then stop with a success status and output the triple $(y, v, \varepsilon)$; otherwise, stop with a failure status.

It is well-known (see, for example, [8, Proposition 2.3]) that the scalar $B_{j}$ updated in step 1 satisfies

$$
\begin{equation*}
B_{j} \geq \frac{1}{M} \max \left\{\frac{j^{2}}{4},\left(1+\sqrt{\frac{\mu}{4 M}}\right)^{2(j-1)}\right\} \quad \forall j \geq 1 \tag{18}
\end{equation*}
$$

The next result presents the key properties about the R-ACG algorithm.
Proposition 2 The R-ACG algorithm has the following properties:
(a) it stops with either failure or success in

$$
\begin{equation*}
\mathcal{O}\left(\left[1+\sqrt{\frac{L}{\mu}}\right] \log _{1}^{+}\left[L K_{\theta}\left(1+\mu K_{\theta}\right)\right]\right) \tag{19}
\end{equation*}
$$

iterations, where $K_{\theta}:=1+\sqrt{2} / \theta$;
(b) if it stops with success, then its output $(y, v, \varepsilon)$ solves Problem $\mathcal{B}$;
(c) if $\psi_{s}$ is $\mu$-strongly convex then it always stops with success and its output $(y, v, \varepsilon)$ solves Problem $\mathcal{A}$.

Proof (a) See Appendix B.
(b) This follows from the successful checks in step 4 and 5 of the algorithm.
(c) The fact that the algorithm never stops with failure follows from Proposition 20(c)-(d) in Appendix B. The fact that the algorithm stops with success follows from the previous statement, the successful checks in step 4 and 5 of the algorithm, and the fact that the algorithm stops in a finite number of iterations in part (a).

## 3 Inexact composite gradient methods

This section presents the ICG methods and the general problem that they solve. It contains three subsections. The first one presents the problem of interest and gives a general outline of the ICG methods, the second one presents the IA-ICG method, and the third one presents the DA-ICG method. For the ease of presentation, the proofs in this section are deferred to Sect. 6.

### 3.1 Problem of interest and outline of the methods

This subsection describes the problem that the ICG methods solve and outlines their structure.

Instead of considering problems having the spectral structure mentioned in Sect. 1, this section considers a more general NCO problem where its variable $u$ lies in a finite dimensional inner product space $\mathcal{Z}$ (and, hence, can be either a vector and/or matrix) and presents both ICG methods in this more general setting. Section 4 then presents
a modification of the ACG subroutine used by both ICG methods that drastically improves their efficiency in the setting of the spectral problem (1).

More specifically, this section considers the problem

$$
\begin{equation*}
\min _{u \in \mathcal{Z}}\left[\phi(u):=f_{1}(u)+f_{2}(u)+h(u)\right] \tag{20}
\end{equation*}
$$

where the functions $f_{1}, f_{2}$, and $h$ are assumed to satisfy the following assumptions:
(A1) $h \in \overline{\operatorname{Conv}} \mathcal{Z}$;
(A2) $f_{1}, f_{2}$ are continuously differentiable functions and there exists $\left(m_{1}, M_{1}\right) \in \mathbb{R}^{2}$ and $\left(m_{2}, M_{2}\right) \in \mathbb{R}^{2}$ such that, for $i \in\{1,2\}$, we have

$$
\begin{equation*}
-\frac{m_{i}}{2}\|u-y\|^{2} \leq f_{i}(u)-\ell_{f_{i}}(u ; y) \leq \frac{M_{i}}{2}\|u-y\|^{2} \quad \forall u, y \in \operatorname{dom} h ; \tag{21}
\end{equation*}
$$

(A3) for $i \in\{1,2\}$, we have

$$
\left\|\nabla f_{i}(u)-\nabla f_{i}(y)\right\| \leq L_{i}\|u-y\| \quad \forall u, y \in \operatorname{dom} h,
$$

where $L_{i}:=\max \left\{\left|m_{i}\right|,\left|M_{i}\right|\right\}$;
(A4) $\phi_{*}:=\inf _{u \in \mathcal{Z}} \phi(u)>-\infty$.
Note that assumption (A2) implies that assumption (A3) holds when the interior of dom $h$ is nonempty. Under the above assumptions, the proposed ICG methods find an approximate solution $(\hat{y}, \hat{v})$ of (20) as in (8) with $\psi_{s}=f_{1}+f_{2}$ and $\psi_{n}=h$, i.e.

$$
\begin{equation*}
\hat{v} \in \nabla f_{1}(\hat{y})+\nabla f_{2}(\hat{y})+\partial h(\hat{y}), \quad\|\hat{v}\| \leq \hat{\rho} . \tag{22}
\end{equation*}
$$

We now outline the ICG methods. Given a starting point $y_{0} \in \operatorname{dom} \psi_{n}$ and a special stepsize $\lambda>0$, each method continually calls the R-ACG algorithm of Sect. 2 to find an approximate solution of a prox-linear form of (20). More specifically, each R-ACG call is used to tentatively find an approximate solution of

$$
\begin{equation*}
\min _{u \in \mathcal{Z}}\left[\psi(u)=\lambda\left[\ell_{f_{1}}\left(u ; z_{0}\right)+f_{2}(u)+h(u)\right]+\frac{1}{2}\left\|u-z_{0}\right\|^{2}\right], \tag{23}
\end{equation*}
$$

for some reference point $z_{0}$. For the IA-ICG method, the point $z_{0}$ is $y_{0}$ for the first R-ACG call and is the last obtained approximate solution for the other R-ACG calls. For the DA-ICG method, the point $z_{0}$ is chosen in an accelerated manner.

From the output of the $k^{\text {th }}$ R-ACG call, a refined pair $(\hat{y}, \hat{v})=\left(\hat{y}_{k}, \hat{v}_{k}\right)$ is generated which: (i) always satisfies the inclusion of (22); and (ii) is such that $\min _{i \leq k}\left\|\hat{v}_{i}\right\| \rightarrow 0$ as $k \rightarrow \infty$. More specifically, this refined pair is generated by applying the refinement procedure of Sect. 2 and adding some adjustments to the resulting output to conform with our goal of finding an approximate solution as in (22). For the ease of future reference, we now state this specialized refinement procedure. Before proceeding, we introduce the shorthand notation

$$
\begin{equation*}
M_{i}^{+}:=\max \left\{M_{i}, 0\right\}, \quad m_{i}^{+}:=\max \left\{m_{i}, 0\right\}, \quad L_{i}(x, y):=L_{f_{i}}(x, y), \tag{24}
\end{equation*}
$$

for $i \in\{1,2\}$, to keep its presentation (and future results) concise.

## Specialized Refinement Procedure

Input: a quadruple $\left(M_{2}, f_{1}, f_{2}, h\right)$ satisfying (A1)-(A2), a scalar $\lambda>0$, and a triple $\left(y, v, z_{0}\right) \in$ $\operatorname{dom} \psi_{n} \times \mathcal{Z} \times \mathcal{Z}$;
Output: a pair $(\hat{y}, \hat{v})$ satisfying the inclusion of (22);

1. compute $\left(\hat{y}, v_{r}\right)=R P(y, v)$ using the refinement procedure in Section 2 with

$$
\begin{equation*}
M=\lambda M_{2}^{+}+1, \quad \psi_{s}=\lambda\left[\ell_{f_{1}}\left(\cdot ; z_{0}\right)+f_{2}\right]+\frac{1}{2}\left\|\cdot-z_{0}\right\|^{2}, \quad \psi_{n}=\lambda h \tag{25}
\end{equation*}
$$

2. compute the residual

$$
\hat{v}=\frac{1}{\lambda}\left(v_{r}+z_{0}-y\right)+\nabla f_{1}(\hat{y})-\nabla f_{1}\left(z_{0}\right),
$$

and output $(\hat{y}, \hat{v})$.

The result below states some properties about the above procedure. For the sake of brevity, we write $(\hat{y}, \hat{v})=\operatorname{SRP}\left(y, v, z_{0}\right)$ to indicate that the pair $(\hat{y}, \hat{v})$ is the output of the above procedure with inputs ( $M_{2}, f_{1}, f_{2}, h$ ), $\lambda$, and $\left(y, v, z_{0}\right)$.

Lemma 3 Let $\left(m_{1}, M_{1}\right),\left(m_{2}, M_{2}\right)$, and $\left(f_{1}, f_{2}, h\right)$ satisfying assumptions (A1)-(A3) and a quadruple $\left(z_{0}, y, v, \varepsilon\right) \in \mathcal{Z} \times \operatorname{dom} \psi_{n} \times \mathcal{Z} \times \mathbb{R}_{+}$be given. Moreover, let $(\hat{y}, \hat{v})=\operatorname{SRP}\left(y, v, z_{0}\right)$ and define

$$
\begin{equation*}
C_{\lambda}(x, y):=\frac{1+\lambda\left[M_{2}^{+}+L_{1}(x, y)+L_{2}(x, y)\right]}{\sqrt{1+\lambda M_{2}^{+}}} \tag{26}
\end{equation*}
$$

for every $x, y \in \mathcal{Z}$. Then, the following statements hold:
(a) $\hat{v} \in \nabla f_{1}(\hat{y})+\nabla f_{2}(\hat{y})+\partial h(\hat{y})$;
(b) if $(y, \nu, \varepsilon)$ solves Problem $\mathcal{B}$ with $\left(\mu, \psi_{s}, \psi_{n}\right)$ as in (28), then

$$
\|\hat{v}\| \leq\left[L_{1}(y, w)+\frac{2+\theta C_{\lambda}(y, \hat{y})}{\lambda}\right]\left\|y-z_{0}\right\| .
$$

It is worth recalling from Sect. 1 that in the applications we consider, the cost of the R-ACG call is small compared to SVD computation that is performed before solving each subproblem as in (23). Hence, in the analysis that follows, we present complexity results related to the number of subproblems solved rather than the total number of R-ACG iterations. We do note, however, that the number of R-ACG iterations per subproblem is finite in view of Proposition 2(a).

### 3.2 Static and dynamic IA-ICG methods

This subsection presents the static and dynamic IA-ICG methods.
We first state the static IA-ICG method.

[^1]Note that the static IA-ICG method may fail without obtaining a pair satisfying (22). In Theorem 4(c) below, we state that a sufficient condition for the method to stop successfully is that $f_{2}$ be convex. This property will be important when we present the dynamic IA-ICG method, which: (i) repeatedly calls the static method; and (ii) incrementally transfers convexity from $f_{1}$ to $f_{2}$ between each call until a successful termination is achieved.

We now make some additional remarks about the above method. First, it performs two kinds of iterations, namely, ones that are indexed by $k$ and ones that are performed by the R-ACG algorithm. We refer to the former kind as outer iterations and the latter kind as inner iterations. Second, in view of (27), if $M_{1}>0$ then $0<\lambda<\left(1-2 \theta^{2}\right) /\left(2 M_{1}\right)$ whereas if $M_{1} \leq 0$ then $0<\lambda<\infty$. Finally, the most expensive part of the method is the R-ACG call in step 1. In Sect. 4, we show that this call can be replaced with a call to a spectral version of R-ACG that is dramatically more efficient when the underlying problem has the spectral structure as in (1).

The next result summarizes some facts about the static IA-ICG method. Before proceeding, we first define some useful quantities. For $\lambda>0$ and $u, w \in \mathcal{Z}$, define

$$
\begin{equation*}
\tilde{\ell}_{\phi}(u ; w):=\ell_{f_{1}}(u ; w)+f_{2}(u)+h(u), \quad \bar{C}_{\lambda}:=\frac{1+\lambda\left(M_{2}^{+}+L_{1}+L_{2}\right)}{\sqrt{1+\lambda M_{2}^{+}}} . \tag{29}
\end{equation*}
$$

Theorem 4 The following statements hold about the static IA-ICG method:
(a) it stops in

$$
\begin{equation*}
\mathcal{O}_{1}\left(\left[\sqrt{\lambda} L_{1}+\frac{1+\theta \bar{C}_{\lambda}}{\sqrt{\lambda}}\right]^{2}\left[\frac{\phi\left(z_{0}\right)-\phi_{*}}{\hat{\rho}^{2}}\right]\right) \tag{30}
\end{equation*}
$$

outer iterations, where $\phi_{*}$ is as in (A4);
(b) if it stops with success, then its output pair $(\hat{y}, \hat{v})$ is a $\hat{\rho}$-approximate solution of (20);
(c) if $f_{2}$ is convex, then it always stops with success.

We now make three remarks about the above results. First, if $\theta=\mathcal{O}\left(1 / \bar{C}_{\lambda}\right)$ then (30) is on the order of

$$
\begin{equation*}
\mathcal{O}_{1}\left(\left[\sqrt{\lambda} L_{1}+\frac{1}{\sqrt{\lambda}}\right]^{2}\left[\frac{\phi\left(z_{0}\right)-\phi_{*}}{\hat{\rho}^{2}}\right]\right) . \tag{31}
\end{equation*}
$$

Moreover, comparing the above complexity to the iteration complexity of the ECG method described in Sect. 1, which is known (see, for example, [16]) to obtain an approximate solution of (20) in

$$
\begin{equation*}
\mathcal{O}_{1}\left(\left[\sqrt{\lambda}\left(L_{1}+L_{2}\right)+\frac{1}{\sqrt{\lambda}}\right]^{2}\left[\frac{\phi\left(z_{0}\right)-\phi_{*}}{\hat{\rho}^{2}}\right]\right) \tag{32}
\end{equation*}
$$

iterations, we see that (31) is smaller than (32) in magnitude when $L_{2}$ is large. Notice also that the complexity in (31) corresponds to applying the ECG method to (1) where the composite function is $f_{2}+h$ instead of just $h$. Second, Theorem 4(b) shows that if the method stops with success, regardless of the convexity of $f_{2}$, then its output pair $(\hat{y}, \hat{v})$ is always an approximate solution of (20). Third, in view of Proposition 10 , the quantities $L_{1}$ and $\bar{C}_{\lambda}$ in all of the previous complexity results can be replaced by their averaged counterparts in (48). As these averaged quantities only depend on $\left\{\left(y_{i}, \hat{y}_{i}\right)\right\}_{i=1}^{k}$, we can infer that the static IA-ICG method adapts to the local geometry of its input functions.

We now state the dynamic IA-ICG method that resolves the issue of failure in the static IA-ICG method.

## Dynamic IA-ICG Method

Input: the same as the static IA-ICG method but with an additional parameter $\xi_{0}>0$;
Output: a pair ( $\hat{y}, \hat{v}$ ) satisfying (22);
0 . set $\xi=\xi_{0}, \ell=1$, and

$$
\begin{align*}
f_{1}=f_{1}-\frac{\xi}{2}\|\cdot\|^{2}, & f_{2}=f_{2}+\frac{\xi}{2}\|\cdot\|^{2},  \tag{33}\\
m_{1}=m_{1}+\xi, \quad M_{1}=M_{1}-\xi, & m_{2}=m_{2}-\xi, \quad M_{2}=M_{2}+\xi
\end{align*}
$$

1. call the static IA-ICG method with inputs $\left(f_{1}, f_{2}, h\right),\left(m_{1}, M_{1}, m_{2}, M_{2}\right), \hat{\rho}$, $y_{0}$, and $(\lambda, \theta)$;
2. if the static IA-ICG call stops with a failure status, then set $\xi=2 \xi$, update the quantities in (33) with the new value of $\xi$, increment $\ell=\ell+1$, and go to step 1 ; otherwise, let ( $\hat{y}, \hat{v}$ ) be the output pair returned by the static IA-ICG call, stop, and output this pair.

Some remarks about the above method are in order. First, in view of (27) and the fact that $M_{1}$ is monotonically decreasing, the parameter $\lambda$ does not need to be changed for each IA-ICG call. Second, in view of assumption (A2) and Theorem 4(c), the IA-ICG call in step 1 always terminates with success whenever $m_{2} \leq 0$. As a consequence, the total number of IA-ICG calls is at most $\left\lceil\log \left(2 m_{2}^{+} / \xi_{0}\right)\right\rceil$. Third, in view of the second remark and Theorem 4(b), the method always obtains a $\hat{\rho}$-approximate solution of (20) in a finite number of IA-ICG outer iterations. Finally, in view of second remark again, the total number of IA-ICG outer iterations is as in Theorem 4(a) but with: (i) an additional multiplicative factor of $\left\lceil\log \left(2 m_{2}^{+} / \xi_{0}\right)\right\rceil$; and (ii) the constants $m_{1}$ and $M_{2}$ replaced with $\left(m_{1}+2 m_{2}^{+}\right)$and $\left(M_{2}+2 m_{2}^{+}\right)$, respectively. It is worth mentioning that a more refined analysis, such as the one in [10], can be applied in order to remove the factor of $\left[\log \left(2 m_{2}^{+} / \xi_{0}\right)\right\rceil$ from the previously mentioned complexity.

### 3.3 Static and dynamic DA-ICG methods

This subsection presents the static DA-ICG method, but omits the statement of its dynamic variant for the sake of brevity. We do argue, however, that the dynamic variant can be stated in the same way as the dynamic IA-ICG method of Sect. 6.1 but with the call to the static IA-ICG method replaced with a call to the static DA-ICG method of this subsection.

We start by stating some additional assumptions. It is assumed that:
(i) the set $\operatorname{dom} h$ is closed;
(ii) there exists a bounded set $\Omega \supseteq \operatorname{dom} h$ for which a projection oracle exists.

We now state the static DA-ICG method.

## Static DA-ICG Method

Input: function triple $\left(f_{1}, f_{2}, h\right)$ and scalar quadruple $\left(m_{1}, M_{1}, m_{2}, M_{2}\right) \in \mathbb{R}^{4}$ satisfying (A1)-(A4), tolerance $\hat{\rho}>0$, initial point $y_{0} \in \operatorname{dom} h$, and scalar pair $(\lambda, \theta) \in \mathbb{R}_{++} \times(0,1)$ satisfying

$$
\begin{equation*}
\lambda M_{1}+\theta^{2} \leq \frac{1}{2} \tag{34}
\end{equation*}
$$

Output: a pair $(\hat{y}, \hat{v})$ satisfying (22) or a failure status;
0 . let $\Delta_{1}(\cdot ; \cdot, \cdot)$ be as in (10) with $\mu=1$, and set $A_{0}=0, x_{0}=y_{0}$, and $k=1$;

1. compute the quantities

$$
\begin{align*}
& a_{k-1}=\frac{1+\sqrt{1+4 A_{k-1}}}{2}, \quad A_{k}=A_{k-1}+a_{k-1} \\
& \tilde{x}_{k-1}=\frac{A_{k-1} y_{k-1}+a_{k-1} x_{k-1}}{A_{k}} \tag{35}
\end{align*}
$$

2. use the R-ACG algorithm to tentatively solve Problem $\mathcal{B}$ associated with (23), i.e., with inputs $\left(\mu, M, \psi_{s}, \psi_{n}\right)$ and $\left(\theta, z_{0}\right)$ where the former is as in (28) and $z_{0}=\tilde{x}_{k-1}$; if the R-ACG stops with success, then let $\left(y_{k}^{a}, v_{k}, \varepsilon_{k}\right)$ denote its output and go to step 3 ; otherwise, stop with a failure status;
3. if the inequality $\Delta_{1}\left(y_{k-1} ; y_{k}^{a}, v_{k}\right) \leq \varepsilon_{k}$ holds, then go to step 4 ; otherwise, stop with a failure status;
4. set $\left(\hat{y}_{k}, \hat{v}_{k}\right)=S R P\left(y_{k}^{a}, v_{k}, \tilde{x}_{k-1}\right)$ where $\operatorname{SRP}(\cdot, \cdot, \cdot)$ is described in Subsection 3.1; if $\left\|\hat{v}_{k}\right\| \leq \hat{\rho}$ then stop with a success status and output $(\hat{y}, \hat{v})=\left(\hat{y}_{k}, \hat{v}_{k}\right)$; otherwise, compute

$$
\begin{align*}
x_{k} & =\underset{u \in \Omega}{\operatorname{argmin}} \frac{1}{2}\left\|u-\left[x_{k-1}-a_{k-1}\left(v_{k}+\tilde{x}_{k-1}-y_{k}^{a}\right)\right]\right\|^{2}, \\
y_{k} & =\underset{u \in\left\{y_{k-1}, y_{k}^{a}\right\}}{\operatorname{argmin}}\left[f_{1}(u)+f_{2}(u)+h(u)\right], \tag{36}
\end{align*}
$$

update $k \leftarrow k+1$, and go to step 1 .

Note that, similar to the static IA-ICG method, the static DA-ICG method may fail without obtaining a pair satisfying (22). Proposition 5(c) shows that a sufficient condition for the method to stop successfully is that $f_{2}$ be convex. Using arguments similar to the ones employed to derive the dynamic IA-ICG method, a dynamic version of DA-ICG method can also be developed that repeatedly invokes the static DA-ICG in place of the static IA-ICG.

We now make some additional remarks about the above method. First, it performs two kinds of iterations, namely, ones that are indexed by $k$ and ones that are performed by the R-ACG algorithm. We refer to the former kind as outer iterations and the latter kind as inner iterations. Second, in view of the update for $y_{k}$ in (36), the collection of function values $\left\{\phi\left(y_{i}\right)\right\}_{i=0}^{k}$ is non-increasing. Third, in view of (34), if $M_{1}>0$ then $0<\lambda<\left(1-2 \theta^{2}\right) /\left(2 M_{1}\right)$ whereas if $M_{1} \leq 0$ then $0<\lambda<\infty$. Finally, the most expensive part of the method is the R-ACG call in step 2. In Sect. 4, we show that this call can be replaced with a call to a spectral
version of R-ACG that is dramatically more efficient when the underlying problem has the spectral structure as in (1).

It is worth mentioning that the outer iteration scheme of the DA-ICG method is a monotone and inexact generalization of the A-ECG method in [6]. More specifically, this A-ECG method is a version of the DA-ICG method where: (i) $\theta=0$; (ii) the R-ACG algorithm in step 2 is replaced by an exact solver of (23); (iii) the update of $x_{k}$ in (36) is replaced by an update involving the prox evaluation of the function $a_{k-1} h$; and (iv) both $f_{1}$ and $f_{2}$ are linearized instead of just $f_{2}$ in the DA-ICG method. Hence, the DA-ICG method can be significantly more efficient when its R-ACG call is more efficient than an exact solver of (23) and/or when the projection onto $\Omega$ is more efficient than evaluating the prox of $a_{k-1} h$.

The next result summarizes some facts about the DA-ICG method. Before proceeding, we introduce the useful constants

$$
\begin{align*}
D_{h} & :=\sup _{u, z \in \operatorname{dom} h}\|u-z\|, \quad D_{\Omega}:=\sup _{u, z \in \Omega}\|u-z\|, \quad \Delta_{\phi}^{0}:=\phi\left(y_{0}\right)-\phi_{*}, \\
d_{0} & :=\inf _{u^{*} \in \mathcal{Z}}\left\{\left\|y_{0}-u^{*}\right\|: \phi\left(u^{*}\right)=\phi_{*}\right\}, \quad E_{\lambda, \theta}:=\sqrt{\lambda} L_{1}+\frac{1+\theta \bar{C}_{\lambda}}{\sqrt{\lambda}} . \tag{37}
\end{align*}
$$

Theorem 5 The following statements hold about the static DA-ICG method:
(a) it stops in

$$
\begin{equation*}
\mathcal{O}_{1}\left(\frac{E_{\lambda, \theta}^{2}\left[m_{1}^{+} D_{h}^{2}+\Delta_{\phi}^{0}\right]}{\hat{\rho}^{2}}+\frac{E_{\lambda, \theta}\left[m_{1}^{+}+1 / \lambda\right]^{1 / 2} D_{\Omega}}{\hat{\rho}}\right) \tag{38}
\end{equation*}
$$

outer iterations;
(b) if it stops with success, then its output pair $(\hat{y}, \hat{v})$ is a $\hat{\rho}$-approximate solution of (20);
(c) if $f_{2}$ is convex, then it always stops with success in

$$
\begin{equation*}
\mathcal{O}_{1}\left(\frac{E_{\lambda, \theta}^{2} m_{1}^{+} D_{h}^{2}}{\hat{\rho}^{2}}+\frac{E_{\lambda, \theta}\left[m_{1}^{+}\right]^{1 / 2} D_{\Omega}}{\hat{\rho}}+\frac{E_{\lambda, \theta}^{2 / 3} d_{0}^{2 / 3} \lambda^{-1 / 3}}{\hat{\rho}^{2 / 3}}\right) \tag{39}
\end{equation*}
$$

outer iterations.
We now make three remarks about the above results. First, in the "best" scenario of $\max \left\{m_{1}, m_{2}\right\} \leq 0$, i.e., $f_{1}$ and $f_{2}$ are convex, we have that (39) reduces to

$$
\mathcal{O}_{1}\left(\left[L_{1}+\frac{1}{\lambda}\right]^{2 / 3}\left[\frac{d_{0}^{2 / 3}}{\hat{\rho}^{2 / 3}}\right]\right)
$$

which has a smaller dependence on $\hat{\rho}$ when compared to (31). In the "worst" scenario of $\min \left\{m_{1}, m_{2}\right\}>0$, if we take $\theta=\mathcal{O}\left(1 / \bar{C}_{\lambda}\right)$, then (38) reduces to

$$
\mathcal{O}_{1}\left(\left[\sqrt{\lambda} L_{1}+\frac{1}{\sqrt{\lambda}}\right]^{2}\left[\frac{m_{1}^{+} D_{h}^{2}+\phi\left(y_{0}\right)-\phi_{*}}{\hat{\rho}^{2}}\right]\right)
$$

which has the same dependence on $\hat{\rho}$ as in (31). Second, part (c) shows that if the method stops with an output pair $(\hat{y}, \hat{v})$, regardless of the convexity of $f_{2}$, then that pair is always an approximate solution of (20). Third, in view of Proposition 18, the quantities $L_{1}$ and $\bar{C}_{\lambda}$ in all of the previous complexity results can be replaced by their averaged counterparts in (63). As these averaged quantities only depend on $\left\{\left(y_{i}^{a}, \hat{y}_{i}, \tilde{x}_{i-1}\right)\right\}_{i=1}^{k}$, we can infer that the static DA-ICG method, like the static IA-ICG method of the previous subsection, also adapts to the local geometry of its input functions.

## 4 Exploiting the spectral decomposition

Recall that at every outer iteration of the ICG methods in Sect. 3, a call to the R-ACG algorithm is made to tentatively solve Problem $\mathcal{B}$ (see Sect. 3.1) associated with (23). Our goal in this section is to present a more efficient version of R-ACG (based on the idea outlined in Sect. 1) when the underlying problem has the spectral structure as in (1).

The content of this section is divided into two subsections. The first one presents the aforementioned algorithm, whereas the second one proves its key properties.

### 4.1 Spectral R-ACG algorithm

This subsection presents the R-ACG algorithm mentioned above. Throughout our presentation, we let $Z_{0}$ represent the starting point given to the R-ACG algorithm by the two ICG methods.

We first state the aforementioned efficient algorithm.

Spectral R-ACG Algorithm
Input: a quadruple $\left(M_{2}, f_{1}, f_{2}^{\mathcal{V}}, h^{\mathcal{V}}\right)$ satisfying (A1)-(A3) with $\left(f_{2}, h\right)=\left(f_{2}^{\mathcal{V}}, h^{\mathcal{V}}\right)$ and a triple $\left(\lambda, \theta, Z_{0}\right)$;
Output: a triple $(Y, V, \varepsilon)$ that solves Problem $\mathcal{B}$ associated with (23) or a failure status;

1. compute

$$
\begin{equation*}
Z_{0}^{\lambda}:=Z_{0}-\lambda \nabla f_{1}\left(Z_{0}\right), \tag{40}
\end{equation*}
$$

and a pair $(P, Q) \in \mathcal{U}^{m} \times \mathcal{U}^{n}$ satisfying $Z_{0}^{\lambda}=P\left[\operatorname{dg} \sigma\left(Z_{0}^{\lambda}\right)\right] Q^{*}$;
2. use the R-ACG algorithm to tentatively solve Problem $\mathcal{B}$ associated with (4), i.e., with inputs $\left(\mu, M, \psi_{s}^{\mathcal{V}}, \psi_{n}^{\mathcal{V}}\right)$ and $\left(\theta, z_{0}\right)$ where the former is given by

$$
\begin{gather*}
\mu=1, \quad M:=\lambda M_{2}^{+}+1, \\
\psi_{s}^{\mathcal{V}}:=\lambda f_{2}^{\mathcal{V}}-\left\langle\sigma\left(Z_{0}^{\lambda}\right), \cdot\right\rangle+\frac{1}{2}\|\cdot\|^{2}, \quad \psi_{n}^{\mathcal{V}}:=\lambda h^{\mathcal{V}}, \tag{41}
\end{gather*}
$$

and $z_{0}=\operatorname{Dg}\left(P^{*} Z_{0} Q\right)$; if the R-ACG stops with success, then let $(y, v, \varepsilon)$ denote its output and go to step 3; otherwise, stop with a failure status;
3. set $Y=P(\operatorname{dg} y) Q^{*}$ and $V=P(\operatorname{dg} v) Q^{*}$, and output the triple $(Y, V, \varepsilon)$.

We now make three remarks about the above algorithm. First, the matrices $P$ and $Q$ in step 1 can be obtained by computing an SVD of $Z_{0}^{\lambda}$. Second, in view of Proposition 20(a) and the fact that $(\mu, M)$ in (41) and (28) are the same, the iteration complexity is the same as the vanilla R-ACG algorithm. Finally, because the functions $\psi_{s}^{\mathcal{V}}$ and $\psi_{n}^{\mathcal{V}}$ in (41) have vector inputs over $\mathbb{R}^{r}$, the steps in the spectral R-ACG algorithm are significantly less costly than the ones in the R-ACG algorithm, which use functions with matrix inputs over $\mathbb{R}^{m \times n}$.

The following result, whose proof is in the next subsection, presents the key properties of this algorithm.

Proposition 6 The spectral $R$-ACG algorithm has the following properties:
(a) if it stops with success, then its output triple $(Y, V, \varepsilon)$ solves Problem $\mathcal{B}$ associated with (23);
(b) if $f_{2}$ is convex, then it always stops with success and its output $(Y, V, \varepsilon)$ solves Problem $\mathcal{A}$ associated with (23).

### 4.2 Proof of Proposition 6

For the sake of brevity, let $\left(\psi_{s}, \psi_{n}\right)$ be as in (28) and, using $P$ and $Q$ from the spectral R-ACG algorithm, define for every $(u, U) \in \mathbb{R}^{r} \times \mathbb{R}^{m \times n}$, the functions

$$
\begin{aligned}
& \mathcal{M}(u):=P(\operatorname{dg} u) Q^{*}, \quad \mathcal{V}(U):=\operatorname{Dg}\left(P^{*} U Q\right), \\
& \psi(U):=\psi_{s}(U)+\psi_{n}(U), \quad \psi^{\mathcal{V}}(u):=\psi_{s}^{\mathcal{V}}(u)+\psi_{n}^{\mathcal{V}}(u) .
\end{aligned}
$$

The first result relates $\left(\psi_{s}, \psi_{n}\right)$ to $\left(\psi_{s}^{\mathcal{V}}, \psi_{n}^{\mathcal{V}}\right)$.
Lemma 7 Let $(y, v, \varepsilon)$ and $(Y, V)$ be as in the spectral $R$ - $A C G$ algorithm. Then, the following properties hold:
(a) we have

$$
\psi_{n}^{\mathcal{V}}(y)=\psi_{n}(Y), \quad \psi_{s}^{\mathcal{V}}(y)+B_{0}^{\lambda}=\psi_{s}(Y),
$$

$$
\text { where } B_{0}^{\lambda}:=\lambda f_{1}\left(Z_{0}\right)-\lambda\left\langle\nabla f_{1}\left(Z_{0}\right), Z_{0}\right\rangle+\left\|Z_{0}\right\|_{F}^{2} / 2
$$

(b) we have

$$
\begin{equation*}
V \in \partial_{\varepsilon}\left(\psi-\frac{1}{2}\|\cdot-Y\|_{F}^{2}\right)(Y) \Longleftrightarrow v \in \partial_{\varepsilon}\left(\psi^{\mathcal{V}}-\frac{1}{2}\|\cdot-y\|^{2}\right)(y) . \tag{42}
\end{equation*}
$$

Proof (a) The relationship between $\psi_{n}^{\mathcal{V}}$ and $\psi_{n}$ is immediate. On the other hand, using the definitions of $Y, f_{2}$, and $B_{0}^{\lambda}$, we have

$$
\begin{aligned}
& \psi_{s}^{\nu}(y)+B_{0}^{\lambda}=\lambda f_{2}(Y)-\left\langle Z_{0}^{\lambda}, Y\right\rangle+\frac{1}{2}\|Y\|_{F}^{2}+B_{0}^{\lambda} \\
& \quad=\lambda\left[f_{2}(Y)+f_{1}\left(Z_{0}\right)+\left\langle\nabla f_{1}\left(Z_{0}\right), Y-Z_{0}\right\rangle\right]+\frac{1}{2}\left\|Y-Z_{0}\right\|_{F}^{2}=\psi_{s}(Y)
\end{aligned}
$$

(b) Let $S_{0}=V+Z_{0}^{\lambda}-Y$ and $s_{0}=v+\sigma\left(Z_{0}^{\lambda}\right)-y$, and note that $S_{0}=\mathcal{M}\left(s_{0}\right)$. Moreover, in view of part (a) and the definition of $\psi$, observe that the left inclusion in (42) is equivalent to $S_{0} \in \partial_{\varepsilon}\left(\lambda\left[f_{2}+h\right]\right)(Y)$. Using this observation, the fact that $S_{0}$ and $Y$ have a simultaneous SVD, and Theorem 23 with $(S, s)=\left(S_{0}, s_{0}\right), \Psi=\lambda\left[f_{2}+h\right]$, and $\Psi^{\mathcal{V}}=\lambda\left[f_{2}^{\mathcal{V}}+h^{\mathcal{L}}\right]$, we have that the left inclusion in (42) is also equivalent to $s_{0} \in \partial_{\varepsilon}\left(\lambda\left[f_{2}^{\mathcal{L}}+h^{\mathcal{\nu}}\right]\right)(y)$. The conclusion now follows from the observation that the latter inclusion is equivalent to the the right inclusion in (42).

We are now ready to give the proof of Proposition 6.
Proof of Proposition 6 (a) Since $(y, v)=(\mathcal{V}(Y), \mathcal{V}(V))$, notice that the successful termination of the algorithm implies that the inequality in (9) and (15) hold. Using this remark, the fact that $\|V\|_{F}^{2}=\|v\|^{2}$, and the bound

$$
\begin{align*}
& \theta^{2}\left\|z_{j}-z_{0}\right\|^{2}=\theta^{2}\left(\left\|z_{j}\right\|^{2}-2\left\langle z_{j}, \mathcal{V}\left(z_{0}\right)\right\rangle+\left\|Z_{0}\right\|_{F}^{2}\right)+\theta^{2}\left(\left\|\mathcal{V}\left(z_{0}\right)\right\|^{2}-\left\|Z_{0}\right\|_{F}^{2}\right) \\
& \quad \leq \theta^{2}\left(\left\|Z_{j}\right\|_{F}^{2}-2\left\langle Z_{j}, Z_{0}\right\rangle+\left\|Z_{0}\right\|_{F}^{2}\right)=\theta^{2}\left\|Z_{j}-Z_{0}\right\|_{F}^{2} \tag{43}
\end{align*}
$$

we then have that the inequality in (9) also holds with $(y, v)=(Y, V)$.

To show the corresponding inequality for (15), let $\left(Y_{r}, V_{r}\right)=R P(Y, V)$ using the refinement procedure in Sect. 2. Moreover, let $\left(y_{r}, v_{r}\right)=R P(y, v)$ and $\Delta_{1}^{\mathcal{V}}(\cdot ; \cdot, \cdot)$ be as in (10), where $\left(\psi_{s}, \psi_{n}\right)=\left(\psi_{s}^{\mathcal{V}}, \psi_{n}^{\mathcal{V}}\right)$. It now follows from (11), (12), Lemma 22 with $\Psi=\psi_{n}$ and $S=V+M Y-\nabla \psi_{s}(Y)$, and Lemma 21(b) that $Y_{r}, Y, V$, and $V_{r}$ have a simultaneous SVD. As a consequence of this, the first remark, and Lemma 7(a), we have that

$$
\begin{aligned}
\varepsilon \geq & \Delta_{1}^{\mathcal{V}}\left(y_{r} ; y, v\right)=\psi^{\mathcal{V}}(y)-\psi^{\mathcal{V}}\left(y_{r}\right)-\left\langle v, y-y_{r}\right\rangle+\frac{1}{2}\left\|y_{r}-y\right\|^{2} \\
& =\psi(Y)-\psi\left(Y_{r}\right)-\left\langle V, Y-Y_{r}\right\rangle+\frac{1}{2}\left\|Y_{r}-Y\right\|^{2}=\Delta_{1}\left(Y_{r} ; Y, V\right),
\end{aligned}
$$

and hence that (15) holds with $(y, v)=(Y, V)$.
(b) This follows from part (a), Proposition 2(c), and Lemma 7(b).

## 5 Computational results

This section presents computational results that highlight the performance of the dynamic IA-ICG and dynamic DA-ICG methods, and it contains three subsections. The first one describes the implementation details, the second presents computational results related to a set of spectral composite problem, while the third gives some general comments about the computational results.

### 5.1 Implementation details

This subsection precisely describes the implementation of the methods and experiments of this section. Moreover, all of the code needed to replicate these experiments is readily available online. ${ }^{1}$

We first describe some practical modifications to the dynamic IA-ICG method. Given $\lambda>0$ and $\left(z_{j}, z_{0}\right) \in \mathcal{Z}^{2}$, denote

$$
\Delta_{\phi}^{\lambda}=4 \lambda\left[\phi\left(z_{0}\right)-\widetilde{\ell}_{\phi}\left(z_{j} ; z_{0}\right)-\frac{M_{1}}{2}\left\|z_{j}-z_{0}\right\|^{2}\right]
$$

where $\widetilde{\ell}_{\phi}$ is as in (29). Motivated by the first inequality in the descent condition (46), we relax (17) in the R-ACG call to the three separate conditions: $\left\|z_{j}-z_{0}\right\|^{2} \leq \Delta_{\phi}^{\lambda},\left\|r_{j}\right\|^{2} \leq \Delta_{\phi}^{\lambda}$, and $2 \eta_{j} \leq \Delta_{\phi}^{\lambda}$.

We now describe some modifications and parameter choices that are common to both methods. First, both ICG methods use the spectral R-ACG algorithm of Sect. 4.1 in place of the R-ACG algorithm of Sect. 2. Moreover, this R-ACG variant uses a line search subroutine for estimating the upper curvature $M$ that is used

[^2]during its execution. Second, when each of the dynamic ICG methods invokes their static counterparts, the parameters $A_{0}$ and $y_{0}$ are set to be the last obtained parameters of the previous invocation or the original parameters if it is the first invocation, i.e., we implement a warm-start strategy. Third, we adaptively update $\lambda$ at each outer iteration as follows: given the old value of $\lambda=\lambda_{\text {old }}$ at the $k^{\text {th }}$ outer iteration, the new value of $\lambda=\lambda_{\text {new }}$ at the $(k+1)^{\text {th }}$ iteration is given by
\[

\lambda_{new}=\left\{$$
\begin{array}{ll}
\lambda_{\text {old }}, & r_{k} \in[0.5,2.0], \\
\lambda_{\text {old }} \cdot \sqrt{0.5}, & r_{k}<0.5, \\
\lambda_{\text {old }} \cdot \sqrt{2}, & r_{k}>2.0,
\end{array}
$$ \quad r_{k}=\frac{\left[\lambda\left(M_{2}^{+}+2 m_{2}^{+}\right)+1\right]\left\|y_{k}-\hat{y}_{k}\right\|}{\left\|\hat{v}_{k}-\left[\lambda\left(M_{2}^{+}+2 m_{2}^{+}\right)+1\right]\left(y_{k}-\hat{y}_{k}\right)\right\|} .\right.
\]

Fourth, we take $\mu=1 / 2$ rather than $\mu=1$ for each of R-ACG calls in order to reduce the possibility of a failure from the R-ACG algorithm. Fifth, in view of (43), we relax condition (17) in the vector-based R-ACG call of Sect. 4.1 to

$$
\left\|r_{j}\right\|^{2}+2 \eta_{j} \leq \theta^{2}\left\|z_{j}-z_{0}\right\|^{2}+\tau
$$

where $\tau:=\theta^{2}\left(\left\|Z_{0}\right\|_{F}^{2}-\left\|z_{0}\right\|^{2}\right) \geq 0$. Finally, both ICG methods choose the common hyperparameters $\left(\xi_{0}, \lambda, \theta\right)=\left(M_{1}, 5 / M_{1}, 1 / 2\right)$ at initialization.

We now describe the five other benchmark methods considered. Throughout their descriptions, we let $m=m_{1}+m_{2}, M=M_{1}+M_{2}$, and $L=\max \{m, M\}$. The first method is Nesterov's efficient ECG method of [15] with $\left(\lambda, \gamma_{u}, \gamma_{d}\right)=(100 / L, 2,2)$. The second method is the accelerated inexact proximal point (AIPP) method of [10] with $(\lambda, \theta, \tau)=(1 / m, 4,10[\lambda M+1])$ and the R-AIPPv2 stepsize scheme. The third method is a variant of the A-ECG method of [6, Algorithm 2], which we abbreviate as AG. In particular, this variant chooses its parameters as in [6, Corollary 2] with $L_{\Psi}$ replaced by $M$, i.e., $\beta_{k}=1 /(2 M)$ for every $k$ (implying a more aggressive stepsize policy). It is worth mentioning that we tested the more conservative AG variant with $\beta_{k}=1 /\left(2 L_{\Psi}\right)$ and observed that it performed substantially less efficient than the above aggressive variant. The fourth method is a special implementation of the adaptive A-ECG method in [7] with $\left(\gamma_{1}, \gamma_{2}, \gamma_{3}\right)=(0.4,0.4,1.0)$ and $(\delta, \sigma)=\left(10^{-2}, 10^{-10}\right)$, which we abbreviate as UP. More specifically, we consider the UPFAG-fullBB method described in [7, Section 4], which uses a Barzilai-Borwein type stepsize selection strategy. The last is the A-ECG method of [13], named NC-FISTA, with $(\xi, \lambda)=(1.05 m, 0.99 / M)$, which we abbreviate as NCF.

Finally, we state some additional details about the numerical experiments. First, the problems considered are of the form in (1) and satisfy assumptions (A1)-(A4) with $f_{2}=f_{2}^{\mathcal{V}} \circ \sigma$ and $h=h^{\mathcal{V}} \circ \sigma$. Second, given a tolerance $\hat{\rho}>0$ and an initial point $Y_{0} \in \operatorname{dom} h$, every method in this section seeks a pair $(\hat{Y}, \hat{V}) \in \operatorname{dom} h \times \mathbb{R}^{m \times n}$ satisfying

$$
\hat{V} \in \nabla f_{1}(\hat{Y})+\nabla\left(f_{2}^{\mathcal{V}} \circ \sigma\right)(\hat{Y})+\partial\left(h^{\mathcal{V}} \circ \sigma\right)(\hat{Y}), \quad \frac{\|\hat{V}\|}{\left\|\nabla f_{1}\left(Y_{0}\right)+\left(f_{2}^{\mathcal{V}} \circ \sigma\right)\left(Y_{0}\right)\right\|+1} \leq \hat{\rho}
$$

and stops after 1000 seconds if such a point cannot be found. Third, to be concise, we abbreviate the IA-ICG and DA-ICG methods as IA and DA, respectively. Finally,

Table 1 Description of the MC data matrices $A \in \mathbb{R}^{m \times n}$

| Name | $\ell$ | $n$ | $\%$ nonzero | $\min _{i, j} A_{i j}$ | $\max _{i, j} A_{i j}$ |
| :--- | :--- | :--- | :--- | :--- | :--- |
| Anime $^{2}$ | 506 | 9437 | $10.50 \%$ | 1 | 10 |
| FilmTrust $^{3}$ | 1508 | 2071 | $1.14 \%$ | 0.5 | 8 |

${ }^{2}$ See the subset of the ratings from https://www.kaggle.com/Coope rUnion/anime-recommendations-database where each user has rated at least 720 items.
${ }^{3}$ See the ratings in the file "ratings.txt" under the FilmTrust section in https://www.librec.net/datasets.html
all described algorithms are implemented in MATLAB 2020a and are run on Linux 64-bit machines that contain at least 8 GB of memory.

### 5.2 Spectral composite problems

This subsection presents computational results of a set of spectral composite optimization problems and contains two sub-subsections. The first one examines a class of nonconvex matrix completion problems, while the second one examines a class of blockwise matrix completion problems.

### 5.2.1 Matrix completion

Given a quadruple $(\alpha, \beta, \mu, \theta) \in \mathbb{R}_{++}^{4}$, a data matrix $A \in \mathbb{R}^{\ell \times n}$, and indices $\Omega$, this subsection considers the following constrained matrix completion (MC) problem:

$$
\begin{aligned}
& \min _{U \in \mathbb{R}^{m \times n}} \frac{1}{2}\left\|P_{\Omega}(U-A)\right\|_{F}^{2}+\kappa_{\mu} \circ \sigma(U)+\tau_{\alpha} \circ \sigma(U) \\
& \text { s.t. }\|U\|_{F}^{2} \leq \sqrt{\ell n} \cdot \max _{i, j}\left|A_{i j}\right|,
\end{aligned}
$$

where $P_{\Omega}$ is the linear operator that zeros out any entry that is not in $\Omega$ and

$$
\kappa_{\mu}(z)=\frac{\mu \beta}{\theta} \sum_{i=1}^{n} \log \left(1+\frac{\left|z_{i}\right|}{\theta}\right), \quad \tau_{\alpha}(z)=\alpha \beta\left[1-\exp \left(-\frac{\|z\|_{2}^{2}}{2 \theta}\right)\right]
$$

for every $z \in \mathbb{R}^{n}$. Here, the function $\kappa_{\mu}+\tau_{\alpha}$ is a nonconvex generalization of the convex elastic net regularizer (see, for example, [17]), and it is well-known (see, for example, [20]) that the function $\kappa_{\mu}-\mu\|\cdot\|_{*}$ is concave, differentiable, and has a $\left(2 \beta \mu / \theta^{2}\right)$-Lipschitz continuous gradient.

We now describe the different data matrices that are considered. Each matrix $A \in \mathbb{R}^{e \times n}$ is obtained from a different collaborative filtering system where each row represents a unique user, each column represents a unique item, and each entry represents a particular rating. Table 1 lists the names of each data set, where the data originates from (in the footnotes), and some basic statistics about the matrices.

We now describe the experiment parameters considered. First the starting point $Z_{0}$ is randomly generated from a shifted binomial distribution that closely follows


Fig. 1 Function values and minimum subgradients for the Anime dataset with $\theta=10^{-1}$
the data matrix $A$. More specifically, the entries of $Z_{0}$ are distributed according to a $\operatorname{BinOmiaL}(a, \mu / a)-\underline{A}$ distribution, where $\mu$ is the sample average of the nonzero entries in $A$, the integer $a$ is the ceiling of the range of ratings in $A$, and $\underline{A}$ is the minimum rating in $A$. Second, the decomposition of the objective function is as follows

$$
\begin{equation*}
f_{1}=\frac{1}{2}\left\|P_{\Omega}(\cdot-A)\right\|_{F}^{2}, \quad f_{2}^{\mathcal{V}}=\mu\left[\kappa_{\mu}(\cdot)-\frac{\beta}{\theta}\|\cdot\|_{1}\right]+\tau_{\alpha}(\cdot), \quad h^{\mathcal{V}}=\frac{\mu \beta}{\theta}\|\cdot\|_{1}+\delta_{\mathcal{F}}(\cdot) \tag{44}
\end{equation*}
$$

where $\mathcal{F}=\left\{U \in \mathbb{R}^{m \times n}:\|U\|_{F} \leq \sqrt{\ell n} \cdot \max _{i, j}\left|A_{i j}\right|\right\}$ is the set of feasible solutions. Third, in view of the previous decomposition, the curvature parameters are set to be

$$
\begin{equation*}
m_{1}=0, \quad M_{1}=1, \quad m_{2}=\frac{2 \beta \mu}{\theta^{2}}+\frac{2 \alpha \beta}{\theta} \exp \left(\frac{-3 \theta}{2}\right), \quad M_{2}=\frac{\alpha \beta}{\theta} \tag{45}
\end{equation*}
$$

where it can be shown that the smallest and largest eigenvalues of $\nabla^{2} \tau_{\alpha}(z)$ are bounded below and above by $-2 \alpha \beta \exp (-3 \theta / 2) / \theta$ and $\alpha \beta / \theta$, respectively, for every $z \in \mathbb{R}^{n}$. Finally, each problem instance uses a specific data matrix $A$ from Table 1, the hyperparameters $(\alpha, \beta, \mu)=(10,20,2)$ and $\hat{\rho}=10^{-6}$, different values of the parameter $\theta$, and $\Omega$ to be the index set of nonzero entries in the chosen matrix $A$.

We now present the results. Figure 1 presents two subplots for the results of the Anime dataset under a value of $\theta=10^{-1}$. The first subplot contains the $\log$ objective value against runtime, while the second one contains the log of the minimal subgradients, i.e. $\min _{i \leq k}\left\|\hat{V}_{i}\right\|$, against runtime. Tables 2 to 3 present the minimal subgradient size obtained within the time limit of 1000 . Moreover, each row of these tables corresponds to a different choice of $\theta$ and the bolded numbers highlight which algorithm performed the best in terms of the size obtained in a run.

### 5.2.2 Blockwise matrix completion

Given a quadruple $(\alpha, \beta, \mu, \theta) \in \mathbb{R}_{++}^{4}$, a block decomposable data matrix $A \in \mathbb{R}^{\ell \times n}$ with blocks $\left\{A_{i}\right\}_{i=1}^{k} \subseteq \mathbb{R}^{p \times q}$, and indices $\Omega$, this subsection considers the following constrained blockwise matrix completion (BMC) problem:

Table 2 Minimum subgradient sizes for the Anime dataset. Times are in seconds and "-" indicates a run that did not generate a subgradient within the given time limit

| Parameters$(\theta, m, M)^{T}$ | Time <br> $t$ | Minimum Subgradient Size $\left(\min _{i \leq k}\left\\|\hat{V}_{i}\right\\|\right)$ |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | ECG | AIPP | AG | UP | NCF | IA | DA |
| $\left[\begin{array}{c}1 \\ 169 \\ 201\end{array}\right]$ | 100 | 1088.6 | 1421.0 | 1568.9 | 1599.4 | 1488.8 | 13.0 | 78.5 |
|  | 200 | 1088.6 | 221.9 | 1510.2 | 132.6 | 1362.4 | 11.6 | 39.2 |
|  | 400 | 1088.6 | 55.6 | 1284.6 | 7.5 | 1147.9 | 11.6 | 11.1 |
|  | 800 | 1088.6 | 7.7 | 716.7 | 7.5 | 862.7 | 11.6 | 11.1 |
| $\left[\begin{array}{c}0.1 \\ 11443 \\ 2001\end{array}\right]$ | 100 | 1542.0 | 1595.8 | 1593.8 | - | 1595.0 | 189.7 | 1345.1 |
|  | 200 | 1489.9 | 1595.0 | 1591.5 | 1595.2 | 1594.2 | 23.1 | 378.1 |
|  | 400 | 1391.0 | 1587.8 | 1584.3 | 1595.1 | 1592.3 | 13.0 | 60.6 |
|  | 800 | 1276.3 | 990.5 | 1557.2 | 1594.3 | 1589.2 | 13.0 | 13.0 |
| $\left[\begin{array}{c}0.01 \\ 839400 \\ 20001\end{array}\right]$ | 100 | 1594.6 | 1595.9 | 1595.6 | 1595.8 | 1595.9 | 162.9 | 452.0 |
|  | 200 | 1592.8 | 1595.6 | 1595.0 | 1595.8 | 1595.8 | 33.5 | 68.0 |
|  | 400 | 1589.8 | 1569.5 | 1592.2 | 1595.8 | 1595.8 | 15.3 | 14.7 |
|  | 800 | 1583.8 | 861.8 | 1582.3 | 1595.7 | 1595.7 | 15.3 | 14.1 |

Table 3 Minimum subgradient sizes for the FilmTrust dataset. Times are in seconds and "-" indicates a run that did not generate a subgradient within the given time limit

| Parameters$(\theta, m, M)^{T}$ | Time | Minimum Subgradient Size $\left(\min _{i \leq k}\left\\|\hat{V}_{i}\right\\|\right)$ |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | ECG | AIPP | AG | UP | NCF | IA | DA |
| $\left[\begin{array}{c}1 \\ 169 \\ 201\end{array}\right]$ | 100 | 127.1 | 328.6 | 328.5 | - | 326.2 | 77.7 | 342.4 |
|  | 200 | 106.7 | 326.2 | 326.8 | 330.0 | 319.4 | 60.7 | 203.1 |
|  | 400 | 106.7 | 294.6 | 319.2 | 330.0 | 305.9 | 60.7 | 186.4 |
|  | 800 | 106.7 | 107.4 | 291.0 | 251.9 | 280.5 | 60.7 | 186.4 |
| $\left[\begin{array}{c}0.1 \\ 11443 \\ 2001\end{array}\right]$ | 100 | 309.0 | 330.0 | 329.6 | 329.9 | 329.9 | 71.0 | 242.3 |
|  | 200 | 287.0 | 326.9 | 327.8 | 329.9 | 329.5 | 71.0 | 235.4 |
|  | 400 | 248.0 | 188.7 | 321.9 | 329.8 | 328.8 | 71.0 | 202.7 |
|  | 800 | 186.9 | 188.7 | 301.8 | 329.4 | 327.4 | 71.0 | 202.7 |
| $\left[\begin{array}{c} 0.01 \\ 839400 \\ 20001 \end{array}\right]$ | 100 | 330.1 | 330.2 | 330.2 | - | 330.2 | 91.8 | 263.9 |
|  | 200 | 330.0 | 330.2 | 330.2 | 330.2 | 330.2 | 91.8 | 262.1 |
|  | 400 | 329.7 | 330.2 | 330.1 | 330.2 | 330.2 | 91.8 | 262.1 |
|  | 800 | 329.2 | 328.7 | 329.7 | 330.2 | 330.2 | 91.8 | 262.1 |



Fig. 2 Function values and minimum subgradients for the truncated normal dataset with $\theta=10^{-1}$

$$
\begin{aligned}
& \min _{U \in \mathbb{R}^{m \times n}} \frac{1}{2}\left\|P_{\Omega}(U-A)\right\|_{F}^{2}+\sum_{i=1}^{k}\left[\kappa_{\mu} \circ \sigma\left(U_{i}\right)+\tau_{\alpha} \circ \sigma\left(U_{i}\right)\right] \\
& \text { s.t. }\|U\|_{F}^{2} \leq \sqrt{\ell n} \cdot \max _{i, j}\left|A_{i j}\right|
\end{aligned}
$$

where $P_{\Omega}, \kappa_{\mu}$, and $\tau_{\alpha}$ are as in Sect. 5.2.1 and $U_{i} \in \mathbb{R}^{p \times q}$ is the $i^{\text {th }}$ block of $U$ with the same indices as $A_{i}$ with respect to $A$.

We now describe the two classes of data matrices that are considered. Every data matrix is a 5 -by- 5 block matrix consisting of 50 -by-100 sized submatrices. Every submatrix contains only $25 \%$ nonzero entries and each data matrix generates its submatrix entries from different probability distributions. More specifically, for a sampled probability $p \sim \operatorname{UNIFORM}[0,1]$ specific to a fixed submatrix, one class uses a $\operatorname{BinOMIAL}(n, p)$ distribution with $n=10$, while the other uses a $\operatorname{TrunCATEDNORMAL}(\mu, \sigma)$ distribution with $\mu=10 p, \sigma^{2}=10 p(1-p)$, and upper and lower bounds 0 and 10 , respectively.

We now describe the experiment parameters considered. First, the the decomposition of the objective function and the quantities $Z_{0},\left(m_{1}, M_{1}\right),\left(m_{2}, M_{2}\right), \hat{\rho}$, and $\Omega$ are the same as in Sect. 5.2.1. Second, we fix $(\alpha, \beta, \mu)=(10,20,2)$ and vary $(\theta, A)$ across the different problem instances.

We now present the results. Figure 2 contains the plots of the log objective function value against the runtime for the binomial data set, listed in increasing order of $M_{2}$. The corresponding plots for the truncated normal data set are similar to the binomial plots so we omit them for the sake of brevity. Tables 4 and 5 present the minimal subgradient size obtained within the time limit of 1000 . Moreover, each row of these tables corresponds to a different choice of $\theta$ and the bolded numbers highlight which algorithm performed the best in terms of the size obtained in a run.

### 5.3 General comments

This subsection makes two comments about the results obtained in the previous subsection. First, within the alloted time (i.e., 1000 seconds), the DA-ICG and IA-ICG

Table 4 Minimum subgradient sizes for the binomial dataset. Times are in seconds and "-" indicates a run that did not generate a subgradient within the given time limit

| Parameters$(\theta, m, M)^{T}$ | Time | Minimum Subgradient Size $\left(\min _{i \leq k}\left\\|\hat{V}_{i}\right\\|\right)$ |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | ECG | AIPP | AG | UP | NCF | IA | DA |
| $\left[\begin{array}{c} 1 \\ 169 \\ 201 \end{array}\right]$ | 100 | 392.4 | 500.3 | 501.2 | 506.0 | 482.5 | 33.9 | 75.5 |
|  | 200 | 392.4 | 478.4 | 492.3 | 506.0 | 465.0 | 33.9 | 43.2 |
|  | 400 | 392.4 | 182.2 | 455.9 | 57.1 | 407.0 | 33.9 | 43.2 |
|  | 800 | 392.4 | 36.7 | 320.6 | 57.1 | 284.3 | 33.9 | 43.2 |
| $\left[\begin{array}{c} 0.1 \\ 11443 \\ 2001 \end{array}\right]$ | 100 | 489.1 | 505.9 | 505.7 | - | 505.8 | 43.4 | 416.0 |
|  | 200 | 476.9 | 505.6 | 505.3 | 505.5 | 505.5 | 43.4 | 76.9 |
|  | 400 | 449.5 | 503.4 | 503.1 | 505.5 | 505.0 | 43.4 | 53.8 |
|  | 800 | 399.4 | 240.8 | 496.2 | 505.3 | 503.9 | 43.4 | 53.8 |
| $\left[\begin{array}{c} 0.01 \\ 839400 \\ 20001 \end{array}\right]$ | 100 | 505.6 | 505.9 | 505.8 | 505.9 | 505.9 | 48.6 | 137.5 |
|  | 200 | 505.1 | 505.9 | 505.7 | 505.9 | 505.9 | 48.6 | 58.6 |
|  | 400 | 504.1 | 498.1 | 504.9 | 505.9 | 505.9 | 48.6 | 58.6 |
|  | 800 | 502.2 | 176.9 | 502.1 | 505.9 | 505.9 | 48.6 | 58.6 |

Table 5 Minimum subgradient sizes for the truncated normal dataset. Times are in seconds and "-" indicates a run that did not generate a subgradient within the given time limit

| Parameters$(\theta, m, M)^{T}$ | Time <br> $t$ | Minimum Subgradient Size $\left(\min _{i \leq k}\left\\|\hat{V}_{i}\right\\|\right)$ |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | ECG | AIPP | AG | UP | NCF | IA | DA |
| $\left[\begin{array}{c}1 \\ 169 \\ 201\end{array}\right]$ | 100 | - | 564.3 | 562.7 | - | 552.2 | 39.1 | 362.3 |
|  | 200 | 433.5 | 551.8 | 554.1 | 566.6 | 536.2 | 30.0 | 80.3 |
|  | 400 | 433.5 | 351.5 | 526.6 | 566.6 | 501.7 | 30.0 | 40.8 |
|  | 800 | 433.5 | 35.6 | 433.7 | 55.8 | 435.7 | 30.0 | 40.8 |
| $\left[\begin{array}{c}0.1 \\ 11443 \\ 2001\end{array}\right]$ | 100 | 533.8 | 566.4 | 566.2 | - | 566.2 | 41.0 | 465.0 |
|  | 200 | 507.4 | 566.1 | 565.7 | 566.0 | 566.0 | 41.0 | 81.4 |
|  | 400 | 478.2 | 563.6 | 561.8 | 566.0 | 565.6 | 41.0 | 50.0 |
|  | 800 | 417.6 | 159.0 | 549.9 | 565.8 | 564.4 | 41.0 | 50.0 |
| $\left[\begin{array}{c} 0.01 \\ 839400 \\ 20001 \end{array}\right]$ | 100 | 565.5 | 566.4 | 566.2 | 566.4 | 566.4 | 45.8 | 54.3 |
|  | 200 | 564.6 | 563.9 | 565.5 | 566.4 | 566.4 | 45.8 | 54.3 |
|  | 400 | 562.7 | 186.1 | 563.1 | 566.3 | 566.4 | 45.8 | 54.3 |
|  | 800 | 559.1 | 143.6 | 555.6 | 566.3 | 566.3 | 45.8 | 54.3 |

methods obtained approximate solutions with small primal residual $\left\|\hat{V}_{k}\right\|$ much faster than the other first-order methods. More specifically, the former methods were able to obtain higher quality solutions much sooner than the latter ones, i.e, within the first 100 seconds. Second, the larger the ratio $m / M$ is, the more efficient the ICG methods are compared to the other benchmarked methods.

## 6 Static ICG iteration complexities

This section establishes the iteration complexities for each of the static ICG methods in Sect. 3.

### 6.1 Static IA-ICG iteration complexity

This subsection establishes the key properties of the static IA-ICG method.
Lemma 8 Let $\left\{\left(y_{i}, \hat{y}_{i}, \hat{v}_{i}\right)\right\}_{i=1}^{k}$ be the collection of iterates generated by the static IAICG method. For every $i \geq 1$, we have

$$
\begin{equation*}
\frac{1}{4 \lambda}\left\|y_{i-1}-y_{i}\right\|^{2} \leq \phi\left(y_{i-1}\right)-\tilde{\ell}_{\phi}\left(y_{i} ; y_{i-1}\right)-\frac{M_{1}}{2}\left\|y_{i}-y_{i-1}\right\|^{2} \leq \phi\left(y_{i-1}\right)-\phi\left(y_{i}\right), \tag{46}
\end{equation*}
$$

where $\tilde{\ell}_{\phi}$ is as in (29).
Proof Let $i \geq 1$ be fixed and let $\left(y_{i}, v_{i}, \varepsilon_{i}\right)$ be the point output by the $i^{\text {th }}$ successful call to the R-ACG algorithm. Moreover, let $\Delta_{1}(\cdot ; \cdot, \cdot)$ be as in (10) with ( $\psi_{s}, \psi_{n}$ ) given by (28). Using the definition of $\tilde{\ell}_{\phi}$, step 2 of the method, and fact that $\left(y^{a}, v, \varepsilon\right)=\left(y_{i}, v_{i}, \varepsilon_{i}\right)$ solves Problem $\mathcal{B}$ in Sect. 2 with $\left(\mu, \psi_{s}, \psi_{n}\right)$ as in (28), we have that

$$
\varepsilon_{i} \geq \Delta_{1}\left(y_{i-1} ; y_{i}, v_{i}\right)=\lambda \tilde{e}_{\phi}\left(y_{i} ; y_{i-1}\right)-\lambda \phi\left(y_{i-1}\right)-\left\langle v_{i}, y_{i}-y_{i-1}\right\rangle+\left\|y_{i}-y_{i-1}\right\|^{2}
$$

Rearranging the above inequality and using assumption (A2), (27), and the fact that $\langle a, b\rangle \geq-\|a\|^{2} / 2-\|b\|^{2} / 2$ for every $a, b \in \mathcal{Z}$ yields

$$
\begin{align*}
& \lambda \phi\left(y_{i-1}\right)-\lambda \tilde{\ell}_{\phi}\left(y_{i} ; y_{i-1}\right) \geq\left\langle v_{i}, y_{i-1}-y_{i}\right\rangle-\varepsilon_{i}+\left\|y_{i}-y_{i-1}\right\|^{2} \\
& \quad=\frac{1}{2}\left\|y_{i}-y_{i-1}\right\|^{2}-\frac{1}{2}\left(\left\|v_{i}\right\|^{2}+2 \varepsilon_{i}\right) \geq\left(\frac{1-\theta^{2}}{2}\right)\left\|y_{i}-y_{i-1}\right\|^{2} \\
& \quad=\frac{\lambda M_{1}}{2}\left\|y_{i}-y_{i-1}\right\|^{2}+\left(\frac{1-\lambda M_{1}-\theta^{2}}{2}\right)\left\|y_{i}-y_{i-1}\right\|^{2}  \tag{47}\\
& \quad=\frac{\lambda M_{1}}{2}\left\|y_{i}-y_{i-1}\right\|^{2}+\frac{1}{4}\left\|y_{i}-y_{i-1}\right\|^{2} .
\end{align*}
$$

Rearranging terms yields the first inequality of (46). The second inequality of (46) follows from the first inequality, the fact that $\tilde{\mathcal{L}}_{\phi}\left(y_{i} ; y_{i-1}\right)+M_{1}\left\|y_{i}-y_{i-1}\right\|^{2} / 2 \geq \phi\left(y_{i}\right)$ from assumption (A2), and the definition of $\tilde{\ell}_{\phi}$.

The next results establish the rate at which the residual $\left\|\hat{v}_{i}\right\|$ tends to 0 .
Lemma 9 Let $p>1$ be given. Then, for every $a, b \in \mathbb{R}^{k}$, we have

$$
\min _{1 \leq i \leq k}\left\{\left|a_{i} b_{i}\right|\right\} \leq k^{-p}\|a\|_{1}\|b\|_{1 /(p-1)}
$$

Proof Let $p>1$ and $a, b \in \mathbb{R}^{k}$ be fixed and let $q \geq 1$ be such that $p^{-1}+q^{-1}=1$. Using the fact that $\langle x, y\rangle \leq\|x\|_{p}\|y\|_{q}$ for every $x, y \in \mathbb{R}^{k}$, and denoting $\tilde{a}$ and $\tilde{b}$ to be vectors with entries $\left|a_{i}\right|^{1 / p}$ and $\left|b_{i}\right|^{1 / p}$, respectively, we have that

$$
\begin{aligned}
k \min _{1 \leq i \leq k}\left\{\left|a_{i} b_{i}\right|\right\}^{1 / p} & \leq \sum_{i=1}^{k}\left|a_{i} b_{i}\right|^{1 / p} \\
& \leq\|\tilde{a}\|_{p}\|\tilde{b}\|_{q}=\|a\|_{1}^{1 / p}\left(\sum_{i=1}^{k}\left|b_{i}\right|^{q / p}\right)^{1 / q}=\left(\|a\|_{1}\|b\|_{q / p}\right)^{1 / p}
\end{aligned}
$$

Dividing by $k$, taking the $p^{\text {th }}$ power on both sides, and using the fact that $p / q=p-1$, yields

$$
\min _{1 \leq i \leq k}\left\{\left|a_{i} b_{i}\right|\right\} \leq k^{-p}\|a\|_{1}\|b\|_{q / p}=k^{-p}\|a\|_{1}\|b\|_{1 /(p-1)}
$$

Proposition $10 \operatorname{Let}\left\{\left(y_{i}, \hat{y}_{i}, \hat{v}_{i}\right)\right\}_{i=1}^{k}$ be as in Lemma 8 and define the quantities

$$
\begin{align*}
L_{1, k}^{\mathrm{avg}} & :=\frac{1}{k} \sum_{i=1}^{k} L_{1}\left(y_{i}, y_{i-1}\right), \quad C_{\lambda, k}^{\mathrm{avg}}:=\frac{1}{k} \sum_{i=1}^{k} C_{\lambda}\left(\hat{y}_{i}, y_{i}\right) \\
D_{k}^{\mathrm{avg}} & :=L_{1, k}^{\mathrm{avg}}+\frac{\theta}{\lambda} C_{\lambda, k}^{\mathrm{avg}}, \quad \beta_{1}:=\left(\frac{1+\bar{C}_{\lambda}}{\lambda}\right)+\sqrt{2}\left(\frac{2+\lambda L_{1}+\theta \bar{C}_{\lambda}}{\lambda}\right) \tag{48}
\end{align*}
$$

where $C_{\lambda}(\cdot, \cdot)$ and $\bar{C}_{\lambda}$ are as in (26) and (29), respectively. Then, we have

$$
\min _{i \leq k}\left\|\hat{v}_{i}\right\|=\mathcal{O}_{1}\left(\left[\sqrt{\lambda} L_{1, k}^{\mathrm{avg}}+\frac{1+\theta C_{\lambda, k}^{\mathrm{avg}}}{\sqrt{\lambda}}\right]\left[\frac{\phi\left(z_{0}\right)-\phi_{*}}{k}\right]^{1 / 2}\right)+\frac{\hat{\rho}}{2}
$$

Proof Using Lemma 3 with $(y, w)=\left(y_{i}, y_{i-1}\right)$ and the fact that $C_{\lambda}(\cdot, \cdot) \leq \bar{C}_{\lambda}$ and $L_{1}(\cdot, \cdot) \leq L_{1}$, we have $\left\|\hat{v}_{i}\right\| \leq \mathcal{E}_{i}\left\|y_{i}-y_{i-1}\right\|$, for every $i \leq k$, where

$$
\mathcal{E}_{i}:=\frac{2+\lambda L_{1}\left(y_{i}, y_{i-1}\right)+\theta C_{\lambda}\left(\hat{y}_{i}, y_{i}\right)}{\lambda} \quad \forall i \geq 1 .
$$

As a consequence, using the sum of the second bound in Lemma 8 from $i=1$ to $k$, the definitions in (48), and Lemma 9 with $p=3 / 2, a_{i}=\mathcal{E}_{i}$, and $b_{i}=\left\|y_{i}-y_{i-1}\right\|$ for $i=1$ to $k$, yields

$$
\begin{align*}
& \min _{i \leq k}\left\|\hat{v}_{i}\right\| \leq \min _{i \leq k} \mathcal{E}_{i}\left\|y_{i}-y_{i-1}\right\| \leq \frac{1}{k^{3 / 2}}\left(\sum_{i=1}^{k} \mathcal{E}_{i}\right)\left(\sum_{i=1}^{k}\left\|y_{i}-y_{i-1}\right\|^{2}\right)^{1 / 2} \\
& \quad=\mathcal{O}_{1}\left(\left[\sqrt{\lambda} L_{1, k}^{\mathrm{avg}}+\frac{1+\theta C_{\lambda, k}^{\mathrm{avg}}}{\sqrt{\lambda}}\right]\left[\frac{\phi\left(z_{0}\right)-\phi_{*}}{k}\right]^{1 / 2}\right) \tag{49}
\end{align*}
$$

We are now ready to give the proof of Theorem 4.
Proof of Theorem 4 (a) This follows from Proposition 10, the fact that $C_{\lambda}(\cdot, \cdot) \leq \bar{C}_{\lambda}$ and $L_{f_{1}}(\cdot, \cdot) \leq L_{1}$, and the stopping condition in step 3 .
(b) The fact that $(\hat{y}, \hat{v})=\left(\hat{y}_{k}, \hat{v}_{k}\right)$ satisfies the inclusion of (22) follows from Lemma 3 with $(y, v, w)=\left(y_{k}, v_{k}, y_{k-1}\right)$. The fact that $\|\hat{v}\| \leq \hat{\rho}$ follows from the stopping condition in step 3 .
(c) This follows from Proposition 2(c) and the fact that method stops in finite number of iterations from part (a).

### 6.2 Static DA-ICG iteration complexity

This subsection establishes several key properties of static DA-ICG method.
To avoid repetition, we assume throughout this subsection that $k \geq 1$ denotes an arbitrary successful outer iteration of the DA-ICG method and let

$$
\left\{\left(a_{i}, A_{i}, y_{i}, y_{i}^{a}, x_{i}, \tilde{x}_{i-1}, \hat{y}_{i}, \hat{v}_{i}, v_{i}, \varepsilon_{i}\right)\right\}_{i=1}^{k}
$$

denote the sequence of all iterates generated by it up to and including the $k^{\text {th }}$ iteration. Observe that this implies that the $i^{\text {th }}$ DA-ICG outer iteration for any $1 \leq i \leq k$ is successful, i.e., the (only) R-ACG call in step 2 of the DA-ICG method does not stop with failure and $\Delta_{1}\left(y_{i-1} ; y_{i}^{a}, v_{i}\right) \leq \varepsilon_{i}$. Moreover, throughout this subsection we let

$$
\begin{equation*}
\widetilde{\gamma}_{i}(u)=\ell_{f_{1}}\left(u ; \tilde{x}_{i-1}\right)+f_{2}(u)+h(u), \quad \gamma_{i}(u)=\widetilde{\gamma}_{i}\left(y_{i}^{a}\right)+\frac{1}{\lambda}\left\langle v_{i}+\tilde{x}_{i-1}-y_{i}^{a}, u-y_{i}^{a}\right\rangle . \tag{50}
\end{equation*}
$$

The first set of results present some basic properties about the functions $\widetilde{\gamma}_{i}$ and $\gamma_{i}$ as well as the iterates generated by the method.

Lemma 11 Let $\Delta_{1}(\cdot ; \cdot, \cdot)$ be as in (10) with $\left(\psi_{s}, \psi_{n}\right)$ given by (28). Then, the following statements hold for any $s \in \operatorname{dom} h$ and $1 \leq i \leq k$ :
(a) $\gamma_{i}\left(y_{i}^{a}\right)=\widetilde{\gamma}_{i}\left(y_{i}^{a}\right)$;
(b) $x_{i}=\operatorname{argmin}_{u \in \Omega}\left\{\lambda a_{i-1} \gamma_{i}(u)+\left\|u-x_{i-1}\right\|^{2} / 2\right\}$;
(c) $y_{i}^{a}-v_{i}=\operatorname{argmin}_{u \in \mathcal{Z}}\left\{\lambda \gamma_{i}(u)+\left\|u-\tilde{x}_{i-1}\right\|^{2} / 2\right\}$;
(d) $-M_{1}\left\|u-\tilde{x}_{i-1}\right\|^{2} / 2 \leq \widetilde{\gamma}_{i}(u)-\phi(u) \leq m_{1}\left\|u-\tilde{x}_{i-1}\right\|^{2} / 2$;
(e) $\phi\left(y_{i-1}\right) \geq \phi\left(y_{i}\right)$ and $\phi\left(y_{i}^{u}\right) \geq \phi\left(y_{i}\right)$.

Proof To keep the notation simple, denote

$$
\begin{align*}
\left(y_{+}^{a}, y_{+}, y, \tilde{x}\right) & =\left(y_{i}^{a}, y_{i}, y_{i-1}, \tilde{x}_{i-1}\right), \quad\left(x_{+}, x\right)=\left(x_{i}, x_{i-1}\right), \\
\left(A_{+}, A, a\right) & =\left(A_{i}, A_{i-1}, a_{i-1}\right), \quad(v, \varepsilon)=\left(v_{i}, \varepsilon_{i}\right) . \tag{51}
\end{align*}
$$

(a) This is immediate from the definitions of $\gamma$ and $\widetilde{\gamma}$ in (50).
(b) Define $\widehat{x}_{i}:=x_{k-1}-a_{k-1}\left(v_{k}+\tilde{x}_{k-1}-y_{k}^{a}\right)$. Using the definition of $\gamma$ in (50), we have that

$$
\begin{aligned}
\underset{u \in \Omega}{\operatorname{argmin}}\left\{\lambda a \gamma(u)+\frac{1}{2}\|u-x\|^{2}\right\} & =\underset{u \in \Omega}{\operatorname{argmin}}\left\{a\left\langle v+\tilde{x}-y_{+}^{a}, u-x\right\rangle+\frac{1}{2}\|u-x\|^{2}\right\} \\
& =\underset{u \in \Omega}{\operatorname{argmin}} \frac{1}{2}\left\|u-\left(x-a\left[v+\tilde{x}-y_{+}^{a}\right]\right)\right\|^{2} \\
& =\underset{u \in \Omega}{\operatorname{argmin}} \frac{1}{2}\left\|u-\hat{x}_{+}\right\|^{2}=x_{+} .
\end{aligned}
$$

(c) Using the definition of $\gamma$ in (50), we have that

$$
\lambda \nabla \gamma\left(y_{+}^{a}-v\right)+\left(y_{+}^{a}-v\right)-\tilde{x}=\left(v+\tilde{x}-y_{+}^{a}\right)+\left(y_{+}^{a}-v\right)-\tilde{x}=0,
$$

and hence, the point $y_{+}^{a}-v$ is the global minimum of $\lambda \gamma+\|\cdot-\tilde{x}\|^{2} / 2$.
(d) This follows from inequality (21) with $i=1$ and the definition of $\widetilde{\gamma}$ in (50).
(e) This follows immediately from the update rule of $y_{i}$ in (36).

Lemma 12 Let $w=\tilde{x}_{i-1}$, the pair $\left(\psi_{n}, \psi_{s}\right)$ be as in (28), and $\Delta_{1}(\cdot ; \cdot, \cdot)$ be as in (10) with $\left(\psi_{s}, \psi_{n}\right)$ given by (28). Then, following statements hold:
(a) the triple $\left(y_{i}^{a}, v_{i}, \varepsilon_{i}\right)$ solves Problem $\mathcal{B}$ and satisfies $\Delta_{1}\left(y_{i-1} ; y_{i}^{a}, v_{i}\right) \leq \varepsilon$, and hence

$$
\begin{equation*}
\left\|v_{i}\right\|+2 \varepsilon_{i} \leq \theta^{2}\left\|y_{i}^{a}-\tilde{x}_{i-1}\right\|^{2}, \quad \Delta_{1}\left(u ; y_{i}^{a}, v_{i}\right) \leq \varepsilon_{i} \quad \forall u \in\left\{\hat{y}_{i}, y_{i-1}\right\} \tag{52}
\end{equation*}
$$

(b) if $f_{2}$ is convex, then $\left(y_{i}^{a}, v_{i}, \varepsilon_{i}\right)$ solves Problem $\mathcal{A}$;
(c) $\Delta_{1}\left(s ; y_{i}^{a}, v_{i}\right)=\lambda\left[\gamma_{i}(s)-\tilde{\gamma}_{i}(s)\right]$;
(d) $\Delta_{1}\left(y_{i} ; y_{i}^{a}, v_{i}\right) \leq \varepsilon$.

Proof (a) This follows from step 2 of the DA-ICG method and Proposition 2(b).
(b) This follows from steps 2 and 3 of the DA-ICG method, the fact that $h$ is convex, and Proposition 2(c) with $\psi_{s}=\widetilde{\gamma}_{i}+\left\|\cdot-\tilde{x}_{i-1}\right\|^{2} / 2$.
(c) Using the definitions of $\left(\psi_{s}, \psi_{n}\right)$ and ( $\gamma, \widetilde{\gamma}$ ) in (28) and (50), respectively, we have that

$$
\begin{aligned}
& \Delta_{1}\left(s ; y_{+}^{a}, v\right)=\left(\psi_{s}+\psi_{n}\right)\left(y_{+}^{a}\right)-\left(\psi_{s}+\psi_{n}\right)(s)-\left\langle v, y_{+}^{a}-s\right\rangle+\frac{1}{2}\left\|s-y_{+}^{a}\right\|^{2} \\
& \quad=\left[\lambda \widetilde{\gamma}\left(y_{+}^{a}\right)+\frac{1}{2}\left\|y_{+}^{a}-\tilde{x}\right\|^{2}\right]-\left[\lambda \widetilde{\gamma}(s)+\frac{1}{2}\|s-\tilde{x}\|^{2}\right]-\left\langle v, y_{+}^{a}-s\right\rangle+\frac{1}{2}\left\|s-y_{+}^{a}\right\|^{2} \\
& \quad=\left[\lambda \gamma(s)+\frac{1}{2}\|s-\tilde{x}\|^{2}\right]-\left[\lambda \widetilde{\gamma}(s)+\frac{1}{2}\|s-\tilde{x}\|^{2}\right]=\lambda \gamma(s)-\lambda \widetilde{\gamma}(s) .
\end{aligned}
$$

(d) If $y_{i}=y_{i-1}$, then this follows from step 3 of the method. On the other hand, if $y_{i}=y_{i}^{a}$, then this follows from part (c).

We now state (without proof) some well-known properties of $A_{i}$ and $a_{i-1}$.
Lemma 13 For every $1 \leq i \leq k$, we have that:
(a) $a_{i-1}^{2}=A_{i}$;
(b) $i^{2} / 4 \leq A_{i} \leq i^{2}$.

The next two lemmas are technical results that are needed to establish the key inequality in Proposition 16.

Lemma 14 For every $u \in \operatorname{dom} h$ and $1 \leq i \leq k$, we have that

$$
\frac{1}{2}\left(A_{i-1}\left\|y_{i-1}-\tilde{x}_{i-1}\right\|^{2}+a_{i-1}\left\|u-\tilde{x}_{i-1}\right\|^{2}\right) \leq 2 D_{\Omega}^{2}+a_{i-1} D_{h}^{2}
$$

Proof Throughout the proof, we use the notation in (51). Using the relation $(p+q)^{2} \leq 2 p^{2}+2 q^{2}$ for every $p, q \in \mathbb{R}$, Lemma 13(a), the fact that $A \leq A^{+}, x \in \Omega$, and $y \in \operatorname{dom} h$, and the definitions of $\tilde{x}$ in (35) and of $D_{\Omega}$ and $D_{h}$ in (37), we conclude that

$$
\begin{aligned}
A \| y & -\tilde{x}\left\|^{2}+a\right\| u-\tilde{x}\left\|^{2}=A\right\| \frac{a}{A_{+}}(y-x)\left\|^{2}+a\right\| \frac{A}{A_{+}}(u-y)+\frac{a}{A_{+}}(u-x) \|^{2} \\
& \leq \frac{A}{A_{+}}\left(\|(y-u)+(u-x)\|^{2}+2 a\left[\frac{A^{2}}{A_{+}^{2}}\|u-y\|^{2}+\frac{a^{2}}{A_{+}^{2}}\|u-x\|^{2}\right]\right) \\
& \leq \frac{2 A}{A^{+}}\left(\|u-y\|^{2}+\|u-x\|^{2}\right)+2 a\|u-y\|^{2}+\frac{2 a}{A_{+}}\|u-x\|^{2} \\
& \leq 2\left[\|u-x\|^{2}+(1+a)\|u-y\|^{2}\right] \leq 2\left[D_{\Omega}^{2}+(1+a) D_{h}^{2}\right] .
\end{aligned}
$$

The conclusion now follows from dividing both sides of the above inequalities by 2 and using the fact that $D_{h} \leq D_{\Omega}$.

Lemma 15 For every $u \in \operatorname{dom} h$ and $1 \leq i \leq k$, we have that

$$
\begin{align*}
& A_{i}\left[\phi\left(y_{i}\right)+\left(\frac{1-\lambda M_{1}}{2 \lambda}\right)\left\|y_{i}^{a}-\tilde{x}_{i-1}\right\|^{2}-\frac{\left\|v_{i}\right\|^{2}}{2 \lambda}\right]+\frac{1}{2 \lambda}\left\|u-x_{i}\right\|^{2}  \tag{53}\\
& \quad \leq A_{i-1} \gamma_{i}\left(y_{i-1}\right)+a_{i-1} \gamma_{i}(u)+\frac{1}{2 \lambda}\left\|u-x_{i-1}\right\|^{2} .
\end{align*}
$$

Proof Throughout the proof, we use the notation in (51). We first present two key expressions. First, using the definition of $\gamma$ in (50) and Lemma 11(c), it follows that

$$
\begin{align*}
& \min _{u \in \mathcal{Z}}\left\{\lambda \gamma(u)+\frac{1}{2}\|u-\tilde{x}\|^{2}\right\}=\lambda \widetilde{\gamma}\left(y_{+}^{a}\right)-\left\langle v+\tilde{x}-y_{+}^{a}, v\right\rangle+\frac{1}{2}\left\|v+\tilde{x}-y_{+}^{a}\right\|^{2} \\
& \quad=\lambda \widetilde{\gamma}\left(y_{+}^{a}\right)-\|v\|^{2}-\left\langle v, \tilde{x}-y_{+}^{a}\right\rangle+\frac{1}{2}\left\|v+\tilde{x}-y_{+}^{a}\right\|^{2}  \tag{54}\\
& \quad=\lambda \widetilde{\gamma}\left(y_{+}^{a}\right)-\frac{1}{2}\|v\|^{2}+\frac{1}{2}\left\|\tilde{x}-y_{+}^{a}\right\|^{2} .
\end{align*}
$$

Second, Lemma 11(b) and the fact that the function $a \gamma+\|\cdot-x\|^{2} /(2 \lambda)$ is $(1 / \lambda)-$ strongly convex imply that

$$
\begin{equation*}
a \gamma\left(x_{+}\right)+\frac{1}{2 \lambda}\left\|x_{+}-x\right\|^{2} \leq a \gamma(u)+\frac{1}{2 \lambda}\|u-x\|^{2}-\frac{1}{2 \lambda}\left\|u-x_{+}\right\|^{2} . \tag{55}
\end{equation*}
$$

Using (54), Lemma 11(d)-(e), Lemma 13(a), and the fact that $\gamma$ is affine, we have that

$$
\begin{align*}
A_{+} & {\left[\phi\left(y_{+}\right)+\left(\frac{1-\lambda M_{1}}{2 \lambda}\right)\left\|y_{+}^{a}-\tilde{x}\right\|^{2}\right] \leq A_{+}\left[\widetilde{\gamma}\left(y_{+}^{a}\right)+\frac{1}{2 \lambda}\left\|y_{+}^{a}-\tilde{x}\right\|^{2}\right] } \\
& =A_{+}\left[\min _{u \in \mathcal{Z}}\left\{\gamma(u)+\frac{1}{2 \lambda}\|u-\tilde{x}\|^{2}\right\}+\frac{\|v\|^{2}}{2 \lambda}\right] \\
& \leq A_{+}\left[\gamma\left(\frac{A y+a x_{+}}{A_{+}}\right)+\frac{1}{2 \lambda}\left\|\frac{A y+a x_{+}}{A_{+}}-\frac{A y+a x}{A_{+}}\right\|^{2}+\frac{\|v\|^{2}}{2 \lambda}\right]  \tag{56}\\
& =A \gamma(y)+a \gamma\left(x_{+}\right)+\frac{a^{2}}{2 \lambda A_{+}}\left\|x-x_{+}\right\|^{2}+\frac{A_{+}}{2 \lambda}\|v\|^{2} \\
& =A \gamma(y)+a \gamma\left(x_{+}\right)+\frac{1}{2 \lambda}\left\|x-x_{+}\right\|^{2}+\frac{A_{+}}{2 \lambda}\|v\|^{2}
\end{align*}
$$

The conclusion now follows from combining (55) with (56).
We now present an inequality that plays an important role in the analysis of the DA-ICG method.

Proposition 16 Let $\Delta_{1}(\cdot ; \cdot, \cdot)$ be as in (10) with $\left(\psi_{s}, \psi_{n}\right)$ as in (28), and define

$$
\begin{equation*}
\theta_{i}(u):=A_{i}\left[\phi\left(y_{i}\right)-\phi(u)\right]+\frac{1}{2 \lambda}\left\|u-x_{i}\right\|^{2} \quad \forall i \geq 0 . \tag{57}
\end{equation*}
$$

For every $u \in \operatorname{dom} h$ satisfying $\Delta_{1}\left(u ; y_{i}^{a}, v_{i}\right) \leq \varepsilon$ and $1 \leq i \leq k$, we have that

$$
\begin{equation*}
\frac{A_{i}}{4 \lambda}\left\|y_{i}^{a}-\tilde{x}_{i-1}\right\|^{2} \leq m_{1}^{+}\left(a_{i-1} D_{h}^{2}+2 D_{\Omega}^{2}\right)+\theta_{i-1}(u)-\theta_{i}(u) \tag{58}
\end{equation*}
$$

Proof Throughout the proof, we use the notation in (51) together with the notation $\pi=\pi_{i-1}$ and $\pi_{+}=\pi_{i}$. Let $u \in \operatorname{dom} h$ be such that $\Delta_{1}\left(u ; y_{+}^{a}, v\right) \leq \varepsilon$. Subtracting $A \phi(u)$ from both sides of the inequality in (53) and using the definition of $\pi_{+}$we have

$$
\begin{align*}
& \frac{A_{+}}{2 \lambda}\left[\left(1-\lambda M_{1}\right)\left\|y_{+}^{a}-\tilde{x}\right\|^{2}-\|v\|^{2}\right]+\pi_{+}(u) \\
& \quad=\frac{A_{+}}{2 \lambda}\left[\left(1-\lambda M_{1}\right)\left\|y_{+}^{a}-\tilde{x}\right\|^{2}-\|v\|^{2}\right]+A_{+}\left[\phi\left(y_{+}\right)-\phi(u)\right]+\frac{1}{2 \lambda}\left\|u-y_{+}^{a}\right\|^{2} \\
& \quad \leq A \gamma(y)+a \gamma(u)-A \phi(u)+\frac{1}{2 \lambda}\|u-x\|^{2} \\
& \quad=a[\gamma(u)-\phi(u)]+A[\gamma(y)-\phi(y)]+\pi(u) . \tag{59}
\end{align*}
$$

Moreover, using Lemma 12(a) and (c), and with our assumption that $\Delta_{1}\left(u ; y_{+}^{a}, v\right) \leq \varepsilon$, we have that

$$
\begin{equation*}
\gamma(s)-\phi(s)=\tilde{\gamma}(s)-\phi(s)+\frac{\Delta_{1}\left(s ; y_{+}^{a}, v\right)}{\lambda} \leq \frac{m_{1}^{+}}{2}\|s-\tilde{x}\|^{2}+\frac{\varepsilon}{\lambda} \quad \forall s \in\{u, y\} . \tag{60}
\end{equation*}
$$

Combining (59), (60), and Lemma 14 then yields

$$
\begin{aligned}
& \frac{A_{+}}{2 \lambda}\left[\left(1-\lambda M_{1}\right)\left\|y_{+}^{a}-\tilde{x}\right\|^{2}-\|v\|^{2}\right]+\pi_{+}(u) \\
& \quad \leq \frac{m_{1}^{+}}{2}\left[a\|u-\tilde{x}\|^{2}+A\|y-\tilde{x}\|^{2}\right]+\frac{\varepsilon A_{+}}{\lambda}+\pi(u) \leq m_{1}^{+}\left(a D_{h}^{2}+2 D_{\Omega}^{2}\right)+\frac{\varepsilon A_{+}}{\lambda}+\pi(u) .
\end{aligned}
$$

Re-arranging the above terms and using (34) together with the first inequality in (52), we conclude that

$$
\begin{aligned}
& m_{1}^{+}\left(a D_{h}^{2}+2 D_{\Omega}^{2}\right)+\pi(u)-\pi_{+}(u) \geq \frac{A_{+}}{2 \lambda}\left[\left(1-\lambda M_{1}\right)\left\|y_{+}^{a}-\tilde{x}\right\|^{2}-\|v\|^{2}-2 \varepsilon\right] \\
& \quad \geq \frac{A_{+}\left(1-\lambda M_{1}-\theta^{2}\right)}{2 \lambda}\left\|y_{+}^{a}-\tilde{x}\right\|^{2} \geq \frac{A_{+}}{4 \lambda}\left\|y_{+}^{a}-\tilde{x}\right\|^{2} .
\end{aligned}
$$

The following result describes some important technical bounds obtained by summing (58) for two different choices of $u$ (possibly changing with $i$ ) from $i=1$ to k.

Proposition 17 Let $\Delta_{\phi}^{0}$ and $d_{0}$ be as in (37) and define

$$
\begin{equation*}
S_{k}:=\frac{1}{4 \lambda} \sum_{i=1}^{k} A_{i}\left\|y_{i}^{a}-\tilde{x}_{i-1}\right\|^{2} \tag{61}
\end{equation*}
$$

Then, the following statements hold:
(a) $S_{k}=\mathcal{O}_{1}\left(k^{2}\left[m_{1}^{+} D_{h}^{2}+\Delta_{\phi}^{0}\right]+k\left[m_{1}^{+}+1 / \lambda\right] D_{\Omega}^{2}\right)$;
(b) if $f_{2}$ is convex, then $S_{k}=\mathcal{O}_{1}\left(k^{2} m_{1}^{+} D_{h}^{2}+k m_{1}^{+} D_{\Omega}^{2}+d_{0}^{2} / \lambda\right)$.

Proof (a) Let $\Delta_{1}(\cdot ; \cdot, \cdot)$ be defined as in (10) with ( $\psi_{s}, \psi_{n}$ ) given by (28). Using (57), the fact that $x_{i}, y_{i}^{a} \in \Omega$, the fact that $A_{i}$ is nonnegative and increasing, and the definitions of $\theta_{i}$ and $D_{\Omega}$ in (57) and (37), respectively, we have that

$$
\begin{align*}
& \sum_{i=1}^{k}\left[\theta_{i-1}\left(y_{i}\right)-\theta_{i}\left(y_{i}\right)\right] \leq \sum_{i=1}^{k} A_{i-1}\left[\phi\left(y_{i-1}\right)-\phi\left(y_{i}\right)\right]+\frac{1}{2 \lambda} \sum_{i=1}^{k}\left\|y_{i}-x_{i-1}\right\|^{2} \\
& \leq A_{k} \sum_{i=1}^{k}\left[\phi\left(y_{i-1}\right)-\phi\left(y_{i}\right)\right]+\frac{k}{2 \lambda} D_{\Omega}^{2} \leq A_{k}\left[\phi\left(y_{0}\right)-\phi_{*}\right]+\frac{k}{2 \lambda} D_{\Omega}^{2} \tag{62}
\end{align*}
$$

Moreover, noting Lemma 12(d) and using Proposition 16 with $u=y_{i}$, we conclude that (58) holds with $u=y_{i}$ for every $1 \leq i \leq k$. Summing these $k$ inequalities and using (62), the definition of $S_{k}$ in (61), and Lemma 13(b) yields the desired conclusion.
(b) Assume now that $f_{2}$ is convex and let $y_{*}$ be a point such that $\phi\left(y_{*}\right)=\phi_{*}$ and $\left\|y_{0}-y_{*}\right\|=d_{0}$. It then follows from Lemma 12(b) and Proposition 1(d) with $(y, v)=\left(y_{i}^{a}, v_{i}\right)$ that $\Delta_{1}\left(y_{*} ; y_{i}^{a}, v_{i}\right) \leq \varepsilon$ for every $1 \leq i \leq k$. The conclusion now follows by using an argument similar to the one in (a) but which instead sums (58) with $u=y_{*}$ from $i=1$ to $k$, and uses the fact that

$$
\sum_{i=1}^{k}\left[\theta_{i-1}\left(y_{*}\right)-\theta_{i}\left(y_{*}\right)\right]=\theta_{0}\left(y_{*}\right)-\theta_{k}\left(y_{*}\right) \leq \frac{1}{2 \lambda}\left\|y_{0}-y_{*}\right\|^{2}=\frac{d_{0}}{2 \lambda},
$$

where the inequality is due to the fact that $\theta_{k}\left(y_{*}\right) \geq 0$ (see (57)) and $A_{0}=0$.
We now establish the rate at which the residual $\left\|\hat{v}_{i}\right\|$ tends to 0 .
Proposition 18 Let $S_{k}$ be as in (61). Moreover, define the quantities

$$
\begin{gather*}
L_{1, k}^{\mathrm{avg}}:=\frac{1}{k} \sum_{i=1}^{k} L_{1}\left(y_{i}^{a}, \tilde{x}_{i-1}\right), \quad C_{\lambda, k}^{\mathrm{avg}}:=\frac{1}{k} \sum_{i=1}^{k} C_{\lambda}\left(\hat{y}_{i}, y_{i}^{a}\right), \\
D_{k}^{\mathrm{erg}}:=L_{1, k}^{\mathrm{erg}}+\frac{\theta}{\lambda} C_{\lambda, k}^{\mathrm{erg}}, \quad 8 \sqrt{2}\left(\frac{2+\lambda L_{1}+\theta \bar{C}_{\lambda}}{\lambda}\right), \tag{63}
\end{gather*}
$$

where $C_{\lambda}(\cdot, \cdot)$ and $\bar{C}_{\lambda}$ are as in (26) and (29), respectively. Then, we have

$$
\min _{i \leq k}\left\|\hat{v}_{i}\right\|=\mathcal{O}_{1}\left(\left[\sqrt{\lambda} L_{1, k}^{\mathrm{avg}}+\frac{1+\theta C_{\lambda, k}^{\mathrm{avg}}}{\sqrt{\lambda}}\right]\left[\frac{S_{k}}{k^{3}}\right]^{1 / 2}\right)+\frac{\hat{\rho}}{2} .
$$

Proof Let $\ell=\lceil k / 2\rceil$. Using Lemma 3 with $(z, w)=\left(y_{i}^{a}, \tilde{x}_{i-1}\right)$ and the bounds $C_{\lambda}(\cdot, \cdot) \leq \bar{C}_{\lambda}$ and $L_{1}(\cdot, \cdot) \leq L_{1}$ we have that $\left\|\hat{v}_{i}\right\| \leq \mathcal{E}_{i}\left\|y_{i}^{a}-\tilde{x}_{i-1}\right\|$, for every $\ell \leq i \leq k$, where

$$
\mathcal{E}_{i}=\frac{2+\lambda L_{1}\left(y_{i}^{a}, \tilde{x}_{i-1}\right)+\theta C_{\lambda}\left(\hat{y}_{i}, y_{i}^{a}\right)}{\lambda} \quad \forall i \geq 1 .
$$

As a consequence, using the definition of $S_{k}$ in (61), the definitions in (63), Lemma 9 with $p=3 / 2, a_{i}=\mathcal{E}_{i} / \sqrt{A_{i}}$, and $b_{i}=\sqrt{A_{i}}\left\|y_{i}^{a}-\tilde{x}_{i-1}\right\|$ for $i \in\{\ell, \ldots, k\}$, Lemma 13(b), and the fact that $(k-\ell+1) \geq k / 2$, yields

$$
\begin{aligned}
& \min _{\ell \leq i \leq k}\left\|\hat{v}_{i}\right\| \leq \min _{\ell \leq i \leq k} \mathcal{E}_{i}\left\|y_{i}^{a}-\tilde{x}_{i-1}\right\| \\
& \quad \leq \frac{1}{(k-\ell+1)^{3 / 2}}\left(\sum_{i=\ell}^{k} \frac{\mathcal{E}_{i}}{\sqrt{A_{i}}}\right)\left(\sum_{i=\ell}^{k} A_{i}\left\|y_{i}^{a}-\tilde{x}_{i-1}\right\|^{2}\right)^{1 / 2} \\
& \quad \leq \frac{2^{3 / 2}}{k^{3 / 2}}\left(\frac{2}{k} \sum_{i=1}^{k} \mathcal{E}_{i}\right)\left(4 \lambda S_{k}\right)^{1 / 2}=\mathcal{O}_{1}\left(\left[\sqrt{\lambda} L_{1, k}^{\mathrm{avg}}+\frac{1+\theta C_{\lambda, k}^{\mathrm{avg}}}{\sqrt{\lambda}}\right]\left[\frac{S_{k}}{k^{3}}\right]^{1 / 2}\right) .
\end{aligned}
$$

We are now ready to prove Theorem 5.
Proof of Theorem 5 (a) This follows from Proposition 18, Proposition 17(a), the fact that $C_{\lambda}(\cdot, \cdot) \leq \bar{C}_{\lambda}$ and $L_{f_{1}}(\cdot, \cdot) \leq L_{1}$, and the termination condition in step 4.
(b) The fact that $(\hat{y}, \hat{v})=\left(\hat{y}_{k}, \hat{r}_{k}\right)$ satisfies the inclusion of (22) follows from Lemma 3 with $\left(y, v, z_{0}\right)=\left(y_{k}^{a}, v_{k}, \tilde{x}_{k-1}\right)$. The fact that $\|\hat{v}\| \leq \hat{\rho}$ follows from the stopping condition in step 4.
(c) The fact that the method does not fail follows from Proposition 2(c). The bound in (39) follows from a similar argument as in part (a) except that Proposition 17(a) is replaced with Proposition 17(b).

## A Technical bounds

The result below presents a basic property of the composite gradient step.
Proposition 19 Let $h \in \overline{\operatorname{Conv}}(\mathcal{Z}), z \in \operatorname{dom} h$, and $g$ be a differentiable function on $\operatorname{dom} h$ which satisfies $g(u)-\ell_{g}(u ; z) \leq L\|u-z\|^{2} / 2$ for some $L \geq 0$ and every $u \in \operatorname{dom} g$. Moreover, define

$$
\hat{z}:=\underset{u}{\operatorname{argmin}}\left\{\ell_{g}(u ; z)+h(u)+\frac{L}{2}\|u-z\|^{2}\right\} .
$$

Then, it holds that

$$
\frac{L}{2}\|z-\hat{z}\|^{2} \leq(g+h)(z)-(g+h)(\hat{z}) .
$$

Proof Using the definition of $\hat{z}$, the fact that $\ell_{g}(\cdot ; z)+h(\cdot)+L\|\cdot-z\|^{2} / 2$ is $L$-strongly convex, and the assumed bound $g(u)-\ell_{g}(u ; z) \leq L\|u-z\|^{2} / 2$ at $u=\hat{z}$, we have

$$
(g+h)(z)=\ell_{g}(z ; z)+h(z) \geq \ell_{g}(\hat{z} ; z)+h(\hat{z})+L\|\hat{z}-z\|^{2} \geq(g+h)(\hat{z})+\frac{L}{2}\|\hat{z}-z\|^{2} .
$$

## B R-ACG algorithm

This section presents technical results related to the R-ACG algorithm.
The first set of results describes some basic properties of the generated iterates.
Proposition 20 If $\psi_{s}$ is $\mu$-strongly convex, then the following statements hold:
(a) $z_{j}^{c}=\operatorname{argmin}_{u \in \mathcal{Z}}\left\{B_{j} \Gamma_{j}(u)+\left\|u-z_{0}^{c}\right\|^{2} / 2\right\}$;
(b) $\Gamma_{j} \leq \psi$ and $B_{j} \psi\left(z_{j}\right) \leq \inf _{u \in \mathcal{Z}}\left\{B_{j} \Gamma_{j}(u)+\left\|u-z_{0}^{c}\right\|^{2} / 2\right\}$;
(c) $\eta_{j} \geq 0$ and $r_{j} \in \partial_{\eta_{j}}\left(\psi-\mu\left\|\cdot-z_{j}\right\|^{2} / 2\right)\left(z_{j}\right)$;
(d) it holds that

$$
\left(\frac{1}{1+\mu B_{j}}\right)\left\|B_{j} r_{j}+z_{j}-z_{0}\right\|^{2}+2 B_{j} \eta_{j} \leq\left\|z_{j}-z_{0}\right\|^{2}
$$

Proof (a) See [14, Proposition 1].
(b) See $[14$, Proposition 1(b)].
(c) The optimality of $z_{j}^{c}$ in part (a), the $\mu$-strong convexity of $\Gamma_{j}$, and the definition of $r_{j}$ imply that

$$
\begin{aligned}
r_{j} & =\frac{z_{0}^{c}-z_{j}^{c}}{B_{j}}+\mu\left(z_{j}-z_{j}^{c}\right) \in \partial\left(\Gamma_{j}-\frac{\mu}{2}\left\|\cdot-z_{j}^{c}\right\|^{2}+\mu\left\langle\cdot, z_{j}^{c}-z_{j}\right\rangle\right)\left(z_{j}^{c}\right) \\
& =\partial\left(\Gamma_{j}-\frac{\mu}{2}\left\|\cdot-z_{j}\right\|^{2}\right)\left(z_{j}^{c}\right) .
\end{aligned}
$$

Using the above inclusion, the definition of $\eta_{j}$, the fact that $\Gamma_{j}-\mu\|\cdot\|^{2} / 2$ is affine, and part (b), we now conclude that

$$
\begin{aligned}
\psi(z)-\frac{\mu}{2}\left\|z-z_{j}\right\|^{2} & \geq \Gamma_{j}(z)-\frac{\mu}{2}\left\|z-z_{j}\right\|^{2}=\Gamma_{j}\left(z_{j}^{c}\right)-\frac{\mu}{2}\left\|z_{j}^{c}-z_{j}\right\|^{2}+\left\langle r_{j}, z-z_{j}^{c}\right\rangle \\
& =\psi\left(z_{j}\right)+\left\langle r_{j}, z-z_{j}\right\rangle-\eta_{j},
\end{aligned}
$$

for every $z \in \operatorname{dom} \psi_{n}$, which is exactly the desired inclusion. The fact that $\eta_{j} \geq 0$ follows from the above inequality with $z=z_{j}$.
(d) It follows from parts (a)-(b) and the definition of $\eta_{j}$ that

$$
\begin{aligned}
\eta_{j} & \leq \Gamma_{j}(u)+\frac{1}{2 B_{j}}\left\|u-z_{0}\right\|^{2}-\psi\left(z_{j}\right) \\
& =\frac{\mu}{2}\left\|z_{j}-z_{j}^{c}\right\|^{2}-\frac{1}{B_{j}}\left\langle z_{0}-z_{j}^{c}, z_{j}-z_{j}^{c}\right\rangle+\frac{1}{2 B_{j}}\left\|z_{j}^{c}-z_{0}\right\|^{2} \\
& =\frac{1}{2 B_{j}}\left\|z_{j}-z_{0}\right\|^{2}-\frac{1}{2 B_{j}}\left(1+\mu B_{j}\right)\left\|z_{j}-z_{j}^{c}\right\|^{2} \\
& =\frac{1}{2 B_{j}}\left\|z_{j}-z_{0}\right\|^{2}-\frac{1}{2 B_{j}\left(1+\mu B_{j}\right)}\left\|B_{j} r_{j}+z_{j}-z_{0}\right\|^{2} .
\end{aligned}
$$

Multiplying both sides of the above inequality by $2 B_{j}$ yields the desired conclusion.

The next result presents the general iteration complexity of the algorithm, i.e. Proposition 2(a).

Proof of Proposition 2(a) Let $\ell$ be the first iteration where

$$
\begin{equation*}
\min \left\{\frac{B_{\ell}^{2}}{4\left(1+\mu B_{\ell}\right)}, \frac{B_{\ell}}{2}\right\} \geq K_{\theta}^{2} \tag{64}
\end{equation*}
$$

and suppose that the R-ACG has not stopped with failure before iteration $\ell$. We show that it must stop with success at the end of the $\ell^{\text {th }}$ iteration. Combining the triangle inequality, the successful check in step 3 of the method, (64), and the relation $(a+b)^{2} \leq 2 a^{2}+2 b^{2}$ for all $a, b \in \mathbb{R}$, we first have that

$$
\begin{aligned}
& \left\|r_{\ell}\right\|^{2}+2 \eta_{\ell} \\
& \quad \leq \max \left\{\frac{1+\mu B_{\ell}}{A_{\ell}^{2}}, \frac{1}{2 B_{\ell}}\right\}\left(\frac{1}{1+\mu B_{\ell}}\left\|B_{\ell} r_{\ell}\right\|^{2}+4 B_{\ell} \eta_{\ell}\right) \\
& \quad \leq \max \left\{\frac{1+\mu B_{\ell}}{B_{\ell}^{2}}, \frac{1}{2 B_{\ell}}\right\}\left(\frac{2}{1+\mu B_{\ell}}\left\|B_{\ell} r_{\ell}+z_{\ell}-z_{0}\right\|^{2}+2\left\|z_{\ell}-z_{0}\right\|^{2}+4 B_{\ell} \eta_{\ell}\right) \\
& \quad \leq \max \left\{\frac{4\left(1+\mu B_{\ell}\right)}{B_{\ell}^{2}}, \frac{2}{B_{\ell}}\right\}\left\|z_{\ell}-z_{0}\right\|^{2} \leq \frac{1}{K_{\theta}^{2}}\left\|z_{\ell}-z_{0}\right\|^{2} \leq \theta^{2}\left\|z_{\ell}-z_{0}\right\|^{2},
\end{aligned}
$$

and hence the method must terminate at the $\ell^{\text {th }}$ iteration. We now bound $\ell$ based on the requirement in (64). Solving for the quadratic in $B_{\ell}$ in the first bound of (64), it is easy to see that $B_{\ell} \geq 4 \mu K_{\theta}^{2}+2 K_{\theta}$ implies (64). On the other hand, for the second
condition in (64), it is immediate that $B_{\ell} \geq 2 K_{\theta}^{2}$ implies (64). In view of (18) and the previous two bounds, it follows that

$$
B_{\ell} \geq \frac{1}{L}\left(1+\sqrt{\frac{\mu}{4 L}}\right)^{2(\ell-1)} \geq 2 K_{\theta}\left(1+2 \mu K_{\theta}^{2}\right)
$$

implies (64). Using the bound $\log (1+t) \geq t /(1+t)$ for $t \geq 0$ and the above bound on $\ell$, it is straightforward to see that $\ell$ is on the same order of magnitude as in (19).

## C Refined ICG points

This appendix presents technical results related to the refined points of the ICG methods.

The result below proves Lemma 3 from the main body of the paper.
Proof of Lemma 3 (a) Using Proposition 1(a), the definition of $\hat{v}$, and the definitions of $\psi_{s}$ and $\psi_{n}$ in (28), we have that

$$
\begin{aligned}
\hat{v} & \in \frac{1}{\lambda}\left[\nabla \psi_{s}(\hat{y})+\partial \psi_{n}(\hat{y})+w-y\right]+\nabla f_{1}(\hat{y})-\nabla f_{1}(w) \\
& =\frac{1}{\lambda}\left[\lambda \nabla f_{1}(w)+\lambda f_{2}(\hat{y})+(w-y)+\lambda \partial h(y)\right]+\nabla f_{1}(\hat{y})-\nabla f_{1}(w) \\
& =\nabla f_{1}(\hat{y})+\nabla f_{2}(\hat{y})+\partial h(\hat{y}),
\end{aligned}
$$

(b) Using assumption (A3), Proposition 1(b), the choice of $M$ in (28), and the fact that $\Delta_{\mu}\left(y_{r} ; y, v\right) \leq \varepsilon$, we first observe that

$$
\begin{align*}
& \left\|\nabla f_{1}(\hat{y})-\nabla f_{1}\left(z_{0}\right)\right\|-L_{1}\left(y, z_{0}\right)\left\|y-z_{0}\right\| \leq L_{1}(y, \hat{y})\|\hat{y}-y\| \\
& \leq \frac{L_{1}(y, \hat{y}) \sqrt{2 \Delta_{\mu}\left(y_{r} ; y, v\right)}}{\sqrt{\lambda M_{2}^{+}+1}} \leq \frac{\theta L_{1}(y, \hat{y})}{\sqrt{\lambda M_{2}^{+}+1}}\left\|y-z_{0}\right\| . \tag{65}
\end{align*}
$$

Using now (65), the choice of $M$ in (28), Proposition 1 (c) with $L(\cdot, \cdot)=\lambda L_{2}(\cdot, \cdot)$, the fact that $\sigma \leq 1$, and the definition of $C_{\lambda}(\cdot, \cdot)$, we conclude that

$$
\begin{aligned}
\|\hat{v}\| & \leq \frac{1}{\lambda}\left\|v_{r}\right\|+\frac{1}{\lambda}\left\|y-z_{0}\right\|+\left\|\nabla f_{1}(\hat{y})-\nabla f_{1}\left(z_{0}\right)\right\| \\
& \leq\left[L_{1}\left(y, z_{0}\right)+\frac{1+\theta}{\lambda}+\frac{\theta\left[\lambda M_{2}^{+}+1+\lambda L_{1}(y, \hat{y})+\lambda L_{2}(y, \hat{y})\right]}{\lambda \sqrt{\lambda M_{2}^{+}+1}}\right]\left\|y-z_{0}\right\| \\
& \leq\left[L_{1}\left(y, z_{0}\right)+\frac{2+\theta C_{\lambda}(y, \hat{y})}{\lambda}\right]\left\|y-z_{0}\right\| .
\end{aligned}
$$

## D Spectral functions

This section presents some results about spectral functions as well as the proof of Propositions 6. It is assumed that the reader is familiar with the key quantities given in Sect. 4.1 (e.g., see (40) and (41)).

We first state two well-known results $[2,11]$ about spectral functions.
Lemma 21 Let $\Psi=\Psi^{\mathcal{\nu}} \circ \sigma$ for some absolutely symmetric function $\Psi^{\mathcal{V}}: \mathbb{R}^{r} \mapsto \mathbb{R}$. Then, the following properties hold:
(a) $\Psi^{*}=\left(\Psi^{\nu} \circ \sigma\right)^{*}=\left(\Psi^{\nu}\right)^{*} \circ \sigma$;
(b) $\nabla \Psi=\left(\nabla \Psi^{\nu}\right) \circ \sigma$;

Lemma 22 Let $\left(\Psi, \Psi^{\nu}\right)$ be as in Lemma 21, the pair $(S, Z) \in \mathcal{Z} \times \operatorname{dom} \Psi$ be fixed, and the decomposition $S=P[\operatorname{dg} \sigma(S)] Q^{*}$ be an $S V D$ of $S$, for some $(P, Q) \in \mathcal{U}^{m} \times \mathcal{U}^{n}$. If $\Psi \in \overline{\operatorname{Conv}} \mathbb{R}^{m \times n}$ and $\Psi^{\mathcal{L}} \in \overline{\mathrm{Conv}} \mathbb{R}^{r}$, then for every $M>0$, we have

$$
S \in \partial\left(\Psi+\frac{M}{2}\|\cdot\|_{F}^{2}\right)(Z) \Longleftrightarrow\left\{\begin{array}{l}
\sigma(S) \in \partial\left(\Psi^{\mathcal{V}}+\frac{M}{2}\|\cdot\|^{2}\right)(\sigma(Z)), \\
Z=P[\operatorname{dg} \sigma(Z)] Q^{*} .
\end{array}\right.
$$

We now present a new result about spectral functions.
Theorem 23 Let $\left(\Psi, \Psi^{\mathcal{V}}\right)$ be as in Lemma 21 and the point $Z \in \mathbb{R}^{m \times n}$ be such that $\sigma(Z) \in \operatorname{dom} \Psi^{\nu}$. Then for every $\varepsilon \geq 0$, we have $S \in \partial_{\varepsilon} \Psi(Z)$ if and only if $\sigma(S) \in \partial_{\varepsilon(S)} \Psi^{\mathcal{L}}(\sigma(Z))$, where

$$
\begin{equation*}
\varepsilon(S):=\varepsilon-[\langle\sigma(Z), \sigma(S)\rangle-\langle Z, S\rangle] \geq 0 \tag{66}
\end{equation*}
$$

Moreover, if $S$ and $Z$ have a simultaneous $S V D$, then $\varepsilon(S)=\varepsilon$.

Proof Using Lemma 21(a), (66), and the well-known fact that $S \in \partial_{\varepsilon} \Psi(Z)$ if and only if $\varepsilon \geq \Psi(Z)+\Psi^{*}(S)-\langle Z, S\rangle$, we have that $S \in \partial_{\varepsilon} \Psi(Z)$ if and only if

$$
\begin{aligned}
\varepsilon(S) & =\varepsilon-[\langle\sigma(Z), \sigma(S)\rangle-\langle Z, S\rangle] \\
& \geq \Psi(Z)+\Psi^{*}(S)-\langle Z, S\rangle-[\langle\sigma(Z), \sigma(S)\rangle-\langle Z, S\rangle] \\
& =\Psi^{\nu}(\sigma(Z))+\left(\Psi^{\nu}\right)^{*}(\sigma(S))-\langle\sigma(Z), \sigma(S)\rangle,
\end{aligned}
$$

or, equivalently, $\sigma(S) \in \partial_{\varepsilon(S)} \Psi^{\mathcal{V}}(\sigma(Z))$ and $\varepsilon(S) \geq 0$. To show that the existence of a simultaneous SVD of $S$ and $Z$ implies $\varepsilon(S)=\varepsilon$ it suffices to show that $\langle\sigma(S), \sigma(Z)\rangle=\langle S, Z\rangle$. Indeed, if $S=P[\operatorname{dg} \sigma(S)] Q^{*}$ and $Z=P[\operatorname{dg} \sigma(Z)] Q^{*}$, for some $(P, Q) \in \mathcal{U}^{m} \times \mathcal{U}^{n}$, then we have

$$
\langle S, Z\rangle=\left\langle\operatorname{dg} \sigma(S), P^{*} P[\operatorname{dg} \sigma(Z)] Q^{*} Q\right\rangle=\langle\operatorname{dg} \sigma(S), \operatorname{dg} \sigma(Z)\rangle=\langle\sigma(S), \sigma(Z)\rangle
$$

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[^1]:    Static IA-ICG Method
    Input: function triple $\left(f_{1}, f_{2}, h\right)$ and scalar quadruple $\left(m_{1}, M_{1}, m_{2}, M_{2}\right) \in \mathbb{R}^{4}$ satisfying (A1)-(A4), tolerance $\hat{\rho}>0$, initial point $y_{0} \in \operatorname{dom} h$, and scalar pair $(\lambda, \theta) \in \mathbb{R}_{++} \times(0,1)$ satisfying

    $$
    \begin{equation*}
    \lambda M_{1}+\theta^{2} \leq \frac{1}{2} \tag{27}
    \end{equation*}
    $$

    Output: a pair $(\hat{y}, \hat{v})$ satisfying (22) or a failure status;
    0 . let $\Delta_{1}(\cdot ; \cdot, \cdot)$ be as in (10) with $\mu=1$, and set $k=1$;

    1. use the R-ACG algorithm to tentatively solve Problem $\mathcal{B}$ associated with (23), i.e., with inputs $\left(\mu, M, \psi_{s}, \psi_{n}\right)$ and $\left(\theta, z_{0}\right)$ where the former is given by

    $$
    \begin{gather*}
    \mu=1, \quad M=\lambda M_{2}^{+}+1 \\
    \psi_{s}=\lambda\left[\ell_{f_{1}}\left(\cdot ; z_{0}\right)+f_{2}\right]+\frac{1}{2}\left\|\cdot-z_{0}\right\|^{2}, \quad \psi_{n}=\lambda h \tag{28}
    \end{gather*}
    $$

    and $z_{0}=y_{k-1}$; if the R-ACG stops with failure, then stop with a failure status; otherwise, let $\left(y_{k}, v_{k}, \varepsilon_{k}\right)$ denote its output and go to step 2 ;
    2. if the inequality $\Delta_{1}\left(y_{k-1} ; y_{k}, v_{k}\right) \leq \varepsilon_{k}$ holds, then go to step 3 ; otherwise, stop with a failure status;
    3. set $\left(\hat{y}_{k}, \hat{v}_{k}\right)=S R P\left(y_{k}, v_{k}, y_{k-1}\right)$; if $\left\|\hat{v}_{k}\right\| \leq \hat{\rho}$ then stop with a success status and output $(\hat{y}, \hat{v})=\left(\hat{y}_{k}, \hat{v}_{k}\right)$; otherwise, update $k \leftarrow k+1$ and go to step 1 .

[^2]:    ${ }^{1}$ See https://github.com/wwkong/nc_opt/tree/master/tests/papers/icg.

