# A NEW MATRIX-FREE ALGORITHM FOR THE LARGE-SCALE TRUST-REGION SUBPROBLEM* 

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#### Abstract

We present a new method for the large-scale trust-region subproblem. The method is matrix-free in the sense that only matrix-vector products are required. We recast the trust-region subproblem as a parameterized eigenvalue problem and compute an optimal value for the parameter. We then find the solution of the trust-region subproblem from the eigenvectors associated with two of the smallest eigenvalues of the parameterized eigenvalue problem corresponding to the optimal parameter. The new algorithm uses a different interpolating scheme than existing methods and introduces a unified iteration that naturally includes the so-called hard case. We show that the new iteration is well defined and convergent at a superlinear rate. We present computational results to illustrate convergence properties and robustness of the method.


Key words. regularization, constrained quadratic optimization, trust region, Lanczos method

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1. Introduction. An important problem in optimization and linear algebra is the trust-region subproblem: minimize a quadratic function subject to an ellipsoidal constraint,

$$
\min \frac{1}{2} x^{T} A x+g^{T} x \quad \text { subject to (s.t.) } \quad\|x\|_{2} \leq \Delta
$$

where $A \in \mathbb{R}^{n \times n}, A=A^{T} ; x, g \in \mathbb{R}^{n}$ and $\Delta>0$. Two significant applications of this basic problem are the regularization or smoothing of discrete forms of ill-posed problems and the trust-region globalization strategy used to force convergence in optimization methods.

A solution $x_{*}$ to the problem must satisfy an equation of the form $(A+\mu I) x_{*}=$ $-g$ with $\mu \geq 0$. The parameter $\mu$ is the Tikhonov regularization parameter for ill-posed problems and the Levenberg-Marquardt parameter in optimization. The constraint might also involve a matrix $C \neq I$, where $C$ is often constructed to impose a smoothness condition on the solution $x_{*}$ for ill-posed problems and to incorporate scaling of the variables in optimization. We will not treat this case explicitly here. However, in many applications the matrix $C$ will be nonsingular and therefore with a change of variables we can reduce the problem to the case we are considering.

[^0]If we can afford to compute the Cholesky factorization of matrices of the form $A+\mu I$, then the method proposed by Moré and Sorensen (cf. [10]) is the method of choice to solve the problem. However, in many important applications, factoring or even forming these matrices is prohibitive. This has motivated the development of matrix-free methods that rely only on matrix-vector products. The first method in this class is the method of Steihaug [18] which computes the solution to the problem in a Krylov subspace. This method is very efficient in conjunction with optimization methods; however, it does not compute an optimal solution and cannot handle a special situation known as the hard case, which we will describe later. New methods based on matrix-vector products are the ones by Golub and von Matt [3], Sorensen [17], Rendl and Wolkowicz [13], and Pham Dinh and Hoai An [12]. Recently, Lucidi, Palagi, and Roma [8] presented new properties of the trust-region subproblem that provide useful tools for the development of new classes of algorithms for this problem in the large-scale context. Recently, we became aware of a new method proposed by Hager [5], where an SQP approach is used to solve the trust-region subproblem.

Golub and von Matt [3] base their algorithm on the theory of Gauss quadrature and do not include in their analysis the possibility of the hard case. Pham Dinh and Hoai An [12] develop an algorithm based on the difference of convex functions. Their strategy is very inexpensive, due to its projective nature, but needs a restarting mechanism to ensure convergence to a global solution. The approaches of Sorensen [17] and Rendl and Wolkowicz [13] recast the trust-region subproblem as a parameterized eigenvalue problem and design an iteration to find an optimal value for the parameter. The idea of formulating the trust-region subproblem in terms of an eigenvalue problem is also exploited in Gander, Golub, and von Matt [1]. Rendl and Wolkowicz present a primal-dual semidefinite framework for the trust-region subproblem, where a dual simplex-type method is used in the basic iteration and a primal simplex-type method provides steps for the hard-case iteration. At each iteration, the method computes the smallest eigenvalue and corresponding eigenvector of the parameterized problem using a block Lanczos routine. Sorensen's algorithm provides a superlinearly convergent scheme to adjust the parameter and finds the optimal vector $x_{*}$ from the eigenvector of the parameterized problem, as long as the hard case does not occur. For the hard case, Sorensen's algorithm is linearly convergent. The algorithm uses the implicitly restarted Lanczos method (IRLM) (cf. [16]) to compute the smallest eigenvalue and corresponding eigenvector of the parameterized problem. The IRLM is particularly suitable for large-scale applications since it has low and fixed storage requirements and relies only on matrix-vector products.

In this work we present a new matrix-free algorithm for the large-scale trust-region subproblem. Our algorithm is similar to those proposed in $[13,17]$ in the sense that we solve the trust-region subproblem through a parameterized eigenvalue problem, but it differs from those approaches in that we do not need two different schemes for the standard case and the hard case. Our algorithm can handle all the cases in the same basic iteration. We achieved this improvement over the methods in $[13,17]$ by computing two eigenpairs of the parameterized problem and incorporating information about the second eigenpair whenever it is appropriate. This strategy does not substantially increase the computational cost over the method in [17]. We introduce a two-point interpolating scheme that is different from the one in [17]. We show that this new iteration is also convergent and that the convergence rate is superlinear. Moreover, our convergence results naturally include the hard case, since no special iterations are necessary. Such a unified approach is not achieved in either [13] or [17].

The organization of this work is the following. In section 2 we analyze the structure of the problem and motivate the algorithm. In section 3 we give a complete characterization of the hard case with respect to the parameterized eigenproblems. We describe the algorithm in detail in section 4 . In section 5 we present the convergence analysis. We describe preliminary numerical experiments in section 6 and present some conclusions in section 7 .
2. Structure of the problem. The problem we are interested in solving is

$$
\begin{array}{cl}
\min & \psi(x) \\
\text { s.t. } & \|x\| \leq \Delta \tag{2.1}
\end{array}
$$

where $\psi(x)=\frac{1}{2} x^{T} A x+g^{T} x ; A, g$ as before and $\|\cdot\| \equiv\|\cdot\|_{2}$ throughout the paper.
Due to the structure of (2.1), its optimality conditions are both necessary and sufficient, as stated in the next lemma, where we follow [17] in the nonstandard but notationally more convenient use of a nonpositive multiplier.

Lemma 2.1 (see [15]). A feasible vector $x_{*}$ is a solution to (2.1) with corresponding Lagrange multiplier $\lambda_{*}$ if and only if $x_{*}, \lambda_{*}$ satisfy $\left(A-\lambda_{*} I\right) x_{*}=-g$ with $A-\lambda_{*}$ I positive semidefinite, $\lambda_{*} \leq 0$, and $\lambda_{*}\left(\Delta-\left\|x_{*}\right\|\right)=0$.

Proof. For the proof see [15]. $\quad$.
In order to design efficient methods for solving problem (2.1) we must exploit the tremendous amount of structure of this problem. In particular, the optimality conditions are computationally attractive since they provide a means to reduce the given $n$-dimensional constrained optimization problem into a zero-finding problem in a single scalar variable. For example, we could define the function $\varphi(\lambda)=\left\|(A-\lambda I)^{-1} g\right\|$ and solve the secular equation $\varphi(\lambda)=\Delta$, monitoring $\lambda$ to be no greater than the smallest eigenvalue of $A$, so that the Cholesky factorization of $A-\lambda I$ is well defined. Using Newton's method to solve $\frac{1}{\varphi(\lambda)}-\frac{1}{\Delta}=0$ has a number of computationally attractive features (cf. [10]) and we should use this approach when we can afford to compute the Cholesky factorization of $A-\lambda I$. When computing a Cholesky factorization is too expensive, we need to use a different strategy. The introduction of a new parameter will make it possible to convert the original trust-region subproblem into a scalar problem that is suitable for the large-scale setting. The conversion amounts to embedding the given problem into a parameterized bordered matrix eigenvalue problem. Consider the bordered matrix

$$
B_{\alpha}=\left(\begin{array}{cc}
\alpha & g^{T} \\
g & A
\end{array}\right)
$$

and observe that

$$
\frac{\alpha}{2}+\psi(x)=\frac{1}{2}\left(1, x^{T}\right) B_{\alpha}\binom{1}{x} .
$$

Therefore, there exists a value of the parameter $\alpha$ such that we can rewrite problem (2.1) as

$$
\begin{array}{cl}
\min & \frac{1}{2} y^{T} B_{\alpha} y \\
\mathrm{s.t.} & y^{T} y \leq 1+\Delta^{2}, e_{1}^{T} y=1 \tag{2.2}
\end{array}
$$

where $e_{1}$ is the first canonical unit vector in $\mathbb{R}^{n+1}$. This formulation suggests that we can find the desired solution in terms of an eigenpair of $B_{\alpha}$ in the following way.

Suppose that $\left\{\lambda,\left(1, x^{T}\right)^{T}\right\}$ is an eigenpair of $B_{\alpha}$. Then

$$
\left(\begin{array}{cc}
\alpha & g^{T} \\
g & A
\end{array}\right)\binom{1}{x}=\binom{1}{x} \lambda
$$

which is equivalent to

$$
\begin{equation*}
\alpha-\lambda=-g^{T} x \tag{2.3}
\end{equation*}
$$

and

$$
\begin{equation*}
(A-\lambda I) x=-g \tag{2.4}
\end{equation*}
$$

Now, let $\delta_{1}, \delta_{2}, \ldots, \delta_{d}$ be the distinct eigenvalues of $A$ in nondecreasing order. Then

$$
\begin{equation*}
\alpha-\lambda=-g^{T} x=\sum_{j=1}^{d} \frac{\beta_{j}^{2}}{\delta_{j}-\lambda}, \tag{2.5}
\end{equation*}
$$

where $\beta_{j}^{2}$ is the sum of the squares of the expansion coefficients of $g$ in the eigenvector basis, corresponding to all the eigenvectors associated with $\delta_{j}$.

Observe that as a consequence of Cauchy's interlace theorem (cf. [11, p. 186]), and also from (2.5), the eigenvalues of $A$ interlace the eigenvalues of $B_{\alpha}$. In particular, if $\lambda_{1}(\alpha)$ is the smallest eigenvalue of $B_{\alpha}$, then $\lambda_{1}(\alpha) \leq \delta_{1}$. This implies that the matrix $A-\lambda_{1}(\alpha) I$ is always positive semidefinite independently of the value of $\alpha$.

Equations (2.3)-(2.4) express $\lambda$ and hence $x$ implicitly in terms of $\alpha$, suggesting the definition of a convenient function as follows. Let ${ }^{\dagger}$ denote the pseudoinverse of a matrix and let us define

$$
\phi(\lambda) \equiv g^{T}(A-\lambda I)^{\dagger} g=-g^{T} x
$$

Therefore,

$$
\phi^{\prime}(\lambda)=g^{T}\left[(A-\lambda I)^{\dagger}\right]^{2} g=x^{T} x
$$

where differentiation is with respect to $\lambda$, and $x$ satisfies $(A-\lambda I) x=-g$. The function $\phi$ appears in many contexts [2, 9, 19, 20] and Figure 1(a) shows its typical behavior. It is worth noticing that the values of $\phi$ and $\phi^{\prime}$ at an eigenvalue $\lambda$ of $B_{\alpha}$ are readily available and contain valuable information with respect to problem (2.1), as long as $\lambda$ has a corresponding eigenvector with nonzero first component.

Finding the smallest eigenvalue and a corresponding eigenvector of $B_{\alpha}$ for a given value of $\alpha$, and then normalizing the eigenvector to have its first component equal to one, will provide a means to evaluate the rational function $\phi$ and its derivative at appropriate values of $\lambda$, namely, at $\lambda=\lambda_{1}(\alpha) \leq \delta_{1}$. Moreover, $\lambda_{1}(\alpha)$ is usually well separated from the rest of the spectrum of $B_{\alpha}$, particularly for small values of $\Delta$. In these cases, we expect a Lanczos-type method to be very efficient in computing this eigenvalue and the corresponding eigenvector. If we can adjust $\alpha$ so that the corresponding $x$ satisfies $x^{T} x=\phi^{\prime}(\lambda)=\Delta^{2}$ with $\alpha-\lambda=\phi(\lambda)$, then

$$
(A-\lambda I) x=-g \quad \text { and } \quad \lambda(\Delta-\|x\|)=0
$$

with $A-\lambda I$ positive semidefinite. If $\lambda \leq 0$, then $x$ is a boundary solution for the trust-region subproblem. In case we find $\lambda>0$ with $\|x\|<\Delta$ during the course of


Fig. 1. Example of the typical pattern of $\phi(\lambda)$ (solid) and the straight line $f(\lambda)=\alpha_{*}-\lambda$ (dash-dotted). The three smallest eigenvalues of $A$ are $-2,-0.5$, and 2. (a) General case with the slope at $\lambda_{*}$ also plotted; (b) exact hard case; (c) near hard case; (d) detail of box in (c).
adjusting $\alpha$, then this implies that the matrix $A$ is positive definite and that $\left\|A^{-1} g\right\|<$ $\Delta$. As shown in [10], these two conditions imply that problem (2.1) has an interior solution that satisfies $A x=-g$.

The availability of the values $\lambda, \phi(\lambda), \phi^{\prime}(\lambda)$ makes it possible to use rational interpolation to adjust the parameter using these values as interpolation points. The adjustment of $\alpha$ by means of rational interpolation consists of constructing a rational interpolant $\widehat{\phi}$ and finding a special point $\widehat{\lambda}$ such that $\widehat{\phi}^{\prime}(\widehat{\lambda})=\Delta^{2}$. We then compute the new parameter as $\alpha_{+}=\widehat{\lambda}+\widehat{\phi}(\widehat{\lambda})$. In this approach it is necessary to safeguard $\alpha_{+}$to ensure convergence of the iteration. This idea was discussed in $[6,15]$ and used in [17]. The algorithm in this paper follows this approach.
3. Characterization of the hard case. We assumed in the previous discussion that the smallest eigenvalue of $B_{\alpha}$ had a corresponding eigenvector with nonzero first component. It remains to consider the possibility that all the eigenvectors associated with $\lambda_{1}(\alpha)$ have first component zero so that we cannot normalize any of them to have its first component equal to one. In this case, the proposed strategy for solving problem (2.1) breaks down. However, this can happen only when $g$ is orthogonal to $\mathcal{S}_{1}$, where $\mathcal{S}_{j}=\left\{q \mid A q=\delta_{j} q\right\}, j=1,2, \ldots, d$.

The condition $g \perp \mathcal{S}_{1}$ is a necessary condition for the occurrence of the so-called hard case. Therefore, we call this situation a potential hard case. Observe that in a potential hard case $\delta_{1}$ is no longer a pole of $\phi$, as Figure 1(b) illustrates. We discuss
the hard case in detail at the end of this section. At this point we will concentrate on the potential hard case, which has intriguing consequences. We will show that in a potential hard case, for all values of $\alpha$ greater than certain critical value $\tilde{\alpha}_{1}$, all the eigenvectors corresponding to the smallest eigenvalue of $B_{\alpha}$ will have first component zero. We will also show that for any $\alpha$, there is always a well-defined eigenvector of $B_{\alpha}$, depending continuously on $\alpha$, that we can safely normalize to have first component one. If $g \not \perp \mathcal{S}_{1}$ or $g \perp \mathcal{S}_{1}$ and $\alpha \leq \tilde{\alpha}_{1}$, then this eigenvector corresponds to $\lambda_{1}(\alpha)$. If $g \perp \mathcal{S}_{1}$ and $\alpha$ exceeds the critical value $\tilde{\alpha}_{1}$ by a small amount, this parameterized vector corresponds to the second smallest eigenvalue of $B_{\alpha}$. A complete understanding of this case leads to the main algorithm of this paper. The following results are the basis for this understanding.

Lemma 3.1. For any $\alpha \in \mathbb{R}$ and any $q \in \mathcal{S}_{j}, 1 \leq j \leq d$, $\left\{\delta_{j},\left(0, q^{T}\right)^{T}\right\}$ is an eigenpair of $B_{\alpha}$ if and only if $g$ is orthogonal to $\mathcal{S}_{j}$.

Proof. The proof follows from the observation that $g \perp \mathcal{S}_{j}$ and $A q=\delta_{j} q$ are equivalent to

$$
\left(\begin{array}{cc}
\alpha & g^{T} \\
g & A
\end{array}\right)\binom{0}{q}=\delta_{j}\binom{0}{q} .
$$

If $\mathcal{Z}_{1}(\alpha)$ is the eigenspace of $B_{\alpha}$ corresponding to $\delta_{1}$, Lemma 3.1 establishes that the set $\left\{\left(0, q^{T}\right)^{T} \mid q \in \mathcal{S}_{1}\right\}$ is a subset of $\mathcal{Z}_{1}(\alpha)$. Note that while $\mathcal{S}_{1}$ corresponds to the smallest eigenvalue of $A, \mathcal{Z}_{1}(\alpha)$ does not necessarily correspond to the smallest eigenvalue of $B_{\alpha}$. These subspaces have the same dimension for all but one exceptional value of $\alpha$. The following result states that there is a unique value of $\alpha$ for which $\operatorname{dim} \mathcal{Z}_{1}(\alpha)=\operatorname{dim} \mathcal{S}_{1}+1$.

Lemma 3.2. Suppose that $g$ is orthogonal to $\mathcal{S}_{j}, 1 \leq j \leq d$, and let $p_{j}=$ $-\left(A-\delta_{j} I\right)^{\dagger} g$. The pair $\left\{\delta_{j},\left(1, p_{j}^{T}\right)^{T}\right\}$ is an eigenpair of $B_{\alpha}$ if and only if $\alpha=\tilde{\alpha}_{j}$, where $\tilde{\alpha}_{j}=\delta_{j}-g^{T} p_{j}$.

Proof. First we observe that $g \perp \mathcal{S}_{j}$ implies that $g \in \mathcal{R}\left(A-\delta_{j} I\right)$ and therefore

$$
\begin{equation*}
\left(A-\delta_{j} I\right) p_{j}=-\left(A-\delta_{j} I\right)\left(A-\delta_{j} I\right)^{\dagger} g=-g, \tag{3.1}
\end{equation*}
$$

since $\left(A-\delta_{j} I\right)\left(A-\delta_{j} I\right)^{\dagger}$ is an orthogonal projector onto $\mathcal{R}\left(A-\delta_{j} I\right)$.
Now, let $\alpha=\tilde{\alpha}_{j}$. Then

$$
\left(\begin{array}{cc}
\tilde{\alpha}_{j} & g^{T} \\
g & A
\end{array}\right)\binom{1}{p_{j}}=\binom{\tilde{\alpha}_{j}+g^{T} p_{j}}{g+A p_{j}}=\delta_{j}\binom{1}{p_{j}},
$$

since by definition of $\tilde{\alpha}_{j}$ we have $\tilde{\alpha}_{j}+g^{T} p_{j}=\delta_{j}$ and by (3.1), $g+A p_{j}=\delta_{j} p_{j}$.
Suppose now that $\left\{\delta_{j},\left(1, p_{j}^{T}\right)^{T}\right\}$ is an eigenpair of $B_{\alpha}$, i.e.,

$$
\left(\begin{array}{cc}
\alpha & g^{T} \\
g & A
\end{array}\right)\binom{1}{p_{j}}=\delta_{j}\binom{1}{p_{j}} .
$$

It follows directly from this relationship that $\alpha=\tilde{\alpha}_{j}=\delta_{j}-g^{T} p_{j}$.
The following corollary summarizes the main results from Lemmas 3.1 and 3.2.
Corollary 3.1. Suppose that $g$ is orthogonal to $\mathcal{S}_{j}, 1 \leq j \leq d$, and let $\mathcal{Z}_{j}(\alpha)=$ $\left\{z \in \mathbb{R}^{n+1} \mid B_{\alpha} z=\delta_{j} z\right\}$. If $\tilde{\alpha}_{j}=\delta_{j}+g^{T} p_{j}$ with $p_{j}=-\left(A-\delta_{j} I\right)^{\dagger} g$, then $\operatorname{dim} \mathcal{Z}_{j}\left(\tilde{\alpha}_{j}\right)=\operatorname{dim} \mathcal{S}_{j}+1$ and for any other value of $\alpha, \operatorname{dim} \mathcal{Z}_{j}(\alpha)=\operatorname{dim} \mathcal{S}_{j}$. Moreover, if $m_{j}$ is the multiplicity of $\delta_{j}$ and $\left\{q_{1}, \ldots, q_{m_{j}}\right\}$ is an orthogonal basis for $\mathcal{S}_{j}$, then

$$
\left\{\binom{1}{p_{j}},\binom{0}{q_{1}}, \ldots,\binom{0}{q_{m_{j}}}\right\}
$$

is an orthogonal basis for $\mathcal{Z}_{j}\left(\tilde{\alpha}_{j}\right)$ and

$$
\left\{\binom{0}{q_{1}}, \ldots,\binom{0}{q_{m_{j}}}\right\}
$$

is an orthogonal basis for $\mathcal{Z}_{j}(\alpha)$, for $\alpha \neq \tilde{\alpha}_{j}$.
The result in Lemma 3.1 was also stated in [17]; the idea behind Lemma 3.2 was presented in [13]. We present here a general formulation of these results given in [14]. In the next results from [14], we establish that there always exists an eigenvector of $B_{\alpha}$ that we can normalize to have first component equal to one, and we characterize the eigenvalue to which this eigenvector corresponds.

ThEOREM 3.1 (see [14]). Let $\lambda(\alpha)$ be the smallest solution of the equation

$$
\phi(\lambda)=\alpha-\lambda
$$

Then, for any value of $\alpha, \lambda(\alpha)$ is an eigenvalue of $B_{\alpha}$ with a corresponding eigenvector that can be normalized to have first component one.

Proof. Suppose first that $g$ is orthogonal to $\mathcal{S}_{i}, i=1,2, \ldots, \ell$, with $1 \leq \ell<d$. Then

$$
\begin{aligned}
\phi(\lambda) & =g^{T}(A-\lambda I)^{\dagger} g \\
& =\sum_{j=\ell+1}^{d} \frac{\beta_{j}^{2}}{\delta_{j}-\lambda}
\end{aligned}
$$

Let $\lambda(\alpha)$ be the smallest solution of the equation $\phi(\lambda)=\alpha-\lambda$. Then $\lambda(\alpha) \in$ $\left(-\infty, \delta_{\ell+1}\right)$. Since $\phi(\lambda)$ is strictly increasing on its domain and $f(\lambda)=\alpha-\lambda$ is a decreasing straight line, we conclude that $\lambda(\alpha)$ is unique. Since $\lambda(\alpha)$ depends continuously on $\alpha$, so does $p(\alpha)=-(A-\lambda(\alpha) I)^{\dagger} g$ and also $v(\alpha)=\left(1, p(\alpha)^{T}\right)^{T}$. Let us see now that $v(\alpha)$ is an eigenvector of $B_{\alpha}$ associated with $\lambda(\alpha)$. Consider

$$
\left(\begin{array}{cc}
\alpha & g^{T} \\
g & A
\end{array}\right)\binom{1}{p(\alpha)}=\binom{\alpha+g^{T} p(\alpha)}{g+A p(\alpha)}
$$

and note that

$$
\begin{aligned}
\alpha+g^{T} p(\alpha) & =\alpha-\phi(\lambda(\alpha)) \\
& =\lambda(\alpha) \quad \text { by definition of } \lambda(\alpha)
\end{aligned}
$$

Now, $g \perp \mathcal{S}_{i}, i=1,2, \ldots, \ell$, implies that $g \in \mathcal{R}(A-\lambda I)$ for $\lambda \in\left(-\infty, \delta_{\ell+1}\right)$. Thus, $g \in \mathcal{R}(A-\lambda(\alpha) I)$ and we have $(A-\lambda(\alpha) I) p(\alpha)=-g$. It follows that

$$
g+A p(\alpha)=\lambda(\alpha) p(\alpha)
$$

and therefore, $B_{\alpha} v(\alpha)=\lambda(\alpha) v(\alpha)$.
Suppose now that $g$ is not orthogonal to $\mathcal{S}_{1}$. Then $\lambda(\alpha) \in\left(-\infty, \delta_{1}\right)$ and this implies $A-\lambda(\alpha) I$ is nonsingular and the previous proof holds with $p(\alpha)=$ $-(A-\lambda(\alpha) I)^{-1} g$.

The following result characterizes the smallest $\ell+1$ distinct eigenvalues of $B_{\alpha}$ if $g$ is orthogonal to the eigenspaces corresponding to the smallest $\ell$ distinct eigenvalues of $A$. In case $g$ is not orthogonal to $\mathcal{S}_{1}$, then the lemma characterizes the smallest eigenvalue of $B_{\alpha}$. We will denote by $\lambda_{j}(\alpha), j=1,2, \ldots, n+1$, the eigenvalues of $B_{\alpha}$ in nondecreasing order.

Lemma 3.3 (see [14]). Let $\{\lambda(\alpha), v(\alpha)\}$ be the eigenpair of $B_{\alpha}$ given by Theorem 3.1 and define $\tilde{\alpha}_{i}$ as in Lemma 3.2. Then if $g \not \not \mathcal{\mathcal { S } _ { 1 }}$, then $\lambda_{1}(\alpha)=\lambda(\alpha)$.

If $g \perp \mathcal{S}_{k}$, for $k=1,2, \ldots, \ell$ and $1 \leq \ell<d$, then
(i) if $\alpha=\tilde{\alpha}_{i}, 1 \leq i \leq \ell$, then $\lambda_{j}(\alpha)=\delta_{j}, j=1,2, \ldots, \ell$. In this case, $\lambda_{\ell+1}(\alpha)$ is the second smallest root of equation $\phi(\lambda)=\alpha-\lambda$;
(ii) if $\alpha<\tilde{\alpha}_{1}$, then $\lambda_{1}(\alpha)=\lambda(\alpha)$ and $\lambda_{j}(\alpha)=\delta_{j-1}, j=2, \ldots, \ell+1$;
(iii) if $\tilde{\alpha}_{i-1}<\alpha<\tilde{\alpha}_{i}, 2 \leq i \leq \ell$, then $\lambda_{i}(\alpha)=\lambda(\alpha), \lambda_{j}(\alpha)=\delta_{j}$ for $j=1, \ldots, i-1$, and $\lambda_{j}(\alpha)=\delta_{j-1}$ for $j=i+1, \ldots, \ell+1$;
(iv) if $\alpha>\tilde{\alpha}_{\ell}$, then $\lambda_{j}(\alpha)=\delta_{j}, j=1,2, \ldots, \ell$ and $\lambda_{\ell+1}(\alpha)=\lambda(\alpha)$.

Proof. These results are a direct consequence of Cauchy's interlace theorem, Lemmas 3.1 and 3.2, and the properties of the functions $\phi(\lambda)$ and $\alpha-\lambda$.

We can expect difficulties in practice when the vector $g$ is nearly orthogonal to the eigenspace $\mathcal{S}_{1}$. If this happens, there still exists $\lambda_{*}<\delta_{1}$ and $x_{*}$ such that $\left(A-\lambda_{*} I\right) x_{*}=-g,\left\|x_{*}\right\|=\Delta$, with $\lambda_{*}$ quite close to $\delta_{1}$. We call this situation a near hard case and Figure 1(c) illustrates it. In the detail shown in Figure 1(d), we can see that in this case, the derivative $\phi^{\prime}$ changes rapidly for $\lambda$ close to $\delta_{1}$, so the problem of finding $\lambda_{*}$ satisfying the correct slope $\phi^{\prime}\left(\lambda_{*}\right)=\Delta^{2}$ is very ill-conditioned.

In the remainder of this section we discuss the hard case and present the results that allow us to compute a nearly optimal solution for the trust-region subproblem in this situation. The hard case can occur only when $g \perp \mathcal{S}_{1}$, the matrix $A$ is indefinite or positive semidefinite and singular, and for certain values of $\Delta$. This case was analyzed in [10] for medium-scale problems and discussed in [13, 17] in the large-scale context. The precise statement is contained in Lemma 3.4 from [15]. We present a slightly different proof here.

Lemma 3.4 (see [15]). Assume $g$ is orthogonal to $\mathcal{S}_{1}$ and let $p=-\left(A-\delta_{1} I\right)^{\dagger} g$. If $\delta_{1} \leq 0$ and $\|p\|<\Delta$, then the solutions of (2.1) consist of the set $\{x \mid x=p+z, z \in$ $\left.\mathcal{S}_{1},\|x\|=\Delta\right\}$, with Lagrange multiplier $\lambda_{*}=\delta_{1}$.

Proof. We need to show that if $x=p+z$ with $z \in \mathcal{S}_{1}$ and $\|x\|=\Delta$, then $\left\{x, \delta_{1}\right\}$ satisfy the conditions in Lemma 2.1. It follows directly from the hypothesis and the fact that $\delta_{1}$ is the smallest eigenvalue of $A$ that $A-\delta_{1} I$ is positive semidefinite, that $\delta_{1} \leq 0$, and that $\delta_{1}(\Delta-\|x\|)=0$. It remains only to show that $\left(A-\delta_{1} I\right) x=-g$. To see this, observe

$$
\begin{aligned}
\left(A-\delta_{1} I\right) x & =\left(A-\delta_{1} I\right)(p+z) \\
& =-\left(A-\delta_{1} I\right)\left(A-\delta_{1} I\right)^{\dagger} g
\end{aligned}
$$

since $z \in \mathcal{N}\left(A-\delta_{1} I\right)$. Now, since $g \in \mathcal{R}\left(A-\delta_{1} I\right)$ and $\left(A-\delta_{1} I\right)\left(A-\delta_{1} I\right)^{\dagger}$ is an orthogonal projector onto $\mathcal{R}\left(A-\delta_{1} I\right)$, we have

$$
\left(A-\delta_{1} I\right) x=-g
$$

which completes the proof.
As we can see, it is precisely in the hard case that in the process of adjusting $\alpha$ we will compute values such that $\alpha>\tilde{\alpha}_{1}$. As Lemma 3.3 establishes, in this case all the eigenvectors corresponding to the smallest eigenvalue of $B_{\alpha}$ have first component zero. Moreover, in a near hard case the eigenvectors will have very small first components and dividing by these values will introduce large roundoff errors. Theorem 3.1 and Lemma 3.3 suggest a strategy for handling this situation, namely, using the eigenvector of $B_{\alpha}$ with the desired structure guaranteed by Theorem 3.1 and the corresponding eigenvalue to obtain the interpolation points, so we can proceed with the adjustment
of the parameter $\alpha$. We will need a safeguarding strategy to enforce convergence of this iteration. We will describe this strategy in the next section where we present the algorithm in detail.

The following results provide the theoretical basis for declaring convergence in the hard case. Results within the same philosophy are presented in [10, 17]. The idea behind these results is to exploit the information available at each iteration and, with practically no additional cost, detect a nearly optimal solution in the hard case or near hard case. Theorem 3.2, Lemma 3.5, and Lemma 3.6 contain these results. Theorem 3.2 establishes that, under certain conditions, the last $n$ components of a special linear combination of eigenvectors of $B_{\alpha}$ form a nearly optimal solution for problem (2.1). Lemma 3.5 establishes the conditions under which we can compute the special linear combination, and Lemma 3.6 shows how to compute it. Theorem 3.2 follows from a more general result from [14] but we present a different proof here. Lemma 3.5 is a reformulation of a result from [14] and Lemma 3.6 is from [14].

ThEOREM 3.2 (see [14]). Let $\lambda_{1}(\alpha)$ be the smallest eigenvalue of $B_{\alpha}$ with a corresponding eigenvector $z_{1}=\left(\nu_{1}, \widetilde{z}_{1}^{T}\right)^{T}$. Let $\lambda_{i}(\alpha)$ be any of the remaining $n$ eigenvalues of $B_{\alpha}$ with a corresponding eigenvector $z_{i}=\left(\nu_{i}, \widetilde{z}_{i}^{T}\right)^{T}$. Define $Z=\left[\begin{array}{ll}z_{1} & z_{i}\end{array}\right]$, $\widetilde{Z}=\left[\begin{array}{ll}\widetilde{z}_{1} & \widetilde{z}_{i}\end{array}\right]$, and assume $Z^{T} Z=I$. Let $\eta>0$.

If there exists $t=\left(\tau_{1}, \tau_{2}\right)^{T}$, with $\|t\|=1$ such that
(i) $\left(e_{1}^{T} Z t\right)^{2}=\frac{1}{1+\Delta^{2}}$, and
(ii) $\left(\lambda_{i}(\alpha)-\lambda_{1}(\alpha)\right) \tau_{2}^{2}\left(1+\Delta^{2}\right) \leq-2 \eta \psi(\widetilde{x})$ for $\widetilde{x}=\frac{\widetilde{Z} t}{e_{1}^{T} Z t}$,
then

$$
\psi\left(x_{*}\right) \leq \psi(\widetilde{x}) \leq \frac{1}{1+\eta} \psi\left(x_{*}\right)
$$

where $x_{*}$ is a boundary solution for problem (2.1) with $\psi\left(x_{*}\right) \leq 0$.
Proof. Since $x_{*}$ is a boundary solution of (2.1), we have $\psi\left(x_{*}\right) \leq \psi(x) \forall x \in \mathbb{R}^{n}$ such that $\|x\|=\Delta$. Therefore, in order to prove that $\psi\left(x_{*}\right) \leq \psi(\widetilde{x})$, it will suffice to show that $\|\widetilde{x}\|=\Delta$.

Note that $\frac{Z t}{e_{1}^{T} Z t}=\left(1, \widetilde{x}^{T}\right)^{T}$ and therefore

$$
\begin{aligned}
\left\|\left(1, \widetilde{x}^{T}\right)\right\|^{2}=1+\|\widetilde{x}\|^{2} & =\left\|\frac{Z t}{e_{1}^{T} Z t}\right\|^{2} \\
& =\frac{1}{\left(e_{1}^{T} Z t\right)^{2}}
\end{aligned}
$$

since $\|t\|=1$ and $Z^{T} Z=I$ by hypothesis. Thus, by (i)

$$
1+\|\widetilde{x}\|^{2}=1+\Delta^{2}
$$

This implies $\|\widetilde{x}\|=\Delta$ and therefore $\psi\left(x_{*}\right) \leq \psi(\widetilde{x})$.
To prove the other part of the inequality, observe that $\alpha+2 \psi\left(x_{*}\right)=$ $\left(1, x_{*}^{T}\right) B_{\alpha}\left(1, x_{*}^{T}\right)^{T}$. Thus, by Rayleigh quotient properties

$$
\alpha+2 \psi\left(x_{*}\right) \geq \lambda_{1}(\alpha)\left\|\left(1, x_{*}^{T}\right)^{T}\right\|^{2} .
$$

Since $\left\|x_{*}\right\|=\Delta$ it follows that $\left\|\left(1, x_{*}^{T}\right)^{T}\right\|^{2}=1+\Delta^{2}$, and therefore

$$
\begin{equation*}
\alpha+2 \psi\left(x_{*}\right) \geq \lambda_{1}(\alpha)\left(1+\Delta^{2}\right) \tag{3.2}
\end{equation*}
$$

We now show that $\alpha+2 \psi(\widetilde{x})=\widetilde{\lambda}\left(1+\Delta^{2}\right)$, with $\widetilde{\lambda}=\lambda_{1}(\alpha) \tau_{1}^{2}+\lambda_{i}(\alpha) \tau_{2}^{2}$. Observe that $\alpha+2 \psi(\widetilde{x})=\left(1, \widetilde{x}^{T}\right) B_{\alpha}\left(1, \widetilde{x}^{T}\right)^{T}$, and since $\left(1, \widetilde{x}^{T}\right)^{T}=\frac{1}{e_{1}^{T} Z t} Z t$, it follows that

$$
\begin{aligned}
\alpha+2 \psi(\widetilde{x}) & =t^{T} Z^{T} B_{\alpha} Z t \frac{1}{\left(e_{1}^{T} Z t\right)^{2}} \\
& =\left[\lambda_{1}(\alpha) \tau_{1}^{2}+\lambda_{i}(\alpha) \tau_{2}^{2}\right]\left(1+\Delta^{2}\right)
\end{aligned}
$$

by (i) and the fact that $z_{1}, z_{i}$ are eigenvectors of $B_{\alpha}$.
Thus, since $\tau_{1}^{2}+\tau_{2}^{2}=1$, we have

$$
\begin{aligned}
\alpha+2 \psi(\widetilde{x}) & =\left[\lambda_{1}(\alpha)\left(1-\tau_{2}^{2}\right)+\lambda_{i}(\alpha) \tau_{2}^{2}\right]\left(1+\Delta^{2}\right) \\
& =\left[\lambda_{1}(\alpha)+\left(\lambda_{i}(\alpha)-\lambda_{1}(\alpha)\right) \tau_{2}^{2}\right]\left(1+\Delta^{2}\right)
\end{aligned}
$$

and therefore

$$
\begin{aligned}
\alpha+2 \psi(\widetilde{x})-\left(\lambda_{i}(\alpha)-\lambda_{1}(\alpha)\right) \tau_{2}^{2}\left(1+\Delta^{2}\right) & =\lambda_{1}(\alpha)\left(1+\Delta^{2}\right) \\
& \leq \alpha+2 \psi\left(x_{*}\right) \quad \text { by }(3.2)
\end{aligned}
$$

If $\quad\left(\lambda_{i}(\alpha)-\lambda_{1}(\alpha)\right) \tau_{2}^{2}\left(1+\Delta^{2}\right) \leq-2 \eta \psi(\widetilde{x})$, then

$$
\alpha+2 \psi(\widetilde{x})+2 \eta \psi(\widetilde{x}) \leq \alpha+2 \psi\left(x_{*}\right)
$$

and we can conclude $\psi(\widetilde{x}) \leq \frac{1}{1+\eta} \psi\left(x_{*}\right)$.
Therefore $\psi\left(x_{*}\right) \leq \psi(\widetilde{x}) \leq \frac{1}{1+\eta} \psi\left(x_{*}\right)$ as claimed.
It follows directly from this result that

$$
\begin{align*}
0 \leq \psi(\widetilde{x})-\psi\left(x_{*}\right) & \leq-\frac{\eta}{1+\eta} \psi\left(x_{*}\right) \\
\left|\psi(\widetilde{x})-\psi\left(x_{*}\right)\right| & \leq \frac{\eta}{1+\eta}\left|\psi\left(x_{*}\right)\right| \tag{3.3}
\end{align*}
$$

The inequality (3.3) implies that under the conditions of Theorem 3.2, $\psi(\widetilde{x})$ will be arbitrarily close to $\psi\left(x_{*}\right)$. We will call such $\widetilde{x}$ a quasi-optimal solution for problem (2.1).

The next result establishes conditions for computing the vector $t$ in Theorem 3.2.
Lemma 3.5 (see [14]). Let $z_{i}=\left(\nu_{i}, \widetilde{z}_{i}^{T}\right)^{T}$, with $\nu_{i} \in \mathbb{R}, \widetilde{z}_{i} \in \mathbb{R}^{n}$ for $i=1,2$. Define the matrices $Z=\left[\begin{array}{ll}z_{1} & z_{2}\end{array}\right]$ and $\widetilde{Z}=\left[\begin{array}{ll}\widetilde{z}_{1} & \widetilde{z}_{2}\end{array}\right]$, and assume $Z^{T} Z=I$. If $\left\|Z^{T} e_{1}\right\|^{2} \geq$ $\frac{1}{\beta}$ for $\beta>0$, then there exists $t \in \mathbb{R}^{2}$ with $t \neq 0$ that satisfies

$$
\begin{equation*}
\|Z t\|^{2}=\beta\left(e_{1}^{T} Z t\right)^{2} \tag{3.4}
\end{equation*}
$$

Proof. Observe that we can rewrite (3.4) as

$$
\begin{aligned}
t^{T} Z^{T} Z t & =\beta\left(e_{1}^{T} Z t\right)^{2} \\
& =\beta\left(t^{T} Z^{T} e_{1} e_{1}^{T} Z t\right)
\end{aligned}
$$

which is equivalent to

$$
\begin{equation*}
t^{T}\left[I-\beta Z^{T} e_{1} e_{1}^{T} Z\right] t=0 \tag{3.5}
\end{equation*}
$$

since $Z^{T} Z=I$ by hypothesis. Equation (3.5) will have a nontrivial solution only if the matrix $M=I-\beta Z^{T} e_{1} e_{1}^{T} Z$ is indefinite or positive semidefinite and singular. So, let
us study the eigenvalues of $M$. The two eigenpairs of the matrix $M=I-\beta Z^{T} e_{1} e_{1}^{T} Z$ are given by

$$
\left\{1-\beta e_{1}^{T} Z Z^{T} e_{1}, Z^{T} e_{1}\right\} \quad \text { and } \quad\{1, v\} \quad \text { with } v \perp Z^{T} e_{1}
$$

Therefore, (3.5) will have nontrivial solutions if $1-\beta e_{1}^{T} Z Z^{T} e_{1} \leq 0$. In other words, if $\left\|Z^{T} e_{1}\right\|^{2}=e_{1}^{T} Z Z^{T} e_{1} \geq \frac{1}{\beta}$, then there exists $t \in \mathbb{R}^{2}$ with $t \neq 0$ such that $t$ satisfies (3.5).

Note that choosing $\beta=1+\Delta^{2}$ in Lemma 3.5 and normalizing $t$ such that $\|t\|=1$ will give a vector that satisfies the conditions of Theorem 3.2. The following lemma provides a way of computing such a vector.

Lemma 3.6 (see [14]). Let $\beta \in \mathbb{R}, \beta>0$, and let $z \in \mathbb{R}^{n}$. The equation

$$
\begin{equation*}
t^{T}\left[I-\beta z z^{T}\right] t=0 \tag{3.6}
\end{equation*}
$$

in $t$ with $t \in \mathbb{R}^{n}$ has $2(n-1)$ nontrivial solutions if the matrix $M=I-\beta z z^{T}$ is indefinite and has one nontrivial solution if $M$ is positive semidefinite and singular.

Proof. Let $P \in \mathbb{R}^{n \times n}$ be such that $P^{T} z=\|z\| e_{1}$ with $P^{T} P=I$ and apply this orthogonal transformation to the matrix $M$ to obtain

$$
P^{T}\left[I-\beta z z^{T}\right] P=I-\beta\|z\|^{2} e_{1} e_{1}^{T}
$$

Therefore, the solutions of (3.6) in this new basis are the solutions of

$$
y^{T}\left(\begin{array}{cc}
-\theta & 0 \\
0 & I
\end{array}\right) y=0
$$

where $\theta=-1+\beta\|z\|^{2} e_{1} e_{1}^{T}$.
The nontrivial solutions of (3.6) are then given by $t=P y$, where
(1) $y=\left(1, \sqrt{\theta} e_{i}^{T}\right)^{T}$ and $y=\left(-1, \sqrt{\theta} e_{i}^{T}\right)^{T}$ with $e_{i}$ the $i$ th canonical vector in $\mathbb{R}^{n-1}$, $i=1,2, \ldots, n-1$, if $M$ is indefinite, i.e., if $\theta>0$, or
(2) $y=e_{1}$, if $M$ is positive semidefinite and singular, i.e., if $\theta=0$.

Therefore, (3.6) has $2(n-1)$ nontrivial solutions if $M$ is indefinite and has one nontrivial solution if $M$ is positive semidefinite and singular.

Remark. Suppose $n=2$ and $z=\left(\nu_{1}, \nu_{i}\right)^{T}$ in Lemma 3.6. Then if $\nu_{1}{ }^{2}+\nu_{i}{ }^{2}>\frac{1}{\beta}$, the vector $t=\left(\tau_{1}, \tau_{2}\right)^{T}$ is given by

$$
\tau_{1}=\frac{\nu_{1}-\nu_{i} \sqrt{\beta\left(\nu_{1}^{2}+\nu_{i}^{2}\right)-1}}{\left(\nu_{1}^{2}+\nu_{i}^{2}\right) \sqrt{\beta}}, \quad \tau_{2}=\frac{\nu_{1}+\nu_{i} \sqrt{\beta\left(\nu_{1}^{2}+\nu_{i}^{2}\right)-1}}{\left(\nu_{1}^{2}+\nu_{i}^{2}\right) \sqrt{\beta}}
$$

or

$$
\tau_{1}=\frac{\nu_{i}+\nu_{1} \sqrt{\beta\left(\nu_{1}^{2}+\nu_{i}^{2}\right)-1}}{\left(\nu_{1}^{2}+\nu_{i}^{2}\right) \sqrt{\beta}}, \quad \tau_{2}=\frac{\nu_{1}-\nu_{i} \sqrt{\beta\left(\nu_{1}^{2}+\nu_{i}^{2}\right)-1}}{\left(\nu_{1}^{2}+\nu_{i}^{2}\right) \sqrt{\beta}} .
$$

If $\nu_{1}{ }^{2}+\nu_{i}{ }^{2}=\frac{1}{\beta}$, then $t$ is given by

$$
\tau_{1}=\frac{\nu_{1}}{\sqrt{\nu_{1}^{2}+\nu_{i}^{2}}}, \quad \tau_{2}=\frac{\nu_{i}}{\sqrt{\nu_{1}^{2}+\nu_{i}^{2}}} .
$$

The previous results are the basis for the algorithm in the next section. They provide the necessary tools for handling the hard case and the standard case in the same iteration and for computing a solution in the hard case. Theorem 3.2 and

Lemma 3.5 give mechanisms for approximating the vector $x=p+z$ in Lemma 3.4 from a linear combination of eigenvectors of $B_{\alpha}$. Theorem 3.2 also establishes conditions under which the vector $x$ computed in this way is a quasi-optimal solution for problem (2.1). Moreover, Theorem 3.2 guarantees that if the second smallest eigenvalue of $B_{\alpha}$ belongs to a cluster, as is the case, for example, in discrete ill-posed problems, we can still build the special vector $\hat{x}$ from an eigenvector associated with the smallest eigenvalue of $B_{\alpha}$ and from an eigenvector associated with any eigenvalue of the cluster-not necessarily the second smallest. Observe that Lemmas 3.5 and 3.6, respectively, provide a way of computing the vectors $\hat{x}$ and $t$ needed in Theorem 3.2. We use Theorem 3.2, Lemma 3.5, and Lemma 3.6 in one of the stopping rules in our method. We describe the stopping criteria in section 4.5.
4. The algorithm. Keeping in mind the availability of a well-suited variant of the Lanczos method, namely, the implicitly restarted Lanczos method (cf. [16]), we will develop a rapidly convergent iteration to adjust $\alpha$ based on this process. Our goal is to adjust $\alpha$ so that

$$
\alpha-\lambda=\phi(\lambda), \quad \phi^{\prime}(\lambda)=\Delta^{2}
$$

where

$$
\phi(\lambda)=-g^{T} x, \quad \phi^{\prime}(\lambda)=x^{T} x,
$$

with $(A-\lambda I) x=-g$.
The approach of this work is similar to the one in [17] in the following sense. We compute a function $\widehat{\phi}$ which interpolates $\phi$ and $\phi^{\prime}$ at two properly chosen points. Then, from the interpolating function $\widehat{\phi}$ we determine $\widehat{\lambda}$ satisfying

$$
\begin{equation*}
\widehat{\phi}^{\prime}(\widehat{\lambda})=\Delta^{2} \tag{4.1}
\end{equation*}
$$

Finally, we use $\widehat{\lambda}$ and $\phi(\widehat{\lambda})$ to update the parameter $\alpha$ and compute the next iterates $\{\lambda, x\}$. The new elements in our algorithm are the introduction of safeguards for the sequence in $\alpha$, the use of the information relative to the second smallest eigenvalue of the matrix $B_{\alpha}$, and the introduction of a different interpolating scheme, where the currently available information is exploited to a greater extent. Considering that the interpretation of the primal feasibility equations of [13] can be related to (4.1), the description of our algorithm also has some flavor of the approach in [13], where an inverse interpolation scheme is used to satisfy primal feasibility. However, in the presence of the hard case, we do not need to combine distinct interpolating functions, as in [13], nor switch to another algorithm as in [17]. In this section we will assume that the vector $g$ is nonzero. If $g=0$, then problem (2.1) reduces to solving an eigenvalue problem for the smallest eigenvalue of $A$. We shall first describe the components of the algorithm and then present the complete method.
4.1. Interpolating schemes. To begin the iteration, we need a single-point interpolating scheme. We use the approach derived in [17] which gives the following expression for $\alpha_{1}$ :

$$
\begin{equation*}
\alpha_{1}=\widehat{\lambda}+\widehat{\phi}(\widehat{\lambda})=\alpha_{0}+\frac{\alpha_{0}-\lambda_{0}}{\left\|x_{0}\right\|}\left(\frac{\Delta-\left\|x_{0}\right\|}{\Delta}\right)\left(\Delta+\frac{1}{\left\|x_{0}\right\|}\right) \tag{4.2}
\end{equation*}
$$

where

$$
\widehat{\lambda}=\delta+\frac{g^{T} x_{0}}{\left\|x_{0}\right\| \Delta}
$$

This method is linearly convergent and may be slow in some cases, so we will use it just to obtain a second pair of iterates, which together with $\lambda_{0}, x_{0}$ will be the starting values for a two-point method. In the two-point method we use the four pieces of information available at the $k$ th iteration, namely, $\phi\left(\lambda_{k-1}\right), \phi^{\prime}\left(\lambda_{k-1}\right), \phi\left(\lambda_{k}\right)$, and $\phi^{\prime}\left(\lambda_{k}\right)$ as follows. We compute $\widehat{\lambda}$ such that

$$
\begin{equation*}
\frac{1}{\Delta}=\frac{1}{\sqrt{\phi^{\prime}\left(\lambda_{k-1}\right)}}\left(\frac{\lambda_{k}-\hat{\lambda}}{\lambda_{k}-\lambda_{k-1}}\right)+\frac{1}{\sqrt{\phi^{\prime}\left(\lambda_{k}\right)}}\left(\frac{\hat{\lambda}-\lambda_{k-1}}{\lambda_{k}-\lambda_{k-1}}\right) \tag{4.3}
\end{equation*}
$$

obtaining

$$
\begin{equation*}
\widehat{\lambda}=\frac{\lambda_{k-1}\left\|x_{k-1}\right\|\left(\left\|x_{k}\right\|-\Delta\right)+\lambda_{k}\left\|x_{k}\right\|\left(\Delta-\left\|x_{k-1}\right\|\right)}{\Delta\left(\left\|x_{k}\right\|-\left\|x_{k-1}\right\|\right)} \tag{4.4}
\end{equation*}
$$

This is equivalent to defining

$$
\begin{equation*}
\widehat{\phi}(\lambda)=\frac{\gamma^{2}}{\delta-\lambda}+\eta \tag{4.5}
\end{equation*}
$$

for any $\eta$ and computing $\widehat{\lambda}$ such that $\frac{1}{\sqrt{\widehat{\phi}^{\prime}(\widehat{\lambda})}}=\frac{1}{\Delta}$. It is easy to verify using (4.3) that

$$
\gamma^{2}=\frac{\left(\lambda_{k}-\lambda_{k-1}\right)^{2}\left\|x_{k-1}\right\|^{2}\left\|x_{k}\right\|^{2}}{\left(\left\|x_{k}\right\|-\left\|x_{k-1}\right\|\right)^{2}} \quad \text { and } \quad \delta=\frac{\lambda_{k}\left\|x_{k}\right\|-\lambda_{k-1}\left\|x_{k-1}\right\|}{\left\|x_{k}\right\|-\left\|x_{k-1}\right\|}
$$

Ideally, $\eta=\phi(\widehat{\lambda})-\frac{\gamma^{2}}{\delta-\hat{\lambda}}$, where $\phi(\widehat{\lambda})$ is the value we are going to estimate in order to update $\alpha$. Using the values $\phi\left(\lambda_{k-1}\right)$ and $\phi\left(\lambda_{k}\right)$, we first define $\eta_{j}=\phi\left(\lambda_{j}\right)-\frac{\gamma^{2}}{\delta-\lambda_{j}}$, for $j=k-1, k$. Then, applying the linear interpolation philosophy on $\lambda_{j}, \eta_{j}$, and defining the weights by means of the already computed value $\widehat{\lambda}$, we choose

$$
\eta=\left(\frac{\lambda_{k}-\widehat{\lambda}}{\lambda_{k}-\lambda_{k-1}}\right) \eta_{k-1}+\left(\frac{\widehat{\lambda}-\lambda_{k-1}}{\lambda_{k}-\lambda_{k-1}}\right) \eta_{k}
$$

After some manipulation we can express the updating formula for $\alpha$ as

$$
\begin{align*}
\alpha_{k+1}= & \widehat{\lambda}+\omega \phi\left(\lambda_{k-1}\right)+(1-\omega) \phi\left(\lambda_{k}\right) \\
& +\frac{\left\|x_{k-1}\right\|\left\|x_{k}\right\|\left(\left\|x_{k}\right\|-\left\|x_{k-1}\right\|\right)}{\omega\left\|x_{k}\right\|+(1-\omega)\left\|x_{k-1}\right\|} \frac{\left(\lambda_{k-1}-\widehat{\lambda}\right)\left(\lambda_{k}-\widehat{\lambda}\right)}{\left(\lambda_{k}-\lambda_{k-1}\right)} \\
= & \omega \alpha_{k-1}+(1-\omega) \alpha_{k} \\
& +\frac{\left\|x_{k-1}\right\|\left\|x_{k}\right\|\left(\left\|x_{k}\right\|-\left\|x_{k-1}\right\|\right)}{\omega\left\|x_{k}\right\|+(1-\omega)\left\|x_{k-1}\right\|} \frac{\left(\lambda_{k-1}-\widehat{\lambda}\right)\left(\lambda_{k}-\widehat{\lambda}\right)}{\left(\lambda_{k}-\lambda_{k-1}\right)} \tag{4.6}
\end{align*}
$$

where $\omega=\frac{\lambda_{k}-\widehat{\lambda}}{\lambda_{k}-\lambda_{k-1}}, \alpha_{k-1}=\lambda_{k-1}+\phi\left(\lambda_{k-1}\right)$, and $\alpha_{k}=\lambda_{k}+\phi\left(\lambda_{k}\right)$.
As we discussed in section 3 , we need a special strategy to obtain interpolation points in potential hard cases. We describe this strategy in section 4.2.
4.2. Choice of interpolation points. According to Lemma 3.1, if the first component of the eigenvector corresponding to the smallest eigenvalue of $B_{\alpha_{k}}$ is zero, this will indicate a potential hard case and we will have $\lambda_{1}\left(\alpha_{k}\right)=\delta_{1}$. However,

Lemma 3.3 establishes that for $\alpha_{k}$ slightly larger than $\tilde{\alpha}_{1}$ there is an eigenvector with significant first component that corresponds to the second smallest eigenvalue of $B_{\alpha}$. Therefore, we propose to use an eigenpair corresponding to an eigenvalue that is close to the second smallest eigenvalue of the bordered matrix to obtain the interpolation point whenever we detect a potential hard case. As we shall explain, not only can we keep the size of the iterate $x_{k}$ under control, but we can also ensure convergence of $\left\{\lambda_{k}, x_{k}\right\}$ to $\left\{\delta_{1}, p\right\}$ by driving the parameter $\alpha_{k}$ to the value $\tilde{\alpha}_{1}$ given by Lemma 3.2. Recall that Lemma 3.2 established that there will be an eigenvector with significant first component corresponding to $\lambda_{1}\left(\alpha_{k}\right)$ precisely when $\alpha_{k}$ assumes the special value $\tilde{\alpha}_{1}=\delta_{1}-g^{T} p$. Moreover, the use of a second eigenvector prevents numerical difficulties in a near-hard-case situation.

There is an easy way to detect a potential hard case during an iteration. Let $\left(\nu_{1}, u_{1}^{T}\right)^{T}$ be a unitary eigenvector of $B_{\alpha_{k}}$ corresponding to $\lambda_{1}\left(\alpha_{k}\right)$. Then, we declare $\nu_{1}$ to be "small," indicating a near hard case has been detected, if the condition $\|g\|\left|\nu_{1}\right| \leq \varepsilon \sqrt{1-\nu_{1}{ }^{2}}$ holds for a given $\varepsilon \in(0,1)$. This is motivated as follows. Since $\left(A-\lambda_{1}\left(\alpha_{k}\right) I\right) u_{1}=-g \nu_{1}$, we have

$$
\frac{\left\|\left(A-\lambda_{1}\left(\alpha_{k}\right) I\right) u_{1}\right\|}{\left\|u_{1}\right\|}=\frac{\|g\|\left|\nu_{1}\right|}{\sqrt{1-\nu_{1}{ }^{2}}}
$$

and hence $\|g\|\left|\nu_{1}\right| \leq \varepsilon \sqrt{1-\nu_{1}^{2}}$ ensures that $\left\|\left(A-\lambda_{1}\left(\alpha_{k}\right) I\right) u_{1}\right\| \leq \varepsilon\left\|u_{1}\right\|$. In other words, $\left\{\lambda_{1}\left(\alpha_{k}\right), u_{1}\right\}$ is an approximate eigenpair of $A$ and the eigenvector $\left(\nu_{1}, u_{1}^{T}\right)^{T}$ from the bordered matrix is essentially impossible to normalize. This is approximately the situation described in Lemma 3.1. Of course, this test can be made scale independent by choosing $\varepsilon=\widehat{\varepsilon}\|A\|$, for $\widehat{\varepsilon} \in(0,1)$.

When a near hard case has been detected, we need an alternative way to define the pair $\left\{\lambda_{k}, x_{k}\right\}$. At each iteration, at essentially no extra cost, we compute an eigenpair corresponding to the smallest eigenvalue of $B_{\alpha_{k}}$, which we denote by $\left\{\lambda_{1}\left(\alpha_{k}\right),\left(\nu_{1}, u_{1}^{T}\right)^{T}\right\}$, and also an eigenpair corresponding to an eigenvalue close to the second smallest eigenvalue of $B_{\alpha}$, which we denote by $\left\{\lambda_{i}\left(\alpha_{k}\right),\left(\nu_{2}, u_{2}^{T}\right)^{T}\right\}$. If both $\left|\nu_{1}\right|$ and $\left|\nu_{2}\right|$ are small, that is, if $\|g\|\left|\nu_{1}\right| \leq \varepsilon \sqrt{1-\nu_{1}^{2}}$ and $\|g\|\left|\nu_{2}\right| \leq \varepsilon \sqrt{1-\nu_{2}^{2}}{ }^{2}$, then we decrease the parameter $\alpha_{k}$. According to Theorem 3.1 there always exists an eigenvector of the bordered matrix with significant first component for any value of $\alpha$ and, as we mentioned before, according to Lemma 3.3, as $\alpha_{k}$ approaches the critical value, this normalizable eigenvector will correspond either to the first or to the second smallest eigenvalue of $B_{\alpha_{k}}$. In other words, for values of $\alpha_{k}$ near the critical value, either $\|g\|\left|\nu_{1}\right|>\varepsilon \sqrt{1-\nu_{1}{ }^{2}}$ or $\|g\|\left|\nu_{2}\right|>\varepsilon \sqrt{1-\nu_{2}{ }^{2}}$ will hold. Hence, after a possible reduction of the parameter $\alpha_{k}$, the pair $\left\{\lambda_{k}, x_{k}\right\}$ is well defined if we compute it by the following procedure:

$$
\begin{aligned}
& \text { If }\|g\|\left|\nu_{1}\right| \leq \varepsilon \sqrt{1-\nu_{1}^{2}} \text {, then set } \lambda_{k}=\lambda_{i}\left(\alpha_{k}\right) \text { and } x_{k}=\frac{u_{2}}{\nu_{2}} \\
& \text { Otherwise, set } \lambda_{k}=\lambda_{1}\left(\alpha_{k}\right) \text { and } x_{k}=\frac{u_{1}}{\nu_{1}}
\end{aligned}
$$

Since $\lambda_{k-1}$ and $\lambda_{k}$ are not constrained to ( $-\infty, \delta_{1}$ ] but might belong to the interval $\left(\delta_{1}, \delta_{\ell+1}\right)$, the value $\widehat{\lambda}$ given by (4.4) may be greater than $\delta_{1}$. In this case, we set $\hat{\lambda}=\delta_{U}$, where $\delta_{U}$ is an upper bound for $\delta_{1}$. In section 4.3 we will show how to obtain an initial value for $\delta_{U}$ and how to update this value. We will also show how to safeguard $\alpha$ computed by (4.6).
4.3. Safeguarding. We need to introduce safeguarding to ensure global convergence of the iteration. Let $\lambda_{*}, x_{*}$ be an optimal pair for problem (2.1), satisfying the conditions in Lemma 2.1, except when there is only an interior solution, in which case we define $x_{*}=-\left(A-\lambda_{*} I\right)^{\dagger} g$ such that $\left\|x_{*}\right\|=\Delta$. Let $\alpha_{*}=\lambda_{*}-g^{T} x_{*}$. Rendl and Wolkowicz [13] presented the following bounds for the optimal parameter $\alpha_{*}$ :

$$
\begin{equation*}
\delta_{1}-\frac{\|g\|}{\Delta} \leq \alpha_{*} \leq \delta_{1}+\|g\| \Delta \tag{4.7}
\end{equation*}
$$

Computing a good approximation to $\delta_{1}$ can be nearly as expensive as solving the given trust-region subproblem. For this reason, as observed in [13], we shall replace the above bounds by some simple alternatives. First, note that any Rayleigh quotient $\delta_{U} \equiv \frac{v^{T} A v}{v^{T} v}$ gives an upper bound for $\delta_{1}$. Therefore, if the diagonal of the matrix $A$ is explicitly available, we take $\delta_{U}=\min \left\{a_{i i} \mid i=1, \ldots, n\right\} ;$ otherwise we take $\delta_{U} \equiv \frac{v^{T} A v}{v^{T} v}$, where $v$ is a random vector. From (4.7) we see that $\alpha_{*} \leq \alpha_{U}$, for $\alpha_{U}=\delta_{U}+\|g\| \Delta$. Since $\alpha \leq 0$ implies $B_{\alpha}$ is not positive definite, we set $\alpha_{0}=\min \left\{0, \alpha_{U}\right\}$ to ensure that $\lambda_{1}\left(\alpha_{0}\right) \leq 0$. After solving for $\lambda_{1}\left(\alpha_{0}\right)$ and setting $\delta_{L}=\lambda_{1}\left(\alpha_{0}\right)$ and $\alpha_{L}=\delta_{L}-\frac{\|g\|}{\Delta}$, we immediately have that $\alpha_{L} \leq \alpha_{*}$, since the interlacing property implies $\delta_{L} \leq \delta_{1}$. Using this simple scheme to obtain $\delta_{L}$ and $\delta_{U}$ as initial lower and upper bounds for $\delta_{1}$, we can start with

$$
\begin{equation*}
\alpha_{L}=\delta_{L}-\frac{\|g\|}{\Delta} \quad \text { and } \quad \alpha_{U}=\delta_{U}+\|g\| \Delta \tag{4.8}
\end{equation*}
$$

We update the upper bound $\delta_{U}$ at each iteration using information from the eigenpair corresponding to the smallest eigenvalue of the bordered matrix in the following way: $\delta_{U}=\min \left\{\delta_{U}, \frac{u_{1}^{T} A u_{1}}{u_{1}^{T} u_{1}}\right\}$, where $\frac{u_{1}^{T} A u_{1}}{u_{1}^{T} u_{1}}=\lambda_{1}\left(\alpha_{k}\right)-\nu_{1} \frac{g^{T} u_{1}}{u_{1}^{T} u_{1}}$. As stated in section 4.2, whenever we detect a potential hard case, $\left\{\lambda_{1}\left(\alpha_{k}\right), u_{1}\right\}$ approximates an eigenpair of $A$ and $\lambda_{1}\left(\alpha_{k}\right)$ is a very good approximation to $\delta_{1}$. Thus, $\delta_{U}$ becomes a sharp estimate of $\delta_{1}$ in this case.

At every iteration, we update one of the safeguarding bounds $\alpha_{L}$ or $\alpha_{U}$ so that we always reduce the length of the interval $\left[\alpha_{L}, \alpha_{U}\right]$. In case the value $\alpha_{k+1}$ predicted by the interpolating schemes (4.2) or (4.6) does not belong to the current safeguarding interval, we redefine $\alpha_{k+1}$ by means of a linear adjustment based on the upper bound $\delta_{U}$. If this choice is not in the interval $\left[\alpha_{L}, \alpha_{U}\right]$, we simply set $\alpha_{k+1}=\frac{\alpha_{L}+\alpha_{U}}{2}$.
4.4. Initialization of $\boldsymbol{\alpha}$. As mentioned in section 4.3, there is a simple choice for initializing $\alpha$, given by $\alpha_{0}=\min \left\{0, \alpha_{U}\right\}$, with $\alpha_{U}$ as in (4.8). This ensures that $\lambda_{1}\left(\alpha_{0}\right) \leq 0$ but it has no additional properties. In an attempt to improve this initial guess, we have developed a more sophisticated hot-start strategy based on the Lanczos process. To begin, we compute the following $j$-step Lanczos factorization for the $j$ smallest eigenvalues of $A$ :

$$
\begin{equation*}
A V=V T+f e_{j}^{T} \tag{4.9}
\end{equation*}
$$

where $V^{T} V=I_{j}$, with $I_{j}$ the identity matrix of order $j(j \ll n), T \in \mathbb{R}^{j \times j}$ tridiagonal, $V^{T} f=0$, and $e_{j}$ denotes the $j$ th canonical unit vector in $\mathbb{R}^{j}$.

The hot-start strategy consists of first changing variables in (2.1) using $x=V y$ and solving the $j$-dimensional problem

$$
\begin{aligned}
\min & \frac{1}{2} y^{T} T y+g^{T} V y \\
\text { s.t. } & \|y\| \leq \Delta
\end{aligned}
$$

Then, we compute a solution $\left\{\theta_{*}, y_{*}\right\}$ to this lower dimensional trust-region subproblem by using the algorithm in [10], based on the Cholesky factorization of the tridiagonal matrix $T-\theta I, \theta<\delta_{1}$. The initial value to be used is $\alpha=\theta_{*}-g^{T} V y_{*}$.

We now show that we can use (4.9) to compute an eigenpair corresponding to the smallest eigenvalue of $B_{\alpha_{0}}$. Observe

$$
\left(\begin{array}{cc}
\alpha_{0} & g^{T}  \tag{4.10}\\
g & A
\end{array}\right)\left(\begin{array}{cc}
1 & 0 \\
0 & V
\end{array}\right)=\left(\begin{array}{cc}
1 & 0 \\
0 & V
\end{array}\right)\left(\begin{array}{cc}
\alpha_{0} & g^{T} V \\
V^{T} g & T
\end{array}\right)+\binom{0}{f} e_{j+1}^{T}
$$

If we run the standard Lanczos process for $A$ using $v_{1}=g /\|g\|$ as the initial vector, then we obtain a tridiagonal matrix on the right-hand side of (4.10). This provides a way of computing the smallest eigenvalue of $B_{\alpha_{0}}$.

In numerical experiments, the use of this hot start for $\alpha$ did not substantially improve the performance of the method.
4.5. Stopping criteria. At each iteration we check for a boundary solution, an interior solution, or a quasi-optimal solution according to Theorem 3.2. We can also stop if we reach a maximum number of iterations or if the length of the safeguarding interval is too small. Given the tolerances $\varepsilon_{\Delta}, \varepsilon_{H C}, \varepsilon_{\alpha} \in(0,1)$, and $\varepsilon_{\text {Int }} \in[0,1)$, we declare convergence of the algorithm according to the following criteria. Let $\left(\nu_{1}, u_{1}^{T}\right)^{T}$ be the eigenvector corresponding to $\lambda_{1}\left(\alpha_{k}\right)$ and let $\left\{\lambda_{k}, x_{k}\right\}$ be the current iterates; then we can write the stopping criteria in the following way.

1. Boundary solution.

We detect a boundary solution if

$$
\left(\left|\left\|x_{k}\right\|-\Delta\right| \leq \varepsilon_{\Delta} * \Delta\right) \quad \text { and } \quad\left(\lambda_{1}\left(\alpha_{k}\right) \leq 0\right)
$$

If this condition is satisfied, the solution is

$$
\lambda_{*}=\lambda_{1}\left(\alpha_{k}\right) \quad \text { and } \quad x_{*}=x_{k} .
$$

## 2. Interior solution.

We detect an interior solution if

$$
\left(\left\|u_{1}\right\|<\Delta\left|\nu_{1}\right|\right) \quad \text { and } \quad\left(\lambda_{1}\left(\alpha_{k}\right)>-\varepsilon_{\text {Int }}\right)
$$

In this case, the solution is $\lambda_{*}, x_{*}$, where $\lambda_{*}=0$ and $x_{*}$ satisfies the linear system $A x=-g$, with $A$ positive definite. The conjugate gradient method is a natural choice for solving this system for most large-scale problems.
3. Quasi-optimal solution.

To declare that we have found a quasi-optimal solution, we first compute $t$ and $\widetilde{x}$ as in Lemma 3.5, provided that the conditions of the lemma are satisfied. If $t=\left(\tau_{1}, \tau_{2}\right)^{T}$ and $\widetilde{x}$ satisfy condition (ii) of Theorem 3.2, then $\widetilde{x}$ is a quasi-optimal solution for problem (2.1) and we set $\lambda_{*}=\widetilde{\lambda}$ and $x_{*}=\widetilde{x}$.
4. The safeguarding interval is too small.

If $\left|\alpha_{U}-\alpha_{L}\right| \leq \varepsilon_{\alpha} \max \left\{\left|\alpha_{L}\right|,\left|\alpha_{U}\right|\right\}$, then we stop the iteration and set $\lambda_{*}=\lambda_{1}\left(\alpha_{k}\right)$. If this criterion is satisfied and we do not have a boundary solution, then we are in the hard case and $\alpha_{*}$ is within $\varepsilon_{\alpha}$ of $\tilde{\alpha}_{1}$. If $\nu_{1}$ is large enough, we set $p=\frac{u_{1}}{\nu_{1}}$. Since $\|p\|<\Delta$ in this case, we compute $x_{*}$ as $x_{*}=p+\tau z$ such that $\left\|x_{*}\right\|=\Delta$, where the vector $z$ is an approximate eigenvector associated with the smallest eigenvalue of $A$. Of the two possible choices for $\tau$, we choose the one with smaller magnitude since this
value minimizes $\psi(p+\tau z)$ (see [10, p. 558]). This choice of $\tau$ is given by

$$
\tau=\frac{\Delta^{2}-\|p\|^{2}}{p^{T} z+\operatorname{sign}\left(p^{T} z\right) \sqrt{\left(p^{T} z\right)^{2}-\left(\Delta^{2}-\|p\|^{2}\right)}} .
$$

The vector $z$ is usually available in potential hard cases since in those cases the eigenvectors corresponding to the smallest eigenvalue of $B_{\alpha_{k}}$ will often have a small first component. In the rather unlikely situation where this vector is not available, we increase the parameter and solve an eigenproblem for the smallest eigenvalue of the bordered matrix. This strategy will provide an approximate vector in $\mathcal{S}_{1}$ as Lemma 3.3 guarantees.

If $\nu_{1}$ is too small or zero, we cannot compute a solution. This situation can arise in practice because the eigensolver might not provide the eigenvector with significant first component that the theory guarantees. We have not encountered this case in our experiments.
4.6. The algorithm. Let us now put all these pieces together and present LSTRS, our algorithm for the large-scale trust-region subproblem. We describe steps 2.1 and 2.5 of Algorithm 4.1 separately. In step 2.1 we adjust the parameter $\alpha_{k}$ so that the eigenvector corresponding to the smallest eigenvalue, or to an eigenvalue equal or close to the second smallest eigenvalue of $B_{\alpha_{k}}$, has a significant first component. We might reduce the interval $\left[\alpha_{L}, \alpha_{U}\right]$ during this adjustment. In step 2.5 we correct the parameter predicted by the interpolation schemes in case it does not belong to the current safeguarding interval $\left[\alpha_{L}, \alpha_{U}\right]$. We try a linear adjustment first and adopt the middle point of the current interval as a last resort. Figure 2 shows Algorithm 4.1, while Figures 3 and 4 show steps 2.1 and 2.5, respectively.

## 5. Convergence analysis.

### 5.1. Iterates are well defined.

Lemma 5.1. The iterates generated by Algorithm 4.1 are well defined.
Proof. In order to define the current iterate $x_{k}$ in Algorithm 4.1, we must ensure that we can safely normalize an eigenvector, corresponding to either the smallest eigenvalue or a value equal or close to the second smallest eigenvalue of $B_{\alpha_{k}}$, to have first component one. This is accomplished in step 2.1, where we adjust the parameter $\alpha_{k}$ until we can normalize one of these two eigenvectors to have first component one. Theorem 3.1 and Lemma 3.3 guarantee that the adjusting procedure in step 2.1 yields a value of $\alpha$ such that there exists an eigenvector that has significant first component and is associated with the smallest eigenvalue or a value equal or close to the second smallest eigenvalue of $B_{\alpha}$.

### 5.2. Local convergence.

### 5.2.1. Preliminary results.

Lemma 5.2. Let $\lambda_{k}, x_{k}$ be the iterates at iteration $k$ of Algorithm 4.1. Then

$$
g \in \mathcal{R}\left(A-\lambda_{k} I\right)
$$

Proof. If $\lambda_{k}, x_{k}$ are the iterates at iteration $k$ of Algorithm 4.1, then

$$
\left(\begin{array}{cc}
\alpha_{k} & g^{T} \\
g & A
\end{array}\right)\binom{1}{x_{k}}=\lambda_{k}\binom{1}{x_{k}} .
$$

Therefore, $\left(A-\lambda_{k} I\right) x_{k}=-g$, which implies that $g \in \mathcal{R}\left(A-\lambda_{k} I\right)$.

Algorithm 4.1. LSTRS.
Input: $A \in \mathbb{R}^{n \times n}, g \in \mathbb{R}^{n}, \Delta>0, \varepsilon_{\Delta}, \varepsilon_{\nu}, \varepsilon_{H C}, \varepsilon_{\alpha} \in(0,1), \varepsilon_{\text {Int }} \in[0,1)$.
Output: $\lambda_{*}, x_{*}$ satisfying conditions of Lemma 2.1.

1. Initialization
1.1 Compute $\delta_{U} \geq \delta_{1}$, initialize $\alpha_{U}$ using (4.8), set $\alpha_{0}=\min \left\{0, \alpha_{U}\right\}$
1.2 Compute eigenpairs $\left\{\lambda_{1}\left(\alpha_{0}\right),\left(\nu_{1}, u_{1}^{T}\right)^{T}\right\},\left\{\lambda_{i}\left(\alpha_{0}\right),\left(\nu_{2}, u_{2}^{T}\right)^{T}\right\}$
corresponding to smallest eigenvalue and an eigenvalue close to second smallest eigenvalue of $B_{\alpha_{0}}$
1.3 Initialize $\alpha_{L}$ using (4.8)
1.4 Set $k=0$
2. repeat
2.1 Adjust $\alpha_{k}$
2.2 Update $\delta_{U}=\min \left\{\delta_{U}, \frac{u_{1}^{T} A u_{1}}{u_{1}^{T} u_{1}}\right\}$
2.3 if $\|g\|\left|\nu_{1}\right|>\varepsilon_{\nu} \sqrt{1-\nu_{1}^{2}}$ then
set $\lambda_{k}=\lambda_{1}\left(\alpha_{k}\right)$ and $x_{k}=\frac{u_{1}}{\nu_{1}}$
if $\left\|x_{k}\right\|<\Delta$ then $\alpha_{L}=\alpha_{k}$ end if
if $\left\|x_{k}\right\|>\Delta$ then $\alpha_{U}=\alpha_{k}$
else set $\lambda_{k}=\lambda_{i}\left(\alpha_{k}\right), x_{k}=\frac{u_{2}}{\nu_{2}}$ and $\alpha_{U}=\alpha_{k}$ end if end if
2.4 Compute $\alpha_{k+1}$ by interpolation scheme using (4.2) if $k=0$ or (4.4) and (4.6) otherwise
2.5 Safeguard $\alpha_{k+1}$
2.6 Set $k=k+1$
until convergence

Fig. 2. LSTRS: A method for the large-scale trust-region subproblem.

LEMMA 5.3. Let $\lambda_{*} \leq \delta_{1}$ be the Lagrange multiplier corresponding to a boundary solution of problem (2.1). Then

$$
g \in \mathcal{R}\left(A-\lambda_{*} I\right)
$$

Proof. If $\lambda_{*}<\delta_{1}$, then $A-\lambda_{*} I$ is nonsingular and $g \in \mathcal{R}\left(A-\lambda_{*} I\right)$. If $\lambda_{*}=\delta_{1}$, then $g \perp \mathcal{N}\left(A-\lambda_{*} I\right)$ must hold and therefore $g \in \mathcal{R}\left(A-\lambda_{*} I\right)$. $\quad$.

Remark. Since $(A-\lambda I)(A-\lambda I)^{\dagger}$ and $(A-\lambda I)^{\dagger}(A-\lambda I)$ are orthogonal projectors onto $\mathcal{R}(A-\lambda I)$, we have that

$$
\begin{equation*}
g=(A-\lambda I)(A-\lambda I)^{\dagger} g=(A-\lambda I)^{\dagger}(A-\lambda I) g \tag{5.1}
\end{equation*}
$$

for any $\lambda$ such that $g \in \mathcal{R}(A-\lambda I)$. In particular, Lemmas 5.2 and 5.3 imply that (5.1) holds for $\lambda=\lambda_{k}$ and $\lambda=\lambda_{*}$.
5.2.2. Technical lemmas. We present several technical lemmas that allow us to prove our local convergence result. We will use the following notation:

$$
\begin{equation*}
A_{k} \equiv A-\lambda_{k} I \quad \text { and } \quad A_{*} \equiv A-\lambda_{*} I \tag{5.2}
\end{equation*}
$$

Step 2.1. Adjust $\alpha_{k}$.
Input: $\varepsilon_{\nu}, \varepsilon_{\alpha} \in(0,1), \alpha_{L}, \alpha_{U}, \alpha_{k}$ with $\alpha_{k} \in\left[\alpha_{L}, \alpha_{U}\right]$.
Output: $\alpha_{k},\left\{\lambda_{1}\left(\alpha_{k}\right),\left(\nu_{1}, u_{1}^{T}\right)^{T}\right\}$ and $\left\{\lambda_{i}\left(\alpha_{k}\right),\left(\nu_{2}, u_{2}^{T}\right)^{T}\right\}$.

```
Set \alpha= \alpha
if k>0 then
    compute eigenpairs {}{\mp@subsup{\lambda}{1}{}(\alpha),(\mp@subsup{\nu}{1}{},\mp@subsup{u}{1}{T}\mp@subsup{)}{}{T}}\mathrm{ and {}{\mp@subsup{\lambda}{i}{}(\alpha),(\mp@subsup{\nu}{2}{},\mp@subsup{u}{2}{T}\mp@subsup{)}{}{T}}
    corresponding to smallest eigenvalue and an eigenvalue
    close to second smallest eigenvalue of B\alpha
end if
while
```



```
    and |\mp@subsup{\alpha}{U}{}-\mp@subsup{\alpha}{L}{}|>\mp@subsup{\varepsilon}{\alpha}{}*\operatorname{max}{|\mp@subsup{\alpha}{L}{}|,|\mp@subsup{\alpha}{U}{}|} do
    \alpha
    \alpha=(\mp@subsup{\alpha}{L}{}+\mp@subsup{\alpha}{U}{\prime})/2
    Compute {\mp@subsup{\lambda}{1}{}(\alpha),(\mp@subsup{\nu}{1}{},\mp@subsup{u}{1}{T}\mp@subsup{)}{}{T}}\mathrm{ and {}{\mp@subsup{\lambda}{i}{}(\alpha),(\mp@subsup{\nu}{2}{},\mp@subsup{u}{2}{T}\mp@subsup{)}{}{T}}
end while
Set }\mp@subsup{\alpha}{k}{}=
```

Fig. 3. Adjustment of $\alpha$.

Step 2.5. Safeguard $\alpha_{k+1}$.
Input: $\alpha_{k+1}$ computed by step 2.4 of Algorithm 4.1, $\delta_{U} \geq \delta_{1}, \alpha_{L}, \alpha_{U}$, $\phi_{i}=-g^{T} x_{i}$, and $\phi_{i}^{\prime}=\left\|x_{i}\right\|^{2}$, for $i=k-1, k$.
Output: Safeguarded value for $\alpha_{k+1}$.
if $\alpha_{k+1} \notin\left[\alpha_{L}, \alpha_{U}\right]$
if $k=0$ then $\alpha_{k+1}=\delta_{U}+\phi_{k}+\phi_{k}^{\prime}\left(\delta_{U}-\lambda_{k}\right)$
else if $\left\|x_{k}\right\|<\left\|x_{k-1}\right\|$ then $\alpha_{k+1}=\delta_{U}+\phi_{k}+\phi_{k}^{\prime}\left(\delta_{U}-\lambda_{k}\right)$
else $\alpha_{k+1}=\delta_{U}+\phi_{k-1}+\phi_{k-1}^{\prime}\left(\delta_{U}-\lambda_{k-1}\right)$
end if
if $\alpha_{k+1} \notin\left[\alpha_{L}, \alpha_{U}\right]$ then set $\alpha_{k+1}=\left(\alpha_{L}+\alpha_{U}\right) / 2$ end if
end if

Fig. 4. Safeguarding of $\alpha$.

The first lemma establishes a key relationship satisfied by the iterates computed by Algorithm 4.1.

Lemma 5.4. Let $\lambda_{k}, x_{k}$ be the iterates at iteration $k$ of Algorithm 4.1. Then

$$
x_{k}=-\left(A-\lambda_{k} I\right)^{\dagger} g .
$$

Proof. First note that if $\lambda_{k}, x_{k}$ are the iterates at iteration $k$ of Algorithm 4.1, then they satisfy

$$
\left(\begin{array}{cc}
\alpha_{k} & g^{T} \\
g & A
\end{array}\right)\binom{1}{x_{k}}=\lambda_{k}\binom{1}{x_{k}} .
$$

Therefore

$$
\begin{equation*}
\left(A-\lambda_{k} I\right) x_{k}=-g \tag{5.3}
\end{equation*}
$$

In order to prove the result we need to consider two cases.
Case 1: $\lambda_{k} \neq \delta_{i}, i=1,2, \ldots, d$.
In this case we have that $A-\lambda_{k} I$ is nonsingular, $\left(A-\lambda_{k} I\right)^{-1}=\left(A-\lambda_{k} I\right)^{\dagger}$, and from (5.3) we conclude

$$
x_{k}=-\left(A-\lambda_{k} I\right)^{\dagger} g
$$

Case 2: $\lambda_{k}=\delta_{i}, 1 \leq i \leq d$.
If $\lambda_{k}=\delta_{i}$, then (5.3) implies that $g \perp \mathcal{S}_{i}$. This follows from the observation that for any $q \in \mathcal{S}_{i}$, we have $0=q^{T}\left(A-\delta_{i} I\right) x_{k}=-q^{T} g$. Corollary 3.1 now implies that $\alpha_{k}=\tilde{\alpha}_{i}$ and

$$
\begin{aligned}
x_{k} & =p_{i} \\
& =-\left(A-\delta_{i} I\right)^{\dagger} g,
\end{aligned}
$$

since $\left(1, x_{k}^{T}\right)^{T}$ is an eigenvector of $B_{\alpha_{k}}$. This concludes the proof.
Before presenting the next lemma, which provides useful relationships for the convergence analysis, we introduce the following definition.

Definition 5.1. Let $\lambda_{i}, x_{i}$ and $\lambda_{j}, x_{j}$ be the iterates computed by Algorithm 4.1 at iterations $i$ and $j$, respectively. Then we define

$$
\begin{equation*}
\rho(i, j) \equiv x_{i}^{T} A_{j}^{\dagger} x_{i}+x_{j}^{T} A_{i}^{\dagger} x_{j} \tag{5.4}
\end{equation*}
$$

We can substitute any of the iterates by $\lambda_{*}$, $y$, with $y=-A_{*}^{\dagger} g$. We denote this by $\rho(*, j)$ and $\rho(i, *)$, respectively.

Observe that if $A=Q D Q^{T}$ is an eigendecomposition of $A$, i.e., $Q$ is an orthogonal matrix and $D$ is a diagonal matrix with the eigenvalues of $A$ on the diagonal, we can write $\rho(i, j)$ in the following way:

$$
\rho(i, j)=g^{T} Q D_{i}^{\dagger}\left(D_{i}^{\dagger}+D_{j}^{\dagger}\right) D_{j}^{\dagger} Q^{T} g
$$

where $D_{i}=D-\lambda_{i} I$ and $D_{j}=D-\lambda_{j} I$. From this expression we obtain

$$
\begin{equation*}
\rho(i, j)=\sum_{k=1}^{d} \frac{\beta_{k}^{2}\left(2 \delta_{k}-\lambda_{i}-\lambda_{j}\right)}{\left(\delta_{k}-\lambda_{i}\right)^{2}\left(\delta_{k}-\lambda_{j}\right)^{2}}, \tag{5.5}
\end{equation*}
$$

where $\beta_{k}^{2}$ is the sum of the expansion coefficients of $g$ in the eigenvector basis, corresponding to all the eigenvectors associated with $\delta_{k}$. As before, we assume that $\delta_{1}, \delta_{2}, \ldots, \delta_{d}$ are the distinct eigenvalues of $A$ in nondecreasing order.

Lemma 5.5. Let $\lambda_{i}, x_{i}, \lambda_{j}, x_{j}$, and $\lambda_{k}, x_{k}$ be the iterates computed by Algorithm 4.1 at iterations $i, j$, and $k$, respectively. Then
(i) $\left(x_{i}-x_{j}\right)^{T} g=\left(\lambda_{j}-\lambda_{i}\right) x_{i}^{T} x_{j}$;
(ii) $\left(x_{i}-x_{j}\right)^{T} x_{k}=\left(\lambda_{i}-\lambda_{j}\right) x_{j}^{T} A_{i}^{\dagger} x_{k}$;
(iii) $x_{i}^{T} x_{i}-x_{j}^{T} x_{j}=\left(\lambda_{i}-\lambda_{j}\right) \rho(i, j)$, with $\rho(i, j)$ given by (5.4).

Moreover, (i)-(iii) also hold if we substitute any of the pairs above by $\lambda_{*}, y$, where $\lambda_{*}$ is the Lagrange multiplier corresponding to a boundary solution of problem (2.1) and $y=-A_{*}^{\dagger} g$.

Proof. Let us first prove (i). Observe that by Lemma 5.4

$$
\left(x_{i}-x_{j}\right)^{T} g=\left(A_{j}^{\dagger} g-A_{i}^{\dagger} g\right)^{T} g
$$

Therefore, using (5.1) and the fact that $A_{i}, A_{i}^{\dagger}, A_{j}, A_{j}^{\dagger}$ commute, we have

$$
\begin{aligned}
\left(x_{i}-x_{j}\right)^{T} g & =g^{T}\left(A_{j}^{\dagger}-A_{i}^{\dagger}\right) g \\
& =g^{T} A_{i}^{\dagger}\left(A_{i}-A_{j}\right) A_{j}^{\dagger} g \\
& =\left(\lambda_{j}-\lambda_{i}\right) x_{i}^{T} x_{j} .
\end{aligned}
$$

To prove (ii), we use (5.1), Lemma 5.4, and the fact that $A_{i}, A_{i}^{\dagger}, A_{j}, A_{j}^{\dagger}$ commute, obtaining

$$
\begin{aligned}
\left(x_{i}-x_{j}\right)^{T} x_{k} & =g^{T}\left(A_{i}^{\dagger}-A_{j}^{\dagger}\right) x_{k} \\
& =g^{T} A_{j}^{\dagger}\left(A_{j}-A_{i}\right) A_{i}^{\dagger} x_{k} \\
& =\left(\lambda_{i}-\lambda_{j}\right) x_{j}^{T} A_{i}^{\dagger} x_{k}
\end{aligned}
$$

Finally, let us prove (iii). By (5.1), Lemma 5.4, and the fact that $A_{i}, A_{i}^{\dagger}, A_{j}, A_{j}^{\dagger}$ commute, we have

$$
\begin{aligned}
x_{i}^{T} x_{i}-x_{j}^{T} x_{j} & =g^{T}\left[\left(A_{i}^{\dagger}\right)^{2}-\left(A_{j}^{\dagger}\right)^{2}\right] g \\
& =g^{T}\left[\left(A_{i}^{\dagger}\right)^{2} A_{j}^{2}-A_{i}^{2}\left(A_{j}^{\dagger}\right)^{2}\right] g \\
& =g^{T}\left(A_{i}^{\dagger}\right)^{2}\left(A_{j}-A_{i}\right)\left(A_{j}+A_{i}\right)\left(A_{j}^{\dagger}\right)^{2} g \\
& =\left(\lambda_{i}-\lambda_{j}\right) x_{i}^{T}\left(A_{i}^{\dagger}+A_{j}^{\dagger}\right) x_{j} \\
& =\left(\lambda_{i}-\lambda_{j}\right) \rho(i, j)
\end{aligned}
$$

Observe that (i)-(iii) hold for $\lambda_{*}, y$, since (5.1) holds for $\lambda_{*}, y=-A_{*}^{\dagger} g$, and $A_{*}$ commutes with the matrices above. This observation concludes the proof.

Using the updating formula (4.6), we obtained the following result relating $\lambda_{k+1}-$ $\lambda_{*}$ with $\lambda_{k-1}-\lambda_{*}$ and $\lambda_{k}-\lambda_{*}$. This lemma provides a key relationship for establishing the local convergence properties of Algorithm 4.1.

LEMMA 5.6. Let $\lambda_{*} \leq \delta_{1}$ be the Lagrange multiplier corresponding to a boundary solution of problem (2.1), with $g \neq 0$. Let $\lambda_{k+1}, x_{k+1}$ be the $(k+1)$ st iterates computed by Algorithm 4.1 using the two-point interpolating scheme given by (4.6) to update $\alpha$. Then, there exists a neighborhood $\mathcal{B}$ of $\lambda_{*}$ such that if $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}$, then $\lambda_{k+1}$ satisfies

$$
\begin{equation*}
\left|\lambda_{k+1}-\lambda_{*}\right| \leq \mathcal{C}\left|\lambda_{k-1}-\lambda_{*}\right|\left|\lambda_{k}-\lambda_{*}\right| \tag{5.6}
\end{equation*}
$$

with $\mathcal{C}$ independent of $k$.
Proof. Let $y=-A_{*}^{\dagger} g$ and $\alpha_{*}=\lambda_{*}-g^{T} y$. We divide the proof into two cases $\|y\|=\Delta$ and $\|y\|<\Delta$. In each case, we first find an appropriate neighborhood of $\lambda_{*}$ and then prove (5.6) for $\lambda_{k-1}, \lambda_{k}$ in that neighborhood.

Case 1: $\|y\|=\Delta$.
We will first find a neighborhood $\mathcal{B}$ of $\lambda_{*}$ such that if $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}$, then $\widehat{\lambda} \in \mathcal{B}$, with $\hat{\lambda}$ given by (4.4). In this part of the proof we define the numbers $\ell$ and $m$ in the following way.

Let $0 \leq \ell<n$ and assume that $g \perp \mathcal{S}_{i}, i=1,2, \ldots, \ell$, where $\ell=0$ indicates that $g \not \perp \mathcal{S}_{1}$. Let $m=0$ if $\lambda_{*}<\delta_{1}$ and $m=\ell$ if $\lambda_{*}=\delta_{1}$. Define

$$
r_{1}=\frac{\delta_{m+1}-\lambda_{*}}{2} \quad \text { and } \quad \mathcal{B}_{1}=\left\{\lambda| | \lambda-\lambda_{*} \mid \leq r_{1}\right\}
$$

and suppose that $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}_{1}$. Then by (4.4), Lemma 5.5 (iii), and the fact that $\|y\|=\Delta$, we have

$$
\begin{align*}
\hat{\lambda}-\lambda_{*} & =\lambda_{k}-\lambda_{*}+\frac{\left\|x_{k-1}\right\|\left(\left\|x_{k-1}\right\|+\left\|x_{k}\right\|\right)\left(\Delta-\left\|x_{k}\right\|\right)}{\rho(k-1, k) \Delta} \\
& =\left(\lambda_{k}-\lambda_{*}\right)\left[1-\frac{\rho(*, k)\left\|x_{k-1}\right\|\left(\left\|x_{k-1}\right\|+\left\|x_{k}\right\|\right)}{\rho(k-1, k) \Delta\left(\Delta+\left\|x_{k}\right\|\right)}\right] \tag{5.7}
\end{align*}
$$

We will prove now that $\left|\widehat{\lambda}-\lambda_{*}\right| \leq\left|\lambda_{k}-\lambda_{*}\right| \vartheta$, with $\vartheta \in(0,1)$.
Let $\Delta_{\text {max }}=\max _{\lambda \in \mathcal{B}_{1}}\left\|(A-\lambda I)^{\dagger} g\right\|$ and $\Delta_{\text {min }}=\min _{\lambda \in \mathcal{B}_{1}}\left\|(A-\lambda I)^{\dagger} g\right\|$. Therefore

$$
\begin{equation*}
\frac{\left\|x_{k-1}\right\|\left(\left\|x_{k-1}\right\|+\left\|x_{k}\right\|\right)}{\Delta\left(\Delta+\left\|x_{k}\right\|\right)} \geq\left(\frac{\Delta_{\min }}{\Delta_{\max }}\right)^{2} \tag{5.8}
\end{equation*}
$$

In view of (5.5) we have that for $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}_{1}$

$$
\rho(*, k) \geq \frac{\left(2 \delta_{\ell+1}-\lambda_{*}-\lambda_{k}\right)}{\left(\delta_{d}-\lambda_{*}\right)^{2}\left(\delta_{d}-\lambda_{k}\right)^{2}}\|g\|^{2}
$$

Since $\delta_{m+1} \leq \delta_{\ell+1}$ and since $\frac{-\delta_{m+1}+\lambda_{*}}{2} \leq \lambda_{*}-\lambda_{k} \leq \frac{\delta_{m+1}-\lambda_{*}}{2}$, we have

$$
\begin{align*}
\rho(*, k) & \geq \frac{\left(2 \delta_{m+1}-\lambda_{*}-\lambda_{k}\right)}{\left(\delta_{d}-\lambda_{*}\right)^{2}\left(\delta_{d}-\lambda_{k}\right)^{2}}\|g\|^{2} \\
& =\frac{2\left(\delta_{m+1}-\lambda_{*}\right)+\left(\lambda_{*}-\lambda_{k}\right)}{\left(\delta_{d}-\lambda_{*}\right)^{2}\left[\left(\delta_{d}-\lambda_{*}\right)+\left(\lambda_{*}-\lambda_{k}\right)\right]^{2}}\|g\|^{2} \\
& \geq \frac{2}{3}\|g\|^{2} \frac{\left(\delta_{m+1}-\lambda_{*}\right)}{\left(\delta_{d}-\lambda_{*}\right)^{4}} . \tag{5.9}
\end{align*}
$$

Using similar manipulations we obtain

$$
\begin{align*}
\rho(k-1, k) & \leq \frac{\left(2 \delta_{d}-\lambda_{k}-\lambda_{k-1}\right)\|g\|^{2}}{\left(\delta_{\ell+1}-\lambda_{k}\right)^{2}\left(\delta_{\ell+1}-\lambda_{k-1}\right)^{2}} \\
& \leq \frac{\left(2 \delta_{d}-\lambda_{k}-\lambda_{k-1}\right)\|g\|^{2}}{\left(\delta_{m+1}-\lambda_{k}\right)^{2}\left(\delta_{m+1}-\lambda_{k-1}\right)^{2}} \\
& \leq \frac{3 \cdot 2^{4}\|g\|^{2}\left(\delta_{d}-\lambda_{*}\right)}{\left(\delta_{m+1}-\lambda_{*}\right)^{4}} . \tag{5.10}
\end{align*}
$$

It follows from (5.7), (5.8), (5.9), and (5.10) that

$$
\left|\widehat{\lambda}-\lambda_{*}\right| \leq\left|\lambda_{k}-\lambda_{*}\right|\left|1-\left(\frac{\Delta_{\min }}{\Delta_{\max }}\right)^{2} \frac{1}{72} \frac{\left(\delta_{m+1}-\lambda_{*}\right)^{5}}{\left(\delta_{d}-\lambda_{*}\right)^{5}}\right| \equiv\left|\lambda_{k}-\lambda_{*}\right| \vartheta
$$

with $\vartheta \in(0,1)$. Therefore $\widehat{\lambda} \in \mathcal{B}_{1}$ whenever $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}_{1}$.
Now, we use these results to establish (5.6). Let the neighborhood $\mathcal{B}$ be given by $\mathcal{B}_{1}$ and let $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}$; therefore $\widehat{\lambda} \in \mathcal{B}$, and to prove (5.6) we need to consider two possibilities $\widehat{\lambda}<\delta_{1}$ and $\widehat{\lambda} \geq \delta_{1}$.

Case 1.1: $\|y\|=\Delta$ and $\widehat{\lambda}<\delta_{1}$.
In this case, we use formulas (4.4) and (4.6), obtaining

$$
\alpha_{k+1}=T_{1}+T_{2}
$$

where

$$
T_{1}=\frac{\alpha_{k-1}\left\|x_{k-1}\right\|\left(\left\|x_{k}\right\|-\Delta\right)+\left\|x_{k}\right\|\left(\Delta-\left\|x_{k-1}\right\|\right)}{\Delta\left(\left\|x_{k}\right\|-\left\|x_{k-1}\right\|\right)}
$$

and

$$
T_{2}=\frac{\left\|x_{k}\right\|\left\|x_{k-1}\right\|\left(\Delta-\left\|x_{k}\right\|\right)\left(\Delta-\left\|x_{k-1}\right\|\right)\left(\lambda_{k}-\lambda_{k-1}\right)}{\Delta\left(\left\|x_{k}\right\|-\left\|x_{k-1}\right\|\right)}
$$

We will now find an upper bound for $\left|T_{2}\right|$. From Lemma 5.5(iii), we have

$$
T_{2}=\frac{\left\|x_{k}\right\|\left\|x_{k-1}\right\|\left(\Delta-\left\|x_{k}\right\|\right)\left(\Delta-\left\|x_{k-1}\right\|\right)\left(\left\|x_{k}\right\|+\left\|x_{k-1}\right\|\right)}{\Delta \rho(k-1, k)}
$$

and since $\|y\|=\Delta$ we can write

$$
T_{2}=\frac{\left\|x_{k}\right\|\left\|x_{k-1}\right\|\left(\left\|x_{k}\right\|+\left\|x_{k-1}\right\|\right)}{\rho(k-1, k)} \frac{\left(\|y\|-\left\|x_{k}\right\|\right)\left(\|y\|-\left\|x_{k-1}\right\|\right)}{\Delta}
$$

Using Lemma 5.5 (iii) we obtain

$$
T_{2}=\frac{\left\|x_{k}\right\|\left\|x_{k-1}\right\|\left(\left\|x_{k}\right\|+\left\|x_{k-1}\right\|\right)}{\rho(k-1, k)} \frac{\left(\lambda_{*}-\lambda_{k-1}\right)\left(\lambda_{*}-\lambda_{k}\right) \rho(k-1, *) \rho(k, *)}{\Delta\left(\Delta+\left\|x_{k}\right\|\right)\left(\Delta+\left\|x_{k-1}\right\|\right)} .
$$

Now, since $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}$ and since $\delta_{m+1} \leq \delta_{\ell+1}$,

$$
\begin{align*}
\rho(k-1, k) & \geq \frac{2}{3}\|g\|^{2} \frac{\left(\delta_{m+1}-\lambda_{*}\right)}{\left(\delta_{d}-\lambda_{*}\right)^{4}}  \tag{5.11}\\
\rho(*, k), \rho(*, k-1) & \leq 10 \frac{\delta_{d}-\lambda_{*}}{\left(\delta_{m+1}-\lambda_{*}\right)^{4}} \tag{5.12}
\end{align*}
$$

Since for $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}$ we also have $\Delta_{\min } \leq\left\|x_{k-1}\right\|,\left\|x_{k}\right\| \leq \Delta_{\max }$, we obtain

$$
\begin{equation*}
\left|T_{2}\right| \leq \mathcal{C}_{2}\left|\lambda_{*}-\lambda_{k-1}\right|\left|\lambda_{*}-\lambda_{k}\right| \tag{5.13}
\end{equation*}
$$

We will use this estimate in a moment. First, we need to relate it to $\lambda_{k+1}-\lambda_{*}$. To do this, consider

$$
\begin{equation*}
\alpha_{k+1}-\alpha_{*}=\frac{\left(\alpha_{k-1}-\alpha_{*}\right)\left\|x_{k-1}\right\|\left(\left\|x_{k}\right\|-\Delta\right)+\left(\alpha_{k}-\alpha_{*}\right)\left\|x_{k}\right\|\left(\Delta-\left\|x_{k-1}\right\|\right)}{\Delta\left(\left\|x_{k}\right\|-\left\|x_{k-1}\right\|\right)}+T_{2} \tag{5.14}
\end{equation*}
$$

From Lemma 5.5(i), the definition of $\alpha_{*}$, and since $\alpha_{j}-\lambda_{j}=-g^{T} x_{j}$ for $j>0$, we have

$$
\begin{align*}
\alpha_{j}-\alpha_{*} & =\lambda_{j}-\lambda_{*}-g^{T}\left(x_{j}-y\right) \\
& =\left(\lambda_{j}-\lambda_{*}\right)\left(1+x_{j}^{T} y\right) \tag{5.15}
\end{align*}
$$

Using (5.15) along with Lemma 5.5, (5.14) becomes

$$
\begin{aligned}
\left(\lambda_{k+1}-\lambda_{*}\right)\left(1+x_{k+1}^{T} y\right)= & \frac{\left(\lambda_{k-1}-\lambda_{*}\right)\left\|x_{k-1}\right\|\left(\left\|x_{k}\right\|-\Delta\right)\left(1+x_{k-1}^{T} y\right)}{\Delta\left(\left\|x_{k}\right\|-\left\|x_{k-1}\right\|\right)} \\
& -\frac{\left(\lambda_{k}-\lambda_{*}\right)\left\|x_{k}\right\|\left(\left\|x_{k-1}\right\|-\Delta\right)\left(1+x_{k}^{T} y\right)}{\Delta\left(\left\|x_{k}\right\|-\left\|x_{k-1}\right\|\right)}+T_{2} \\
= & \frac{\left(\lambda_{k-1}-\lambda_{*}\right)\left(\lambda_{k}-\lambda_{*}\right) T_{3}}{\Delta\left(\left\|x_{k}\right\|-\left\|x_{k-1}\right\|\right)\left(\left\|x_{k}\right\|+\Delta\right)\left(\left\|x_{k-1}\right\|+\Delta\right)}+T_{2}
\end{aligned}
$$

where

$$
\begin{aligned}
T_{3}= & \left(\left\|x_{k-1}\right\|-\left\|x_{k}\right\|\right)\left(\left\|x_{k-1}\right\|+\left\|x_{k}\right\|+\Delta\right)\left(1+x_{k}^{T} y\right) \rho(k-1, *) \\
& +y^{T}\left(x_{k-1}-x_{k}\right)\left\|x_{k-1}\right\|\left(\left\|x_{k-1}\right\|+\Delta\right) \rho(k, *) \\
& +(\rho(k, *)-\rho(k-1, *))\left\|x_{k-1}\right\|\left(\left\|x_{k-1}\right\|+\Delta\right)\left(1+x_{k}^{T} y\right)
\end{aligned}
$$

Now, by Lemma 5.5(ii) we have

$$
\begin{equation*}
y^{T}\left(x_{k-1}-x_{k}\right)=\left(\lambda_{k-1}-\lambda_{k}\right) y^{T} A_{k}^{\dagger} x_{k-1} \tag{5.16}
\end{equation*}
$$

and by Lemma 5.5(iii)

$$
\begin{align*}
\rho(k, *)-\rho(k-1, *) & =x_{k}^{T} A_{*}^{\dagger} x_{k}+y^{T} A_{k}^{\dagger} y-x_{k-1}^{T} A_{*}^{\dagger} x_{k-1}-y A_{k-1}^{\dagger} y \\
& =g^{T} A_{k}^{\dagger} A_{*}^{\dagger} A_{k}^{\dagger} g-g^{T} A_{k-1}^{\dagger} A_{*}^{\dagger} A_{k-1}^{\dagger} g+y^{T}\left(A_{k}^{\dagger}-A_{k-1}^{\dagger}\right) y \\
& =g^{T}\left(\left(A_{k}^{\dagger}\right)^{2}-\left(A_{k-1}^{\dagger}\right)^{2}\right) y+\left(\lambda_{k}-\lambda_{k-1}\right) y^{T} A_{k}^{\dagger} A_{k-1}^{\dagger} y \\
& =\left(\lambda_{k}-\lambda_{k-1}\right)\left(y^{T} A_{k-1}^{\dagger} A_{k}^{\dagger} y+x_{k-1}^{T}\left(A_{k}^{\dagger}+A_{k-1}^{\dagger}\right) A_{k}^{\dagger} y\right) . \tag{5.17}
\end{align*}
$$

Therefore by (5.12), (5.16), (5.17), and since $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}$, we have

$$
\left|T_{3}\right| \leq \mathcal{C}_{3}\left|\lambda_{k}-\lambda_{k-1}\right|
$$

We may now combine the estimates we have established for $T_{1}, T_{2}$, and $T_{3}$ to give

$$
\begin{aligned}
\left|\lambda_{k+1}-\lambda_{*}\right|\left|1+x_{k+1}^{T} y\right| & \leq \mathcal{C}_{3} \frac{\left|\lambda_{k-1}-\lambda_{*}\right|\left|\lambda_{k}-\lambda_{*}\right|\left|\lambda_{k}-\lambda_{k-1}\right|}{\Delta\left|\left(\left\|x_{k}\right\|-\left\|x_{k-1}\right\|\right)\left(\left\|x_{k}\right\|+\Delta\right)\left(\left\|x_{k-1}\right\|+\Delta\right)\right|}+\left|T_{2}\right| \\
& \leq\left(\mathcal{C}_{4}+\mathcal{C}_{2}\right)\left|\lambda_{k-1}-\lambda_{*}\right|\left|\lambda_{k}-\lambda_{*}\right|
\end{aligned}
$$

since $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}$, and (5.13) holds. Let us see now that $\frac{1}{1+x_{k+1}^{T} y}<1$. Note that

$$
\begin{align*}
x_{k+1}^{T} y & =g^{T} A_{k+1}^{\dagger} A_{*}^{\dagger} g \\
& =\sum_{j=\ell+1}^{d} \frac{\beta_{j}^{2}}{\left(\delta_{j}-\lambda_{k+1}\right)\left(\delta_{j}-\lambda_{*}\right)} \geq \frac{\|g\|^{2}}{\left(\delta_{d}-\lambda_{k+1}\right)\left(\delta_{d}-\lambda_{*}\right)} . \tag{5.18}
\end{align*}
$$

From this expression we can conclude $x_{k+1}^{T} y>0$ since $\lambda_{*}<\delta_{m+1} \leq \delta_{d}$ and also since $\lambda_{k+1}<\delta_{m+1} \leq \delta_{d}$, by the way we compute the iterates in Algorithm 4.1.

We can now claim (5.6) when $\|y\|=\Delta$ and $\widehat{\lambda}<\delta_{1}$.
Case 1.2: $\|y\|=\Delta$ and $\widehat{\lambda} \geq \delta_{1}$.

In this case we must use the ideal safeguard, setting $\widehat{\lambda}=\delta_{1}$. Before proceeding with the proof, we point out that this can occur only when $\lambda_{*}=\delta_{1}$. To see this recall that, if $\lambda_{*}<\delta_{1}$, we have $\mathcal{B}=\left\{\lambda| | \lambda-\lambda_{*} \left\lvert\, \leq \frac{\delta_{1}-\lambda_{*}}{2}\right.\right\}$, and we proved that $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}$ implies $\widehat{\lambda} \in \mathcal{B}$. Therefore

$$
\begin{aligned}
\hat{\lambda}-\lambda_{*} & \leq \frac{\delta_{1}-\lambda_{*}}{2} \\
\widehat{\lambda} & \leq \frac{\delta_{1}+\lambda_{*}}{2}<\delta_{1}
\end{aligned}
$$

if $\lambda_{*}<\delta_{1}$.
To continue with the proof we write the formula for $\alpha_{k+1}$ in this case as

$$
\alpha_{k+1}=T_{4}+T_{5}
$$

where

$$
T_{4}=\omega \alpha_{k-1}+(1-\omega) \alpha_{k}
$$

and

$$
T_{5}=\frac{\left\|x_{k}\right\|\left\|x_{k-1}\right\|\left(\left\|x_{k}\right\|-\left\|x_{k-1}\right\|\right)}{\omega\left\|x_{k}\right\|+(1-\omega)\left\|x_{k-1}\right\|} \frac{\left(\lambda_{k-1}-\delta_{1}\right)\left(\lambda_{k}-\delta_{1}\right)}{\left(\lambda_{k}-\lambda_{k-1}\right)}
$$

Since $\widehat{\lambda}=\delta_{1}$, we have $\omega=\frac{\lambda_{k}-\delta_{1}}{\lambda_{k}-\lambda_{k-1}}$ and therefore

$$
\begin{align*}
\omega\left\|x_{k}\right\|+(1-\omega)\left\|x_{k-1}\right\| & =\frac{\left(\lambda_{k}-\delta_{1}\right)\left\|x_{k}\right\|+\left(\delta_{1}-\lambda_{k-1}\right)\left\|x_{k-1}\right\|}{\lambda_{k}-\lambda_{k-1}} \\
& =\frac{\left(\lambda_{k}-\delta_{1}\right) \rho(k-1, k)+\left\|x_{k-1}\right\|\left(\left\|x_{k-1}\right\|+\left\|x_{k}\right\|\right)}{\left\|x_{k-1}\right\|+\left\|x_{k}\right\|} \tag{5.19}
\end{align*}
$$

by Lemma $5.5(\mathrm{iii})$.
Using (5.19), (5.4), and Lemma 5.5(iii), we obtain

$$
T_{5}=\frac{\left\|x_{k}\right\|\left\|x_{k-1}\right\| \rho(k, k-1)\left(\lambda_{k-1}-\delta_{1}\right)\left(\lambda_{k}-\delta_{1}\right)}{\left(\lambda_{k}-\delta_{1}\right) \rho(k-1, k)+\left\|x_{k-1}\right\|\left(\left\|x_{k-1}\right\|+\left\|x_{k}\right\|\right)}
$$

By (5.11) and the hypothesis that $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}$, we have

$$
\begin{equation*}
\left|T_{5}\right| \leq \mathcal{C}_{5}\left|\lambda_{k-1}-\delta_{1}\right|\left|\lambda_{k}-\delta_{1}\right| \tag{5.20}
\end{equation*}
$$

We now write

$$
\begin{aligned}
\alpha_{k+1}-\alpha_{*} & =\omega \alpha_{k-1}+(1-\omega) \alpha_{k}-\alpha_{*}+T_{5} \\
& =\omega\left(\alpha_{k-1}-\alpha_{*}\right)+(1-\omega)\left(\alpha_{k}-\alpha_{*}\right)+T_{5}
\end{aligned}
$$

Equation (5.14) and the fact that $\lambda_{*}=\delta_{1}$ yield

$$
\left|\lambda_{k+1}-\lambda_{*}\right|\left|1+y^{T} x_{k+1}\right| \leq \frac{\left|\lambda_{k-1}-\lambda_{*}\right|\left|\lambda_{k}-\lambda_{*}\right|}{\left|\lambda_{k}-\lambda_{k-1}\right|}\left|y^{T}\left(x_{k}-x_{k-1}\right)\right|+\left|T_{5}\right|
$$

Observe that $\left|y^{T}\left(x_{k}-x_{k-1}\right)\right| \leq\left|g^{T}\left(A-\lambda_{k} I\right)^{\dagger}\left(A-\lambda_{k-1} I\right)^{\dagger}\left(A-\lambda_{*} I\right)^{\dagger} g\right|$, and we can compute an upper bound for this term using the Cauchy-Schwarz inequality, continuity of $\|\cdot\|$, and that $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}$. Therefore, by (5.20) and since $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}$

$$
\left|\lambda_{k+1}-\lambda_{*}\right|\left|1+y^{T} x_{k+1}\right| \leq\left(\mathcal{C}_{6}+\mathcal{C}_{5}\right)\left|\lambda_{k-1}-\lambda_{*}\right|\left|\lambda_{k}-\lambda_{*}\right| .
$$

Using (5.18), we can now establish (5.6) when $\|y\|=\Delta$ and $\hat{\lambda} \geq \delta_{1}$.
Case 2: $\|y\|<\Delta$.
In this situation, we are in the hard case and therefore $\lambda_{*}=\delta_{1}$ and $g \perp \mathcal{S}_{i}$, $i=1,2, \ldots, \ell$, with $1 \leq \ell<d$. For this case we will find a neighborhood $\mathcal{B}$ of $\lambda_{*}$ such that $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}$ implies $\hat{\lambda}>\delta_{1}$.

Let the function $\varphi(\lambda) \equiv\left\|(A-\lambda I)^{\dagger} g\right\|$. Then $\varphi(\lambda)$ is strictly increasing in $\left(-\infty, \delta_{\ell+1}\right)$, and there exist $\lambda_{a}, \lambda_{b}$ such that $\varphi\left(\lambda_{a}\right)=\frac{\Delta_{y}}{2}$ and $\varphi\left(\lambda_{b}\right)=\frac{\Delta+\Delta_{y}}{2}$, with $\Delta_{y}=\|y\|$.

Let

$$
r_{2}=\min \left\{\frac{\delta_{1}-\lambda_{a}}{2}, \frac{\delta_{1}-\lambda_{b}}{2}, \frac{\delta_{\ell+1}-\delta_{1}}{2}\right\} \quad \text { and } \quad \mathcal{B}_{2}=\left\{\lambda| | \lambda-\lambda_{*} \mid \leq r_{2}\right\} .
$$

Then for $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}_{2}$

$$
\begin{gathered}
\frac{\Delta_{y}}{2} \leq\left\|x_{k-1}\right\|,\left\|x_{k}\right\| \leq \frac{\Delta+\Delta_{y}}{2} \\
\text { and } \quad \delta_{1}-r_{2} \leq \lambda_{k-1}, \lambda_{k} \leq \delta_{1}+r_{2}<\delta_{\ell+1} .
\end{gathered}
$$

Now observe that using (4.4) and Lemma 5.5 (iii) we can write

$$
\begin{aligned}
\widehat{\lambda}-\lambda_{*} & =\widehat{\lambda}-\delta_{1} \\
& =\frac{\left(\lambda_{k-1}-\delta_{1}\right)\left\|x_{k-1}\right\|\left(\left\|x_{k}\right\|-\Delta\right)+\left(\lambda_{k}-\delta_{1}\right)\left\|x_{k}\right\|\left(\Delta-\left\|x_{k-1}\right\|\right)}{\Delta\left(\left\|x_{k}\right\|-\left\|x_{k-1}\right\|\right)}-\delta_{1} \\
& =\lambda_{k}-\delta_{1}+\frac{\left\|x_{k-1}\right\|\left(\Delta-\left\|x_{k}\right\|\right)\left(\left\|x_{k-1}\right\|+\left\|x_{k}\right\|\right)}{\Delta \rho(k, k-1)} \\
& \geq \lambda_{k}-\delta_{1}+\frac{\Delta_{y}^{2}\left(\Delta-\Delta_{y}\right)}{4 \Delta \rho(k, k-1)} .
\end{aligned}
$$

Observe now that for $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}_{2}$

$$
\begin{equation*}
\rho(k, k-1) \leq \frac{3 \cdot 2^{4}\|g\|^{2}\left(\delta_{d}-\delta_{1}\right)}{\left(\delta_{\ell+1}-\delta_{1}\right)^{4}} . \tag{5.22}
\end{equation*}
$$

Using (5.21) and (5.22) we obtain

$$
\widehat{\lambda}-\delta_{1} \geq \lambda_{k}-\delta_{1}+\frac{\Delta_{y}{ }^{2}\left(\Delta-\Delta_{y}\right)\left(\delta_{\ell+1}-\delta_{1}\right)^{4}}{3 \cdot 2^{6} \Delta\|g\|^{2}\left(\delta_{d}-\delta_{1}\right)} \text {. }
$$

Let

$$
\zeta \equiv \frac{\Delta_{y}{ }^{2}\left(\Delta-\Delta_{y}\right)\left(\delta_{\ell+1}-\delta_{1}\right)^{4}}{3 \cdot 2^{6} \Delta\|g\|^{2}\left(\delta_{d}-\delta_{1}\right)} .
$$

Observe that $\zeta$ is well defined since $\delta_{d}=\delta_{1}$ would imply $g \in \mathbb{R}^{n} \perp=\{0\}$. Observe also that $\zeta>0$, since $\Delta>\Delta_{y}$ and $\delta_{\ell+1}>\delta_{1}$. Then, for $\lambda_{k} \geq \delta_{1}-\zeta$ we have $\widehat{\lambda} \geq \delta_{1}$. So, let

$$
r_{3}=\min \left\{r_{2}, \zeta\right\} \quad \text { and } \quad \mathcal{B}_{3}=\left\{\lambda| | \lambda-\delta_{1} \mid \leq r_{3}\right\}
$$

It follows that for $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}_{3}$, we have $\widehat{\lambda} \geq \delta_{1}$ and we must use the ideal safeguard, setting $\widehat{\lambda}=\delta_{1}$. The proof now proceeds as in Case 1.2, where the neighborhood $\mathcal{B}$ is given by $\mathcal{B}_{3}$, and we use (5.22) instead of (5.10).

The analysis of the two cases concludes the proof.
Note that the assumption in Lemma 5.6 that the trust-region constraint is binding at the solution includes the possibility of the hard case, since in this case $x_{*}=$ $-A_{*}^{\dagger} g+z$, with $z \in \mathcal{S}_{1}$ and $\left\|x_{*}\right\|=\Delta$.

### 5.2.3. Local convergence result.

THEOREM 5.1. Let $\lambda_{*} \leq \delta_{1}$ be the Lagrange multiplier corresponding to a boundary solution of problem (2.1), with $g \neq 0$. Let $\left\{\lambda_{k}\right\},\left\{x_{k}\right\}$ be the sequences of iterates generated by Algorithm 4.1 using the two-point interpolating scheme given by (4.6) to update $\alpha$. There exists a neighborhood $\mathcal{B}$ of $\lambda_{*}$ such that if $\lambda_{i-1}, \lambda_{i} \in \mathcal{B}$, then for $k \geq i-1$
(i) $\left\{\lambda_{k}\right\}$ remains in $\mathcal{B}$ and converges $q$-superlinearly to $\lambda_{*}$;
(ii) $\left\{x_{k}\right\}$ converges $q$-superlinearly to $y=-\left(A-\lambda_{*} I\right)^{\dagger} g$.

Proof. First we show that $\left\{\lambda_{k}\right\}$ converges to $\lambda_{*}$ and that the rate of convergence is superlinear.

Let $r \in \mathbb{R}, r>0$, and $\mathcal{B}=\left\{\lambda| | \lambda-\lambda_{*} \mid<r\right\}$ be the neighborhood of $\lambda_{*}$ stated in Lemma 5.6 and suppose that $\lambda_{i-1}, \lambda_{i} \in \mathcal{B}$, for $i \geq 1$. Then, Lemma 5.6 implies that there exists $\mathcal{C}$ such that

$$
\begin{equation*}
\left|\lambda_{i+1}-\lambda_{*}\right| \leq \mathcal{C}\left|\lambda_{i}-\lambda_{*}\right|\left|\lambda_{i-1}-\lambda_{*}\right| \tag{5.23}
\end{equation*}
$$

Let $\hat{r}=\min \left\{r, \frac{1}{2 C}\right\}$, define $\widehat{\mathcal{B}}=\left\{\lambda| | \lambda-\lambda_{*} \mid<\hat{r}\right\}$, and observe that $\widehat{\mathcal{B}} \subset \mathcal{B}$. Suppose $\lambda_{i-1}, \lambda_{i} \in \widehat{\mathcal{B}}$; then $\lambda_{i-1}, \lambda_{i} \in \mathcal{B}$, and (5.23) holds.

Observe now that for $\lambda_{i-1}, \quad \lambda_{i} \in \widehat{\mathcal{B}}$ we have $\mathcal{C}\left|\lambda_{i-1}-\lambda_{*}\right| \leq \frac{1}{2}$ and therefore

$$
\left|\lambda_{i+1}-\lambda_{*}\right| \leq \frac{1}{2}\left|\lambda_{i}-\lambda_{*}\right|
$$

which implies $\lambda_{i+1} \in \widehat{\mathcal{B}} \subset \mathcal{B}$.
It follows inductively that if $\lambda_{i-1}, \lambda_{i} \in \mathcal{B}$, then $\lambda_{k} \in \mathcal{B}$ for $k \geq i-1$ and this implies

$$
\left|\lambda_{k}-\lambda_{*}\right| \leq \frac{1}{2^{k-i+1}}\left|\lambda_{i-1}-\lambda_{*}\right|
$$

and therefore $\lambda_{k} \rightarrow \lambda_{*}$ as $k \rightarrow \infty$.
To see that the rate of convergence is $q$-superlinear, observe that, by (5.23), for $k \geq i$ we have

$$
\frac{\left|\lambda_{k+1}-\lambda_{*}\right|}{\left|\lambda_{k}-\lambda_{*}\right|}=\mathcal{C}\left|\lambda_{k-1}-\lambda_{*}\right|
$$

which goes to zero as $k$ goes to infinity.

In the second part of the proof we show that the sequence $\left\{x_{k}\right\}$ converges superlinearly to $y=-\left(A-\lambda_{*} I\right)^{\dagger} g$.

Recall from Lemma 5.4 that $x_{k}=-A_{k}^{\dagger} g$ and let us study $x_{k}-y$ which is given by

$$
\begin{aligned}
x_{k}-y & =\left(A-\lambda_{*} I\right)^{\dagger} g-\left(A-\lambda_{k} I\right)^{\dagger} g \\
& =\left(A-\lambda_{*} I\right)^{\dagger}\left(\left(A-\lambda_{k} I\right)-\left(A-\lambda_{*} I\right)\right)\left(A-\lambda_{k} I\right)^{\dagger} g \\
& =\left(\lambda_{*}-\lambda_{k}\right)\left(A-\lambda_{*} I\right)^{\dagger}\left(A-\lambda_{k} I\right)^{\dagger} g
\end{aligned}
$$

using (5.1) and rearranging terms. Taking norms on both sides we have

$$
\begin{align*}
\left\|x_{k}-y\right\| & =\left|\lambda_{k}-\lambda_{*}\right|\left\|\left(A-\lambda_{*} I\right)^{\dagger}\right\|\left\|\left(A-\lambda_{k} I\right)^{\dagger}\right\|\|g\| \\
& \leq \widehat{\mathcal{C}}\left|\lambda_{k}-\lambda_{*}\right| \tag{5.24}
\end{align*}
$$

for a positive constant $\widehat{\mathcal{C}}$, since $\lambda_{k} \in \mathcal{B}, \lambda_{*} \leq \delta_{1}$, and $\|g\|$ is constant. Therefore, since $\lambda_{k} \rightarrow \lambda_{*}$ as $k \rightarrow \infty$, we have that $x_{k} \rightarrow y$ as $k \rightarrow \infty$.

To see that the rate of convergence is q-superlinear, observe that (5.23) and (5.24) imply

$$
\frac{\left\|x_{k+1}-y\right\|}{\left\|x_{k}-y\right\|} \leq\left|\lambda_{k-1}-\lambda_{*}\right|
$$

which goes to zero as $k$ goes to infinity. This completes the proof.
5.2.4. Near hard case. The next lemma provides a relationship between the function $\phi$ and the interpolating function (4.5). We will use this relationship in the analysis of the near hard case.

Lemma 5.7. At iteration $k$ of Algorithm 4.1 the interpolating function (4.5) satisfies

$$
\widehat{\phi}(\widehat{\lambda})-\phi\left(\lambda_{k}\right)=\left(\widehat{\lambda}-\lambda_{k}\right)\left[x_{k}^{T} x_{k-1}+\frac{\left\|x_{k-1}\right\|\left\|x_{k}\right\| \rho(k, k-1)\left(\widehat{\lambda}-\lambda_{k-1}\right)}{\left(\lambda_{k}-\widehat{\lambda}\right) \rho(k, k-1)+\left\|x_{k-1}\right\|\left(\left\|x_{k}\right\|+\left\|x_{k-1}\right\|\right)}\right]
$$

with $\rho(k, k-1)$ as in (5.4) and $\widehat{\lambda}$ given by (4.4).
Proof. By (4.5) and Lemma 5.5(iii), we have

$$
\begin{aligned}
\widehat{\phi}(\widehat{\lambda})= & \frac{\gamma^{2}}{\delta-\widehat{\lambda}}+\omega \phi\left(\lambda_{k-1}\right)-\omega \frac{\gamma^{2}}{\delta-\lambda_{k-1}} \\
& +(1-\omega) \phi\left(\lambda_{k}\right)-(1-\omega) \frac{\gamma^{2}}{\delta-\lambda_{k}} \\
= & \phi\left(\lambda_{k}\right)+\omega g^{T}\left(x_{k}-x_{k-1}\right)+\frac{\omega \gamma^{2}\left(\widehat{\lambda}-\lambda_{k-1}\right)}{(\delta-\widehat{\lambda})\left(\delta-\lambda_{k-1}\right)} \\
& +\frac{(1-\omega) \gamma^{2}\left(\widehat{\lambda}-\lambda_{k}\right)}{(\delta-\widehat{\lambda})\left(\delta-\lambda_{k}\right)} \\
= & \phi\left(\lambda_{k}\right)+\left(\widehat{\lambda}-\lambda_{k}\right) x_{k}^{T} x_{k-1}+\frac{\gamma^{2}\left(\widehat{\lambda}-\lambda_{k}\right)\left(\widehat{\lambda}-\lambda_{k-1}\right)}{(\delta-\widehat{\lambda})\left(\delta-\lambda_{k}\right)\left(\delta-\lambda_{k-1}\right)} \\
= & \phi\left(\lambda_{k}\right)+\left(\widehat{\lambda}-\lambda_{k}\right) x_{k}^{T} x_{k-1} \\
& +\frac{\left\|x_{k-1}\right\|\left\|x_{k}\right\|\left(\left\|x_{k}\right\|-\left\|x_{k-1}\right\|\right)\left(\widehat{\lambda}-\lambda_{k}\right)\left(\widehat{\lambda}-\lambda_{k-1}\right)}{\left(\omega\left\|x_{k}\right\|+(1-\omega)\left\|x_{k-1}\right\|\right)\left(\lambda_{k}-\lambda_{k-1}\right)}
\end{aligned}
$$

where $\omega=\frac{\lambda_{k-1}-\widehat{\lambda}}{\lambda_{k}-\lambda_{k-1}}$. Thus, the result follows from Lemma 5.5.

A few comments are in order concerning the near hard case. As mentioned in section 3, finding $\lambda_{*}<\delta_{1}$ in a near hard case is a very ill-conditioned process. The difference $\delta_{1}-\lambda_{*}$ can be very small to the extent of being undetectable within the given tolerances. The smaller the value $\delta_{1}-\lambda_{*}$, the harder it is to determine $\left\{\lambda_{*}, x_{*}\right\}$. Furthermore, rounding errors generally will convert an exact hard case into a near hard case. Although $\delta_{1}$ is still a pole of $\phi$ when $g$ is not exactly orthogonal to $\mathcal{S}_{1}$, the weight of such a pole is very small in comparison to the other poles because the expansion coefficients of $g$ in the basis of eigenvectors of $A$ are practically zero for those eigenvectors associated with $\delta_{1}$. The strategy that we follow in Algorithm 4.1 for dealing with this case consists of building an interpolating function that ignores the pole $\delta_{1}$ at early stages, using the eigenpair corresponding to the second smallest eigenvalue of $B_{\alpha_{k}}$ to obtain the interpolation points. In addition, we use the second eigenpair to compute a vector that might be a quasi-optimal solution for the trustregion subproblem as established in Theorem 3.2. Moreover, as that theorem and related results established, it is not necessary to compute an eigenpair corresponding to the second smallest eigenvalue. This is especially useful when the vector $g$ is orthogonal or nearly orthogonal to several eigenspaces corresponding to the smallest eigenvalues of $A$, and those eigenvalues are clustered.

If we use information concerning a second eigenpair, then we will have $\lambda_{k}>\delta_{1}$. This occurs because the first component $\nu_{1}$ of the eigenvector $\left(\nu_{1}, u_{2}^{T}\right)^{T}$ associated with $\lambda_{1}\left(\alpha_{k}\right)$ is too small so that $\left\|u_{1} / \nu_{1}\right\|=\left\|x_{k}\right\|$ becomes excessively large. Therefore $\left\{\lambda_{k}, x_{k}\right\}$ is defined as $\left\{\lambda_{i}\left(\alpha_{k}\right), u_{2} / \nu_{2}\right\}$. Intuitively, this is a good strategy since in the exact hard case this would continuously select the correct eigenvector that will approach $\left(1, p_{1}^{T}\right)^{T}$ when $\alpha$ tends to the value $\tilde{\alpha}_{1}$, stated in Lemma 3.2, from either side.

Now, at iteration $k$ the parameter $\alpha_{k}$ is updated as $\alpha_{k+1}=\widehat{\lambda}+\widehat{\phi}(\widehat{\lambda})$ with $\widehat{\lambda} \leq \delta_{U}$, where either $\widehat{\lambda}<\delta_{U}, \widehat{\phi^{\prime}}(\widehat{\lambda})=\Delta^{2}$, or $\widehat{\lambda}=\delta_{U}, \widehat{\phi}^{\prime}\left(\delta_{U}\right)<\Delta^{2}$. By the same arguments of the proof of Case 1 in Lemma 5.6, there exists a neighborhood $\mathcal{B}$ of $\lambda_{*}$ such that if $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}$, then $\widehat{\lambda} \in \mathcal{B}$, with $\left|\hat{\lambda}-\lambda_{k}\right|=\vartheta\left|\lambda_{*}-\lambda_{k}\right|$, for $\vartheta \in(0,1)$. In other words, eventually the safeguarding $\widehat{\lambda}=\delta_{U}$ is no longer necessary. If $\lambda_{k-1}, \lambda_{k} \in \mathcal{B}$, then Lemma 5.7 implies that $\left|\widehat{\phi}(\widehat{\lambda})-\phi\left(\lambda_{k}\right)\right|=\vartheta\left|\widehat{\lambda}-\lambda_{k}\right|\left|\lambda_{*}-\lambda_{k}\right|$. The agreement between $\widehat{\lambda}$ and $\lambda_{k}$ and between $\widehat{\phi}(\widehat{\lambda})$ and $\phi\left(\lambda_{k}\right)$ drive $\alpha_{k}$ toward $\alpha_{*}=\lambda_{*}+\phi\left(\lambda_{*}\right)$. As $\alpha_{k}$ approaches $\alpha_{*}$, the reduction of the safeguarding interval $\left[\alpha_{L}, \alpha_{U}\right]$ at every iteration provides a means to avoid the numerical difficulties associated with a near hard case, and eventually there is no need to use a second eigenpair of $B_{\alpha_{k}}$. At early stages, however, it might be that $\widehat{\lambda}=\delta_{U}$. Although $\phi\left(\delta_{1}\right)$ is infinite, the interpolating function value $\widehat{\phi}\left(\delta_{U}\right)$ is finite. Using $\alpha_{k+1}=\delta_{U}+\widehat{\phi}\left(\delta_{U}\right)$ is essential in keeping the process under control.

### 5.3. Global convergence.

THEOREM 5.2. Algorithm 4.1 is globally convergent.
Proof. The goal of Algorithm 4.1 is to solve the trust-region subproblem either by determining the existence of an interior solution or by computing an optimal value $\alpha_{*}$ for the parameter $\alpha$, such that the solution to the parameterized eigenvalue problem for $B_{\alpha_{*}}$ can be used to compute a boundary solution for the trust-region subproblem. The global convergence of Algorithm 4.1 is achieved by keeping $\alpha_{k}$ in an interval that contains the optimal parameter $\alpha_{*}$.

We first recall that the initial safeguarding interval $\left[\alpha_{L}, \alpha_{U}\right]$ contains the optimal value $\alpha_{*}$. Starting with that interval, the updating procedure for $\alpha_{L}$ and $\alpha_{U}$ guaran-
tees that $\alpha_{*}$ remains in the interval and that the safeguarding interval is reduced at each iteration.

Therefore, since $\alpha_{k}=\lambda_{k}-g^{T} x_{k}$, after a finite number of iterations either the sequence $\left\{\lambda_{k}\right\}$ reaches the neighborhood of $\lambda_{*}$ of Theorem 5.1 that guarantees convergence, or the length of the safeguarding interval $\left|\alpha_{U}-\alpha_{L}\right|$ goes to zero with $\alpha_{L} \leq \alpha_{*} \leq \alpha_{U}$.
6. Numerical experiments. In this section we present numerical experiments to demonstrate the viability of our approach and to illustrate different aspects of our method. We implemented Algorithm 4.1 (LSTRS) in MATLAB 5.3 using a Mexfile interface to access the IRLM [16] implemented in ARPACK [7]. We ran our experiments on a Sun Ultrasparc 10 with a 300 MHz processor and 256 megabytes of RAM, running Solaris 5.6. The floating point arithmetic was IEEE standard double precision with machine precision $2^{-52} \approx 2.2204 \cdot 10^{-16}$. We present five sets of experiments. In the first and second sets we study the sensitivity of LSTRS to different tolerances for the trust-region radius and to different sizes of the trust-region radius, respectively, for problems where the hard case is not present. In order to put our method in context, we include the number of matrix-vector products required by the conjugate gradient method to solve systems of the form $(A-\lambda I) x=-g$. The third set of experiments illustrates the local superlinear rate of convergence. The fourth set shows the behavior of LSTRS in the hard case. In the fifth set we provide a comparison with the semidefinite programming approach presented in [13].

The following tolerances are fixed in all the experiments: $\varepsilon_{\nu}=10^{-2}, \varepsilon_{\alpha}=$ $10^{-8}, \varepsilon_{I n t}=10^{-8}$. We will indicate the values for the rest of the parameters when we describe each particular set of experiments.
6.1. Different tolerances. In the first experiment, we show the behavior of the method when different levels of accuracies of the norm of the trust-region solution are required. The matrix $A$ in (2.1) was $A=L-5 I$, where $L$ is the standard two-dimensional (2D) discrete Laplacian on the unit square based upon a 5 -point stencil with equally spaced mesh points. The shift of $-5 I$ was introduced to make $A$ indefinite. The order of $A$ was $n=1024$. We solved a sequence of 20 related problems, differing only by the vector $g$, randomly generated with entries uniformly distributed on $(0,1)$. We solved each of these problems for a fixed trust-region radius $\Delta=100$ and for $\varepsilon_{\Delta}=10^{-4}, 10^{-6}, 10^{-8}$, where $\varepsilon_{\Delta}$ is the relative accuracy of the norm of the computed solution with respect to $\Delta$. The initial $\delta_{U}$ was the minimum of the diagonal of $A$ and $\alpha_{0}=\delta_{U}$. The tolerance for a quasi-optimal solution was set to $\varepsilon_{H C}=10^{-16}$ in order to allow the method to compute a boundary solution; otherwise the quasi-optimal stopping criterion would be satisfied first.

For $\varepsilon_{\Delta}=10^{-4}, 10^{-6}$ the number of Lanczos basis vectors was limited to 9 , and 6 shifts (i.e., 6 matrix-vector products) were applied on each implicit restart, while for $\varepsilon_{\Delta}=10^{-8}$, the number of vectors was 20 with 14 shifts on each implicit restart. The maximum number of restarts allowed was 45 for $\varepsilon_{\Delta}=10^{-4}, 10^{-6}$ and 100 for $\varepsilon_{\Delta}=10^{-8}$. More basis vectors were needed for $\varepsilon_{\Delta}=10^{-8}$, since in this case the eigenvalues were computed to a higher accuracy. We chose $v_{1}$, the initial vector for the IRLM, in the following way. In the first iteration of LSTRS, $v_{1}=(1,1, \ldots, 1) / \sqrt{n+1}$ and subsequently, $v_{1}$ was the first column of the matrix $V$ containing the Lanczos vectors computed by the IRLM for the previous bordered matrix. This choice standardized the initial vector along the set of tests and performed better than a randomly generated vector, or the eigenvector corresponding to the smallest eigenvalue of $B_{\alpha_{k}}$, or the vector $\left(0, g^{T}\right)^{T}$. Note that the

Table 1
Average behavior for different tolerances.

| $\varepsilon_{\Delta}$ | LSTRS IT | LSTRS MV | CG MV | $\frac{L S T R S M V}{C G M V}$ |
| :---: | :---: | :---: | :---: | :---: |
| $10^{-4}$ | 6.00 | 64.00 | 43.45 | 1.47 |
| $10^{-6}$ | 7.50 | 83.00 | 57.00 | 1.46 |
| $10^{-8}$ | 6.55 | 183.40 | 71.55 | 2.56 |

Table 2
Average behavior for different tolerances allowing quasi-optimal solutions.

| $\varepsilon_{\Delta}$ | LSTRS IT | LSTRS MV | CG MV | $\frac{L S T R S M V}{C G M V}$ |
| :---: | :---: | :---: | :---: | :---: |
| $10^{-4}$ | 5.00 | 54.00 | 48.80 | 1.11 |
| $10^{-6}$ | 5.40 | 62.00 | 62.90 | 0.99 |
| $10^{-8}$ | 5.40 | 64.75 | 76.95 | 0.84 |

last two options have the additional disadvantage of preventing the IRLM from finding the eigenspace of $B_{\alpha}$ corresponding to $\delta_{1}$ whenever a potential hard case is present. As in [17], we relaxed the accuracy required in the eigenvalue solution and made it proportional to the relative accuracy in the computed solution. Specifically, $\left\|B_{\alpha} q-q \lambda\right\|<\varepsilon_{L a n}$, where $\varepsilon_{\text {Lan }}=\max \left\{\min \left\{\varepsilon_{L a n},\left|\frac{\Delta-\|x\|}{\Delta}\right|\right\}, \varepsilon_{\max }\right\}$ and $\varepsilon_{\text {max }}=0.125,0.1,0.075$ for $\varepsilon_{\Delta}=10^{-4}, 10^{-6}, 10^{-8}$, respectively.

In Table 1 we report the average number of iterations of LSTRS (LSTRS IT), the average number of matrix-vector products required by LSTRS (LSTRS MV), and the average number of matrix-vector products required by the conjugate gradient method (CG MV) to solve the system $\left(A-\lambda_{*} I\right) x=-g$ to the same accuracy $\varepsilon_{\Delta}$ in the norm of the computed solution of LSTRS. The value of $\lambda_{*}$ was the optimal value computed by LSTRS.

We observe that for $\varepsilon_{\Delta}=10^{-4}, 10^{-6}$ the behavior in [17] is reproduced: a trustregion solution requires fewer than twice as many matrix-vectors products on average than the number needed to solve a single linear system to the same accuracy using conjugate gradients. For $\varepsilon_{\Delta}=10^{-8}$, even though LSTRS requires more matrix-vector products, the cost of LSTRS is less than three times the cost of solving one system by conjugate gradients.

If we repeat the experiment, setting the tolerance for a quasi-optimal solution to $\varepsilon_{H C}=10^{-6}$, we obtain the results in Table 2, where we observe the low number of matrix-vector products required by LSTRS. In this experiment we used nine Lanczos basis vectors for all cases and allowed a maximum of 45 restarts.
6.2. Different trust-region radii. The second experiment illustrates the behavior of LSTRS for different sizes of the trust-region radius. The matrix $A$ in (2.1) was of the form $A=U D U^{T}$ with $D$ diagonal and $U=I-2 u u^{T}, u^{T} u=1$. The elements of $D$ were randomly selected from a uniform distribution on $(-5,5)$. Both vectors $u$ and $g$ were randomly generated with entries uniformly distributed on $(-0.5,0.5)$ and then $u$ was normalized to have unit length. The order of $A$ was $n=1000$. We solved a sequence of 10 problems generated with different seeds, for a fixed tolerance $\varepsilon_{\Delta}=10^{-6}$ and $\Delta$ varying from 100 to 0.0001 by a factor of 10 , for a total of 70 problems. The initial $\delta_{U}$ was set to -4.5 and $\alpha_{0}=\min \left\{0, \alpha_{U}\right\}$. The tolerance for a quasi-optimal solution was set to $\varepsilon_{H C}=10^{-6}$.

The parameters for the IRLM were the following. For $\Delta=100,10$ the number of Lanczos basis vectors was 30 , and 20 shifts were applied on each implicit restart, while

Table 3
Average behavior for different trust-region radii.

| $\Delta$ | 100 | 10 | 1 | 0.1 | 0.01 | 0.001 | 0.0001 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| IT | 9.9 | 7.7 | 4.6 | 4 | 3.8 | 3 | 3 |
| LSTR MV | 1032.4 | 354 | 46.9 | 38.8 | 37.2 | 29.4 | 29.4 |
| $\frac{\left\\|g+\left(A-\lambda_{*} I\right) x_{*}\right\\|}{\\|g\\|}$ | $10^{-3}$ | $10^{-3}$ | $10^{-2}$ | $10^{-11}$ | $10^{-2}$ | $10^{-2}$ | $10^{-13}$ |
| $\left\|\frac{\Delta-\left\\|x_{*}\right\\|}{\Delta}\right\|$ | $10^{-16}$ | $10^{-16}$ | $10^{-8}$ | $10^{-9}$ | $10^{-12}$ | $10^{-10}$ | $10^{-7}$ |

for $\Delta \leq 1$, the number of vectors was nine with six shifts on each implicit restart. The maximum number of restarts was 150 and 45 , respectively. The difference in the number of basis vectors is due to the fact that for larger radii the hard case and near hard case are more likely to occur, and therefore the smallest eigenvalues of the bordered matrix become more clustered and the IRLM needs more space and iterations to compute the desired eigenpairs to the required accuracy. The initial vector for the IRLM was chosen as in section 6.1. We relaxed the accuracy required in the eigenvalue solution in the following way. The initial values for $\varepsilon_{\text {Lan }}$ were $0.03,0.1$, and 0.25 for $\Delta=100,10$, and $\Delta<10$, respectively. The value of $\varepsilon_{\text {Lan }}$ was kept the same until $\left|\frac{\Delta-\left\|x_{k}\right\|}{\Delta}\right|<0.1$, when $\varepsilon_{\text {Lan }}=0.015,0.05$, and 0.125 for $\Delta=100,10$ and $\Delta<10$, respectively. The results of the experiment are shown in Table 3, where we also report the average number of matrix-vector products required by the conjugate gradient method to solve the systems $\left(A-\lambda_{k} I\right) x=-g$ for $\lambda_{k}$ generated by LSTRS.

As observed in [17], the conjugate gradient method has a much easier time for smaller values of $\Delta$.
6.3. Superlinear convergence. The purpose of the third experiment was to verify superlinear convergence. The matrix $A$ was again set to $A=L-5 I$ with $L$ the 2D discrete Laplacian on the unit square, but now $n=256$. The vector $g$ was randomly generated with entries uniformly distributed on $(-0.5,0.5)$. We studied problems with and without hard case. To generate the hard case, we orthogonalized the vector $g$ randomly generated as before against the eigenvector $q$ corresponding to the smallest eigenvalue of $A$. We accomplished this by setting $g \leftarrow g-q\left(q^{T} g\right)$. For the problem without hard case the trust-region radius was $\Delta=10$ and $\varepsilon_{\Delta}=10^{-11}$. For the problem with hard case the radius was $\Delta=100$ and $\varepsilon_{H C}=10^{-11}$. The eigenproblems were solved with the MATLAB routine eig. The results are shown in Table 4, where we report the quantity $\left|\frac{\Delta-\left\|x_{k}\right\|}{\Delta}\right|$ for the problem without hard case and the quantity $\frac{\left(\lambda_{i}(\alpha)-\lambda_{1}(\alpha)\right) \tau_{2}^{2}\left(1+\Delta^{2}\right)}{-2 \eta \psi(\widetilde{x})}$ from Theorem 3.2 for the problem with hard case.

The quantity $\left\|\left(A-\lambda_{*} I\right) x_{*}+g\right\| /\|g\|$ was of order $10^{-14}$ for problem (a) and $10^{-7}$ for problem (b). An asterisk $*$ in the hard case means that we could not check for a quasi-optimal solution since the conditions of Lemma 3.5 were not satisfied.
6.4. The hard case. The fourth experiment illustrates the behavior of the method in the hard case. The matrix $A$ was of the form $A=U D U^{T}$, with $D=$ $\operatorname{diag}\left(d_{1}, \ldots, d_{n}\right)$ and $U=I-2 u u^{T}, u^{T} u=1$. The elements of $D$ were randomly generated with a uniform distribution on $(-5,5)$ then sorted in nondecreasing order and $d_{i}$ set to -5 for $i=1,2, \ldots, \ell$, allowing multiplicity $\ell$ for the smallest eigenvalue of $A$. Both vectors $u$ and $g$ were randomly generated with entries selected from a

TABLE 4
Verification of superlinear convergence for problems without hard case (a) and with hard case (b).

| $k$ | $\left\|\frac{\Delta-\left\\|x_{k}\right\\|}{\Delta}\right\|$ |
| :---: | :---: |
| 0 | $8.739485 \mathrm{e}-01$ |
| 1 | $1.101152 \mathrm{e}+01$ |
| 2 | $7.790406 \mathrm{e}-01$ |
| 3 | $5.987336 \mathrm{e}-01$ |
| 4 | $1.247129 \mathrm{e}-01$ |
| 5 | $2.593978 \mathrm{e}-02$ |
| 6 | $3.410990 \mathrm{e}-04$ |
| 7 | $1.038581 \mathrm{e}-06$ |
| 8 | $4.049703 \mathrm{e}-11$ |
| 9 | $8.704149 \mathrm{e}-15$ |

(a)

| $k$ | $\frac{\left(\lambda_{i}\left(\alpha_{k}\right)-\lambda_{1}\left(\alpha_{k}\right)\right) \tau_{2}^{2}\left(1+\Delta^{2}\right)}{-2 \eta \psi(x)}$ |
| :---: | :---: |
| 0 | $1.694375 \mathrm{e}-01$ |
| 1 | $*$ |
| 2 | $4.112269 \mathrm{e}-02$ |
| 3 | $9.276102 \mathrm{e}-03$ |
| 4 | $6.306448 \mathrm{e}-04$ |
| 5 | $5.851597 \mathrm{e}-06$ |
| 6 | $4.159997 \mathrm{e}-09$ |
| 7 | $2.116485 \mathrm{e}-09$ |
| 8 | $7.976599 \mathrm{e}-10$ |
| 9 | $1.267130 \mathrm{e}-12$ |

(b)
uniform distribution on $(-0.5,0.5)$ and then $u$ was normalized to have unit length. The order of $A$ was $n=1000$.

In this case, the eigenvectors of the matrix $A$ are of the form $q_{i}=e_{i}-2 u u_{i}, i=$ $1, \ldots, n$, with $e_{i}$ the $i$ th canonical vector in $\mathbb{R}^{n}$ and $u_{i}$ the $i$ th component of the vector $u$. This provides complete control in the generation of the hard case. In fact, if $\ell=1$, the vector $g$ was orthogonalized against $q_{1}$ computed by the formula given above. For $\ell>1, g$ was computed as the sum of the vectors in an orthonormal basis for the orthogonal complement of $\mathcal{S}_{1}$. After this, a noise vector $s$ was added to $g$ and $g \leftarrow \frac{g+s}{\|g+s\|}$. Both hard case and near hard case were generated by adding noise vectors of norms $10^{-8}$ and $10^{-2}$, respectively. To ensure that the hard case really occurred, we computed $\Delta_{\text {min }}=\left\|\left(A-d_{1} I\right)^{\dagger} g\right\|$ and set $\Delta=2 \Delta_{\text {min }}$. The problems were solved to the level $\varepsilon_{H C}=10^{-6}$. The initial $\delta_{U}$ was set to -4.5 and $\alpha_{0}=\min \left\{0, \alpha_{U}\right\}$.

The parameters for the IRLM were chosen as follows: for the hard case, 9 Lanczos basis vectors with 6 shifts on each implicit restart and a maximum of 45 restarts; for the near hard case, 18 Lanczos basis vectors with 12 shifts on each implicit restart and a maximum of 90 restarts. The different number of basis vectors is due to the fact that in the near hard case the smallest eigenvalues of the bordered matrix become more clustered and the IRLM needs more space in order to compute the desired eigenpairs. The tolerance $\varepsilon_{\text {Lan }}$ was fixed at $10^{-2}$.

In Table $5(\mathrm{a})$, (b) we summarize the average results for a sequence of 10 problems, generated with different seeds, for problems with hard case and near hard case, respectively.
6.5. Comparison with the semidefinite programming (SDP) approach. Finally, we compared LSTRS with the SDP approach of [13]. In this experiment, we solved two different families of problems. For each family, we generated 10 problems of each type (easy and hard case) with different seeds and solved them with Algorithm 4.1 (LSTRS) and the SDP approach of [13]. In all cases, the eigenproblems were solved with the function eig of MATLAB, so that the eigenpairs available to both methods had the same level of accuracy and also to avoid the inconsistencies associated with having two different eigensolvers. We report the average number of iterations (IT), average magnitude of the residual $\left\|\left(A-\lambda_{*} I\right) x_{*}+g\right\| /\|g\|$, and average relative accuracy in the norm of the trust-region solution, $\mid \Delta-\left\|x_{*}\right\| \| / \Delta$. Since we were using the function eig as the eigensolver, we are also reporting the average number of calls

Table 5
(a) The hard case and (b) near hard case, when $\mathcal{S}_{1}$ has dimension $\ell \geq 1$.

| $\ell$ | $M V$ | $I T$ | $\frac{\left\\|\left(A-\lambda_{*} I\right) x_{*}+g\right\\|}{\\|g\\|}$ |
| :---: | :---: | :---: | :---: |
| 1 |  | 9.6 | $10^{-2}$ |
| 5 | 790.4 | 12.1 | $10^{-2}$ |
| 10 | 1301.9 | 22.6 | $10^{-2}$ |

(a)

| $\ell$ | $M V$ | $I T$ | $\frac{\left\\|\left(A-\lambda_{*} I\right) x_{*}+g\right\\|}{\\|g\\|}$ |
| :---: | :---: | :---: | :---: |
| 1 | 1153.5 | 10 | $10^{-3}$ |
| 5 | 1039.9 | 9.7 | $10^{-3}$ |
| 10 | 1063.5 | 9.7 | $10^{-3}$ |
| $(\mathrm{~b})$ |  |  |  |

TABLE 6
Comparison with SDP approach. First set of problems, $\varepsilon_{H C}=10^{-8}$.

|  |  | IT | SOLVES | $\frac{\left\\|g+\left(A-\lambda_{*} I\right) x_{*}\right\\|}{\\|g\\|}$ | $\left\|\frac{\Delta-\left\\|x_{*}\right\\|}{\Delta}\right\|$ |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $A=L-5 I$ | Easy <br> case | LSTRS |  | 5.0 | $10^{-13}$ | $10^{-7}$ |
|  |  | SDP | 4.8 | 5.8 | $10^{-3}$ | $10^{-3}$ |
|  |  | SDP | 8.0 | 9.7 | $10^{-9}$ | $10^{-16}$ |
|  |  |  | 10.1 | $10^{-7}$ | $10^{-7}$ |  |

to the eigensolver (SOLVES) to provide a means of comparing the amount of work needed by each method. It is important to point out that in large-scale applications the computational effort will concentrate on solving the eigenvalue problems, and therefore in such situations we should also compare the cost of solving each eigenvalue problem.

In the first family of problems, the matrix $A$ was $A=L-5 I$ of order $n=256$ and the vector $g$ was randomly generated with entries uniformly distributed on $(0,1)$. As in section 6.3 , we orthogonalized $g$ against the eigenvector of $A$ corresponding to $\delta_{1}$ to generate the hard case. For both easy and hard cases we added a noise vector to $g$, of norm $10^{-8}$. The trust-region radius was $\Delta=100$. We used $\varepsilon_{\Delta}=10^{-6}$ and we ran the experiments with $\varepsilon_{H C}=10^{-8}$ and $\varepsilon_{H C}=10^{-6}$. We report these results in Tables 6 and 7 , respectively.

In the second family of problems, $A, g$, and $\Delta_{\min }$ were generated exactly as in section 6.4, where $A=U D U^{T}$ of order $n=256$. For the easy case, $\Delta=0.1 \Delta_{\min }$ and for the hard case $\Delta=5 \Delta_{\min }$. The tolerances used for Algorithm 4.1 were $\varepsilon_{\Delta}=10^{-6}$ and $\varepsilon_{H C}=10^{-6}$. The results are reported in Table 8.

The previous tests indicate a marginal advantage to our algorithm in most cases. We believe this is partially due to the fact that in the SDP approach it is necessary to compute the smallest eigenvalue of $A$ in order to begin the major iteration, while our approach avoids this extra calculation. From the comparative results, we can see that LSTRS obtained solutions with improved feasibility over the ones computed by the SDP approach. Moreover, LSTRS required slightly less computational effort overall to compute the solutions, especially in the hard case.
7. Conclusions. We have presented a new algorithm for the large-scale trustregion subproblem. The algorithm is based upon embedding the trust-region problem into a family of parameterized eigenvalue problems as developed in [17]. The main

TABLE 7
Comparison with SDP approach. First set of problems, $\varepsilon_{H C}=10^{-6}$.

|  |  | IT | SOLVES | $\frac{\left\\|g+\left(A-\lambda_{*} I\right) x_{*}\right\\|}{\\|g\\|}$ | $\left\|\frac{\Delta-\left\\|x_{*}\right\\|}{\Delta}\right\|$ |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $A=L-5 I$ | Easy <br> case |  |  | 4.5 | $10^{-2}$ | $10^{-8}$ |
|  |  | SDP | 4.8 | 5.8 | $10^{-3}$ | $10^{-3}$ |
|  |  | SDP | 7.0 | 8.7 | $10^{-7}$ | $10^{-16}$ |
|  |  | 10.1 | $10^{-7}$ | $10^{-7}$ |  |  |

Table 8
Comparison with SDP approach. Second set of problems.

|  |  | IT | SOLVES | $\frac{\left\\|g+\left(A-\lambda_{*} I\right) x_{*}\right\\|}{\\|g\\|}$ | $\left\|\frac{\Delta-\left\\|x_{*}\right\\|}{\Delta}\right\|$ |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $A=U D U^{T}$ | Easy <br> case |  |  | 8.7 | $10^{-3}$ | $10^{-14}$ |
|  |  | 4.4 | 5.4 | $10^{-4}$ | $10^{-4}$ |  |
|  | Hard | LSTRS | 6.4 | 12.5 | $10^{-3}$ | $10^{-8}$ |
|  | case | SDP | 13.8 | 14.8 | $10^{-5}$ | $10^{-5}$ |

contribution of this paper has been to give a better understanding of the hard-case condition and to utilize this understanding to develop a better treatment of this case. As a result, we have designed a unified algorithm that naturally incorporates both the standard and hard cases.

We have proved that the iterates for this new algorithm converge either to an optimal pair for the trust-region subproblem or to a pair that can be used to construct a quasi-optimal solution. We have proved that the rate of convergence is superlinear and we have demonstrated this computationally for both the standard and hard cases. This result represents a major improvement over the performance of the method originally presented in [17]. That approach used a different iteration for the hard case that was linearly convergent. In practice this behavior seemed to occur often and greatly detracted from the performance. We have also compared our method to the SDP approach presented in [13], obtaining better results in terms of feasibility.

Our motivation for developing the LSTRS method came from some important large-scale applications. In particular, the regularization of ill-posed problems such as those arising in seismic inversion [21] provides an important class of trust-region subproblems. It was shown in [14] that near hard cases are common for this class of problems, where the vector $g$ is nearly orthogonal to eigenspaces corresponding to several of the smallest eigenvalues of $A$. The work in [14] also reports the successful application of LSTRS to the regularization of discrete forms of ill-posed problems from inverse problems, including problems with field data.

Further work should include an analysis of the quasi-optimal solutions computed by LSTRS, the use of LSTRS within a trust-region method for the solution of largescale optimization problems, and an analysis of such a method in light of the work in [4].

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