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## On the numerical solution of large-scale sparse discrete-time Riccati equations

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Institutions: Max Planck Society, Braunschweig University of Technology
Published on: 01 Nov 2011 - Advances in Computational Mathematics (Springer)
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## Chemnitz Scientific Computing Preprints

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Impressum:
Chemnitz Scientific Computing Preprints - ISSN 1864-0087
(1995-2005: Preprintreihe des Chemnitzer SFB393)
```

Herausgeber:
Professuren für
Numerische und Angewandte Mathematik an der Fakultät für Mathematik der Technischen Universität Chemnitz

Postanschrift:
TU Chemnitz, Fakultät für Mathematik 09107 Chemnitz
Sitz:
Reichenhainer Str. 41, 09126 Chemnitz

# Chemnitz Scientific Computing Preprints 

Peter Benner Heike Faßbender

# On the numerical solution of large-scale sparse discrete-time Riccati equations 

CSC/09-11

## Contents

1 Introduction ..... 1
2 Newton's method for discrete-time algebraic Riccati equations ..... 5
3 Smith's Method and ADI Iteration ..... 8
3.1 Smith's method ..... 9
3.2 ADI Iteration ..... 18
4 Numerical Examples ..... 23
5 Conclusions ..... 28
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#### Abstract

The numerical solution of Stein (aka discrete Lyapunov) equations is the primary step in Newton's method for the solution of discrete-time algebraic Riccati equations (DARE). Here we present a low-rank Smith method as well as a low-rank alternating-direction-implicit-iteration to compute lowrank approximations to solutions of Stein equations arising in this context. Numerical results are given to verify the efficiency and accuracy of the proposed algorithms.


Keywords. Riccati equation, discrete-time control, sparse matrices, Newton's method, Smith iteration, ADI iteration, low rank factor.

AMS subject classifications (MSC2010). 65F30, 15A24, 93C20, 93C05.

## 1 Introduction

After full (time and spatial) discretization, the solution of optimal control problems for parabolic partial differential equations (PDEs) leads to generalized largescale, sparse discrete-time algebraic Riccati equation (DARE)

$$
\begin{align*}
0=\mathcal{R}(X)= & C Q C^{T}+A X A^{T}-E X E^{T}  \tag{1}\\
& -\left(A X B^{T}+C S^{T}\right)\left(R+B X B^{T}\right)^{-1}\left(B X A^{T}+S C^{T}\right)
\end{align*}
$$

where $A$ and $E \in \mathbb{R}^{n \times n}$ are large and sparse, $B \in \mathbb{R}^{m \times n}, C \in \mathbb{R}^{n \times p}, Q \in \mathbb{R}^{p \times p}, R \in$ $\mathbb{R}^{m \times m}, S \in \mathbb{R}^{m \times p}$, and $Q$ and $R$ are symmetric. We also assume $E$ to be invertible throughout.

In the last decades, much research has addressed the construction of numerically robust algorithms for the solution of (1). However, these methods generally have at least a memory complexity $\mathcal{O}\left(n^{2}\right)$ and a computational complexity $\mathcal{O}\left(n^{3}\right)$, regardless whether or not the system matrix A is sparse or otherwise structured. Therefore, the majority of numerical algorithms is restricted to systems of moderate order. Of course, the upper limit for this order depends on the problem to be solved as well as the particular computing environment and may vary between a few hundreds and a few thousands. However, a significant number of applications lead to systems of larger order. Large systems arise from the semi-discretization of (possibly linearized) PDEs by means of finite differences or finite elements, see e.g. $[4,36,40]$, and many more.

Consider, e.g., the problem

$$
\text { Minimize } \mathcal{J}(u)=\frac{1}{2} \int_{0}^{\infty} y(t)^{T} Q y(t)+u(t)^{T} R u(t)+2 y(t)^{T} S u(t) d t
$$

subject to the PDE constraint

$$
\begin{array}{rlrl}
\frac{\partial}{\partial t} x(\zeta, t) & =\nabla(k(\zeta) \nabla) x+c(\zeta) \nabla x+r(\zeta) x+b(\zeta) u \quad \text { in } \Omega \times[0, T] \\
x(\zeta, t) & =0 & \text { on } \partial \Omega, \\
x(\zeta, 0) & =x_{0}(\zeta) & \text { on } \Omega, \\
y(t) & =\int_{\Omega_{0}} c(\zeta) x(\zeta, t) d \zeta, & \Omega_{0} \subset \bar{\Omega} .
\end{array}
$$

Discretizing the PDE constraint using a finite-difference scheme or a finite element approach leads to a constraint in form of an ordinary differential equation

$$
\begin{aligned}
M \dot{x}(t) & =K x(t)+F u \\
y(t) & =C x(t)
\end{aligned}
$$

where $M$ and $K$ are large, sparse square matrices, $M$ is positive definite and $K$ is negative definite. Employing further a time discretization a difference equation of the form

$$
E x_{k+1}=A x_{k}+B u_{k}
$$

is obtained. In case the semi-implicit Euler method with stepsize $\Delta t$ is used, one obtains

$$
\begin{equation*}
(M-\Delta t K) x_{k+1}=M x_{k}+\Delta t F u_{k}, \tag{2}
\end{equation*}
$$

that is $A=M, E=M-\Delta t K$ and $B=\Delta t F$. Alternatively, a Crank-Nicholsontype discretization can be used. Here, we will use the average of a forward Euler and a semi-implicit Euler step in time. That is, we average the equations

$$
\begin{aligned}
& M x_{k+1}=(M+\Delta t K) x_{k}+\Delta t F u_{k}, \\
& M x_{k+1}=M x_{k}+\Delta t K x_{k+1}+\Delta t F u_{k},
\end{aligned}
$$

and obtain

$$
\begin{aligned}
\left(M-\frac{\Delta t}{2} K\right) x_{k+1} & =\left(M+\frac{\Delta t}{2} K\right) x_{k}+\Delta t F u_{k} \\
y_{k} & =C x_{k} .
\end{aligned}
$$

The discretized optimal control problem now reads

$$
\text { Minimize } \mathcal{J}(u)=\frac{1}{2} \sum_{k=0}^{\infty} y_{k}^{T} Q y_{k}+u_{k}^{T} R u_{k}+2 y_{k}^{T} S u_{k}
$$

subject to

$$
\begin{aligned}
E x_{k+1} & =A x_{k}+B u_{k}, \\
y_{k} & =C x_{k},
\end{aligned}
$$

where $E=M-\frac{\Delta t}{2} K, A=M+\frac{\Delta t}{2} K, B=\Delta t F$.
Under generic control-theoretic conditions, the optimal solution of this minimization problem is given by the feedback control

$$
u_{k}=-\left(R+B^{T} X_{d} B\right)^{-1}\left(B^{T} X_{d} A+S^{T} C\right) x_{k}, \quad k=0,1, \ldots,
$$

where $X_{d}$ is the stabilizing solution of the DARE (1); see, e.g., [21, 39, 44, 52] and many other textbooks on control theory. DAREs also arise in other applications such as $\mathcal{H}_{\infty}$-control, factorization problems for rational matrix functions, Kalman filtering. An overview of some of these applications is given in [39, Chapter 5]. A detailed discussion of the solution theory for the case $E=I$ (which is equivalent to the case $E$ nonsingular) is given in [39], whereas the case of singular $E$ is treated in [44]. Solutions of the optimal control problem without solving the corresponding Riccati equation are given in $[7,16,35,43,44]$, including the singular $E$ case. It should be noted, though, that DARE-free solutions of the linear-quadratic discrete-time optimal control problem require the full $n$-dimensional deflating subspace of the corresponding (generalized) symplectic pencil. For large-scale problems, this is prohibitive as even if the sparsity structure of the problemdefining matrices could be employed, the subspace itself would require memory of size $2 n^{2}$.

Mostly, systems originating from the application mentioned above possess two interesting properties. First, their order $n$ is large (say, $n>1000$ ), but the dimensions of the input and output spaces are relatively small ( $m, q \ll n$, often $m, q \ll 10$ ). For example, the order of a system arising from a parabolic PDE is about the number of grid points or mesh nodes used for the semi-discretization w.r.t. the spatial coordinates, which is relatively large. In contrast, $m$ and $q$ are often quite small and independent of the fineness of the discretization. Second, the system matrix $A$ is structured. Often, $A$ is a sparse matrix or it is implicitly represented as a product of sparse matrices and inverses of sparse matrices. In general, this structure allows the numerically inexpensive realization of matrixvector products and the efficient solution of systems of linear equations with $A$.

In the sequel we will use the abbreviation

$$
\begin{equation*}
K(X):=\left(A X B^{T}+C S^{T}\right)\left(R+B X B^{T}\right)^{-1}, \tag{3}
\end{equation*}
$$

to simplify the notation and we will make the following assumptions:

1. $E$ is nonsingular.
2. $R=R^{T}>0$.
3. $\left[\begin{array}{ll}Q & S^{T} \\ S & R\end{array}\right] \geq 0$.
4. A stabilizing solution $X_{d}$ of (1) exists, that is, a solution $X_{d}$ exists such that the eigenvalues of $\left(A-K\left(X_{d}\right) B\right)-\lambda E$ lie in the open unit circle: $\left(\sigma\left(A-K\left(X_{d}\right) B, E\right) \subset \mathcal{D}_{1}(0)\right)$. It is unique, and furthermore, $R+B X_{d} B^{T}>0$.

For sufficient conditions for 4 . to hold, see, e.g., [39, Theorem 13.1.3].

In principle, by inverting $E$, (1) can be reduced to the case $E=I$. However, this introduces unnecessary rounding errors and, if $E$ is ill-conditioned, even numerical instability. Therefore, inverting $E$ is avoided here. Considering $E$ is important, as in the PDE control problems $E \neq I$, and as $E$ is large and sparse, its inverse in general would be large and dense, so that the large, sparse generalized DARE would be transformed into a large, dense one for which there do not exist any suitable numerical methods yet.

The solution of DAREs has been an extremely active area of research, see, e.g., [21, 44, 52] for an overview. The usual solution methods for DAREs such as the Schur vector method [46], symplectic SR methods [12, 23], the matrix sign function $[7,15,26,51]$, the matrix disk function $[7,15,39,55]$ or the doubling method [52, 42] do not make (full) use of the sparse structure of $A, E$ and require in general $\mathcal{O}\left(n^{3}\right)$ flops and workspace of size $\mathcal{O}\left(n^{2}\right)$ even for sparse problems, and are therefore not suitable here.

The numerical solution of several types of large-scale matrix equations with sparse coefficient matrices arising in control theory has been the target of numerous papers in the last decade. Significant progress has been made in particular for continuous Lyapunov and algebraic Riccati equations, e.g., [8, 13, 53, 41, 49, 33]. For an overview and further references, see [9]. It is the aim of this paper to extend some of these results to DAREs. We will follow in particular the approach taken in [13]. That is, we will make use of the fact that $\mathcal{R}(X)=0$ defines a system of nonlinear equations and can hence be solved by an appropriate Newton method as proposed in $[32,3]$. Newton's method is reviewed in Section 2. The main computational cost in the algorithm stems from the numerical solution of the Stein equation. Section 3 proposes low rank Smith and ADI iterations for its solution based on [13, 17, 41, 49]. Both methods compute an approximate lowrank Cholesky factor of the desired solution. Numerical experiments are reported in Section 4.

## 2 Newton's method for discrete-time algebraic Riccati equations

The function $\mathcal{R}(X)$ in (1) is a rational matrix function, $\mathcal{R}(X)=0$ defines a system of nonlinear equations. Inspired by Kleinman's formulation of a Newton method for continuous-time algebraic Riccati equations [37], Hewer [32] proposed a Newton method for solving DAREs. The algorithm was extended to the generalized equation as given in (1) by Arnold and Laub [3]. A discussion of its convergence properties can be found in [39, 44].

Newton's method for the numerical solution of DAREs can be formulated as given in Algorithm 1.

```
Algorithm 1 Newton's Method for the DARE
    Input: The coefficient matrices \(A, B, C, E, Q, R, S\) of the DARE (1), and a
    starting guess \(X_{0}\), so that \(\sigma\left(A-K\left(X_{0}\right) B, E\right) \subset \mathcal{D}_{1}(0)\) and \(R+B X_{0} B^{T}>0\).
    Output: An approximate solution \(X_{k+1}\) of the DARE (1) and an estimate
    \(N_{k}\) for the error matrix \(X^{*}-X_{k+1}\), where \(X^{*}\) is the stabilizing solution of
    \(\mathcal{R}(X)=0\).
    for \(k=0,1,2, \ldots\) do
        1. \(K_{k} \leftarrow K\left(X_{k}\right)\).
        2. \(A_{k} \leftarrow A-K_{k} B\).
        3. \(\mathcal{R}_{k} \leftarrow \mathcal{R}\left(X_{k}\right)\).
        4. Solve for \(N_{k}\) in the Stein equation
            \(A_{k} N_{k} A_{k}^{T}-E N_{k} E^{T}=-\mathcal{R}_{k}\).
        5. \(X_{k+1} \leftarrow X_{k}+N_{k}\).
    end for
```

We have the following convergence result for Algorithm 1 [32, 39, 44].
Theorem 2.1 If the assumptions (4) hold, and $X_{0}$ is stabilizing, then for the iterates produced by Algorithm 1 we have:
a) All iterates $X_{k}$ are stabilizing, i.e., $\sigma\left(A-K\left(X_{k}\right) B, E\right) \subset \mathcal{D}_{1}(0)$ for all $k \in \mathbb{N}_{0}$.
b) $X_{d} \leq \ldots \leq X_{k+1} \leq X_{k} \leq \ldots \leq X_{1}$.
c) $\lim _{k \rightarrow \infty} X_{k}=X_{d}$.
d) There exists a constant $\gamma>0$ such that

$$
\left\|X_{k+1}-X_{d}\right\| \leq \gamma\left\|X_{k}-X_{d}\right\|^{2}, \quad k \geq 1
$$

i.e., the $X_{k}$ converge globally quadratic to $X_{d}$.

The formulation of Algorithm 1 is analogous to the standard formulation of Newton's method as given, e.g., in [22, Algorithm 5.1.1] for the solution of nonlinear equations. Because of its robustness in the presence of rounding errors, in dense matrix computations one prefers to calculate the Newton step explicitly as in Algorithm 1 using the Stein equation

$$
\begin{equation*}
A_{k} N_{k} A_{k}^{T}-E N_{k} E^{T}=-\mathcal{R}_{k} \tag{5}
\end{equation*}
$$

rather than to use the mathematically equivalent formulation of the Newton step [3, 32, 39, 44],

$$
\begin{align*}
A_{k} X_{k+1} A_{k}^{T}-E X_{k+1} E^{T} & =-C Q C^{T}+C S^{T} K_{k}^{T}+K_{k} S C^{T}-K_{k} R K_{k}^{T} \\
& =:-\mathcal{C}\left(X_{k}\right)=-\mathcal{C}_{k}, \tag{6}
\end{align*}
$$

which determines $X_{k+1}$ directly. The coefficient matrices of the two Stein equations are the same, but the right-hand-sides are different; see Algorithm 2.

```
Algorithm 2 Newton-Hewer Method for the DARE
    Input: The coefficient matrices \(A, B, C, E, Q, R, S\) of the DARE (1), and a
    starting guess \(X_{0}\), so that \(\sigma\left(A-K\left(X_{0}\right) B, E\right) \subset \mathcal{D}_{1}(0)\) and \(R+B X_{0} B^{T}>0\).
    Output: An approximate solution \(X_{k+1}\) of the DARE (1).
    for \(k=0,1,2, \ldots\) do
        1. \(K_{k} \leftarrow K\left(X_{k}\right)\).
        2. \(A_{k} \leftarrow A-K_{k} B\).
        3. \(\mathcal{C}_{k} \leftarrow \mathcal{C}\left(X_{k}\right)\).
        4. Solve for \(X_{k+1}\) in the Stein equation
\[
A_{k} X_{k+1} A_{k}^{T}-E X_{k+1} E^{T}=-\mathcal{C}_{k} .
\]
    end for
```

The problem for the successful application of the Newton method is to find a stabilizing initial guess $X_{0}$. There exist stabilization procedures for discrete-time linear systems (see, e.g., $[2,38,52]$ ). But these may give large initial errors $\left\|X_{d}-X_{0}\right\|$. The procedure suggested in [38] is even unfeasable for numerical computations as it is based on explicitly summing up $A^{k} B B^{T}\left(A^{T}\right)^{k}$ for $k$ up to $n$, thereby often causing overflow already for small values of $n$. Suitable initialization procedures are suggested in $[10,11]$ which work quite well in practice, but
cannot exploit sparsity of the problem and are therefore not suitable for largescale problems. An iterative procedure for stabilizing large-scale discrete-time systems is suggested in [25] and can be used for our purposes.

Despite the ultimate rapid convergence indicated by Theorem 2.1 d ), the iteration may initially converge slowly. This can be due to a large initial error $\left\|X_{d}-X_{0}\right\|$ or a disastrously large first Newton step resulting in a large error $\left\|X_{d}-X_{1}\right\|$. In both cases, it is possible that many iterations are required to find the region of rapid convergence. An ill-conditioned Stein equation makes it difficult to compute an accurate Newton step. An inaccurately computed Newton step can cause the usual convergence theory to break down in practice. Sometimes rounding errors or a poor choice of $X_{0}$ cause Newton's method to converge to a nonstabilizing solution. Fortunately, in PDE control as considered here, we often have $\sigma(E, A) \subset \mathcal{D}_{1}(0)$, so that $X_{0}=0$ is an appropriate starting guess.

The computational cost for Algorithms 1 and 2 mainly depends upon the cost for the numerical solution of the Stein equation (5), (6), resp.. This can be done using the Bartels-Stewart algorithm [5, 6] or an extension to the case $E \neq I$ [27, 28, 48]. The Bartels-Stewart algorithm is the standard direct method for the solution of Stein equations of small to moderate size. This method requires the computation of a Schur decomposition, and thus is not appropriate for large scale problems. The cost for the solution of the Stein equation is $\approx 73 n^{3}$ flops. See [52] for a discussion of an efficient implementation. In [52], the following idea based on $[31,56]$ how to solve stable nonnegative Stein equations

$$
\begin{equation*}
A X A^{H}-X=-C \tag{7}
\end{equation*}
$$

where $C \in \mathbb{C}^{n \times n}$ is positive semidefinite and $A \in \mathbb{C}^{n \times n}$ is stable, is also discussed: The unique solution $X$ is positive semidefinite and hence allows for a Cholesky factorization $X=L L^{H}$. Let $U$ be the unitary matrix that reduces $A$ to upper Schur form $T=U A U^{H}$, such that we have

$$
\left(U^{H} T U\right)\left(L L^{H}\right)\left(U^{H} T^{H} U\right)-L L^{H}=-B B^{H}
$$

where $B B^{H}=C$ and $B \in \mathbb{C}^{n \times m}, m \leq n$. Or, equivalently,

$$
T(U L)\left(L^{H} U^{H}\right) T^{H}-(U L)\left(L^{H} U^{H}\right)=-(U B)\left(B^{H} U^{H}\right)
$$

Using the $R Q$ factorizations of $U B$ and $U L$ we obtain

$$
T Z Z^{H} T^{H}-Z Z^{H}=-R R^{H}
$$

where $Z$ and $R$ are triangular matrices. This equation can be solved efficiently in $Z$, the Cholesky factor of the solution of the original equation (7), see [52] for details. The memory space needed for the data is $2 n^{2}$, while the computational cost is about $\approx c n^{3}$ flops, where $c$ is of the order of unity. Although this method
is not suitable in the context of large sparse DAREs, we will make use of the idea of computing only a (Cholesky) factor of the desired solution.

Other iterative schemes for solving the Stein equation (5), (6), resp., have been developed. Most of these were first formulated for the continuous-time case. Examples are the Smith method [54], the sign-function method [51], and the alternating direction implicit (ADI) iteration method [57]. Unfortunately, all of these methods compute the solution in dense form and hence require $\mathcal{O}\left(n^{2}\right)$ storage. In case the solution to the Stein equation has low numerical rank (i.e., the eigenvalues decay rapidly) one can take advantage of this low rank structure to obtain approximate solutions in low rank factored form [14]. If the effective rank is $r \ll n$, then the storage is reduced from $\mathcal{O}\left(n^{2}\right)$ to $\mathcal{O}(n r)$. This approach will be discussed in detail in Section 3.

## 3 Smith's Method and ADI Iteration

In this section, iterative methods for solving the equivalent Stein equations

$$
\begin{equation*}
A_{k} Y A_{k}^{T}-E Y E^{T}=-\mathcal{R}_{k} \tag{8}
\end{equation*}
$$

and

$$
\begin{equation*}
A_{k} Y A_{k}^{T}-E Y E^{T}=-\mathcal{C}_{k} \tag{9}
\end{equation*}
$$

with

$$
A_{k}=A-K_{k} B, \quad K_{k}=K\left(X_{k}\right), \quad \mathcal{R}_{k}=\mathcal{R}\left(X_{k}\right), \quad \mathcal{C}_{k}=\mathcal{C}\left(X_{k}\right)
$$

will be discussed. In particular, it will be assumed that the solution $Y$ has numerically low rank. This is usually the case if the DARE stems from a discretized PDE control problem as illustrated in Figure 1, where the eigenvalues of the solution of the DARE corresponding to the discrete heat equation example from the SLICOT Model Reduction Benchmark Collection [19] are shown, see also Example 4.2. Here, the numerical rank of $X$ is 31 as compared to $n=200$.

The symmetric Stein equation

$$
\begin{equation*}
S-A S A^{T}=V, \tag{10}
\end{equation*}
$$

where $A$ and $V$ are given $n \times n$ matrices and $V$ is symmetric, has a unique solution $S$ (necessarily symmetric) provided $z \bar{w} \neq 1$ for every pair of eigenvalues $z, w$ of $A$, see [39, pp. 104-105]. In particular, this is the case when all eigenvalues of $A$ lie in the open unit disk. Then the unique solution is given by

$$
\begin{equation*}
S=\sum_{j=0}^{\infty} A^{j} V A^{T j} \tag{11}
\end{equation*}
$$



Figure 1: Decay of the eigenvalues of the DARE solution corresponding to the discrete heat equation example from [19] with $n=200, m=p=1$.
(the series converges because $\left|z_{0}\right|<1$ for every $z_{0} \in \sigma(A)$ ). Moreover, if all eigenvalues of $A$ are in the open unit disk, then we have for the unique solution $S$ of (10) $S \geq 0$ if $V \geq 0$.

From this we have, that for given matrices $A_{k}, E \in \mathbb{R}^{n \times n}$ the Stein equations (8) and (9) have a unique solution if and only if $\lambda_{r} \lambda_{s} \neq 1$ for any $\lambda_{r}, \lambda_{s} \in \sigma\left(A_{k}, E\right)$. Hence, as the Stein equation (8) (or (9), respectively) arises in the Newton iteration, Theorem 2.1a) guarantees here the existence of a unique symmetric solution of the Stein equation.

### 3.1 Smith's method

The Smith iteration [54] is derived from the symmetric Stein equation (10) and its representation (11)

$$
S_{0}=0, \quad S_{j+1}=V+A S_{j} A^{T}
$$

If all eigenvalues of $A$ lie in the open unit disk, this iteration converges and the iterates can be written as

$$
S_{j}=\sum_{i=0}^{j-1} A^{i} V A^{T i}, \quad j=1,2, \ldots
$$

Its rate of convergence is quite slow, in general. The accelerated version [54] the so-called squared Smith method - is based on the recursion

$$
\begin{equation*}
S_{0}=V, \quad S_{j+1}=S_{j}+A^{2^{j}} S_{j} A^{T 2^{j}} \tag{12}
\end{equation*}
$$

Its iterates can be written as

$$
S_{j}=\sum_{i=1}^{2^{j}} A^{i-1} V A^{T(i-1)}
$$

Despite the quadratic convergence rate, one should be reluctant to apply the squared method to large, sparse equations. The matrices $A^{2^{j}}$, which have to be squared explicitly in each step of the iteration, are dense even if $A$ is sparse.

We will discuss how to use the standard Smith iteration efficiently for solving the Stein equations (8) and (9) in the context of large, sparse equations. The idea behind this is based on a low-rank version of (12) as proposed, e.g., in [14].

The Smith iteration for (8) is given by

$$
\begin{equation*}
E Y_{j+1} E^{T}=A_{k} Y_{j} A_{k}^{T}+\mathcal{R}_{k}, \tag{13}
\end{equation*}
$$

while for (9) it is given by

$$
\begin{equation*}
E Y_{j+1} E^{T}=A_{k} Y_{j} A_{k}^{T}+\mathcal{C}_{k} . \tag{14}
\end{equation*}
$$

The last term on the right hand side of (14) can be expressed as

$$
\begin{aligned}
\mathcal{C}_{k} & =C Q C^{T}-C S^{T} K_{k}^{T}-K_{k} S C^{T}+K_{k} R K_{k}^{T} \\
& =\left[\begin{array}{ll}
C & -K_{k}
\end{array}\right]\left[\begin{array}{cc}
Q & S^{T} \\
S & R
\end{array}\right]\left[\begin{array}{c}
C^{T} \\
-K_{k}^{T}
\end{array}\right] .
\end{aligned}
$$

As $\left[\begin{array}{ll}Q & S^{T} \\ S & R\end{array}\right] \geq 0$, we can obtain a (Cholesky) factorization

$$
\begin{align*}
& {\left[\begin{array}{ll}
C & -K_{k}
\end{array}\right]\left[\begin{array}{cc}
Q & S^{T} \\
S & R
\end{array}\right]\left[\begin{array}{c}
C^{T} \\
-K_{k}^{T}
\end{array}\right]} \\
& =\left[\begin{array}{ll}
C & -K_{k}
\end{array}\right]\left[\begin{array}{cc}
L & 0 \\
S L^{-T} & W
\end{array}\right]\left[\begin{array}{cc}
L^{T} & L^{-1} S^{T} \\
0 & W^{T}
\end{array}\right]\left[\begin{array}{c}
C^{T} \\
-K_{k}^{T}
\end{array}\right]  \tag{15}\\
& =\left[\begin{array}{cc}
C L & 0 \\
-K_{k} S L^{-T} & -K_{k} W
\end{array}\right]\left[\begin{array}{cc}
L^{T} C^{T} & -L^{-1} S^{T} K_{k}^{T} \\
0 & -W^{T} K_{k}^{T}
\end{array}\right]=U U^{T}
\end{align*}
$$

with $Q=L L^{T}$, and $W W^{T}=R-S L^{-1} L^{-T} S^{T}$. The latter expression is positive semidefinite due to Assumption 3. on page 4. Therefore, $U U^{T} \geq 0$. When the
iteration (14) is started with $Y_{0}=0$, then $Y_{1}=E^{-1} U U^{T} E^{-T}$ and all subsequent iterates

$$
Y_{j}=\sum_{i=0}^{j-1}\left(E^{-1} A_{k}\right)^{i} E^{-1} U U^{T} E^{-T}\left(A_{k}^{T} E^{-T}\right)^{i}, \quad j=1,2, \ldots
$$

are symmetric positive semidefinite and converge to a positive semidefinite matrix $Y_{*}$. Based on the assumption that the spectrum of the positive semidefinite matrix $Y_{*}$ decays to zero very rapidly, we can expect that $Y_{*}$ can be written using a factorization $Z Z^{T}$ for some $Z \in \mathbb{R}^{n \times r}, r \ll n$. Now, if we assume that $Y_{j}=Z_{j} Z_{j}^{T}$ for $\operatorname{rank}\left(Z_{j}\right)=r_{j} \ll n, Z_{j} \in \mathbb{R}^{n \times r_{j}}$ and observe that $\operatorname{rank}\left(U U^{T}+A_{k} Y_{j} A_{k}^{T}\right) \leq$ $m+p+r_{j} \ll n$, we see that we can solve (14) for a low-rank representation of the (Cholesky) factor of $Y_{j}$. In particular, we have

$$
\begin{align*}
\left(E Z_{j+1}\right)\left(Z_{j+1}^{T} E^{T}\right) & =\left(A_{k} Z_{j}\right)\left(Z_{j}^{T} A_{k}^{T}\right)+U U^{T}\left[\begin{array}{ll}
A_{k} Z_{j} & U
\end{array}\right]\left[\begin{array}{c}
\left(Z_{j}^{T} A_{k}^{T}\right) \\
U^{T}
\end{array}\right] \\
& =M M^{T} . \tag{16}
\end{align*}
$$

Hence,

$$
E Z_{j+1}=\left[\begin{array}{ll}
A_{k} Z_{j} & U
\end{array}\right]
$$

yields one possible representation of the solution $Y_{j+1}$. Thus, the Smith algorithm can be reformulated in terms of the (Cholesky) factor $Z_{j}$ of $Y_{j}$. There is no need to compute $Y_{j}$ at each iteration; only $Z_{j}$ is needed.
Incorporating this Smith iteration into Algorithm 2 we obtain Algorithm 3.
Remark 3.1 For $Z_{0}=0$ it is straightforward to see that

$$
Z_{j}=E^{-1}\left[\left(A_{k} E^{-1}\right)^{j-1} U \quad\left(A_{k} E^{-1}\right)^{j-2} U \quad \ldots \quad\left(A_{k} E^{-1}\right) U \quad U\right] .
$$

If $Z_{j} \in \mathbb{R}^{n \times r_{j}}$, then $M$ in (16) is in $\mathbb{R}^{n \times\left(r_{j}+p+m\right)}$ and, as $E \in \mathbb{R}^{n \times n}, Z_{j+1} \in$ $\mathbb{R}^{n \times\left(r_{j}+p+m\right)}$. The dimension of $Z_{j+1}$ will increase by $p+m$ in each iteration step. Hence, if $p+m$ is large and/or the convergence is slow, the number of columns of $Z_{j+1}$ will easily reach unmanageable levels of memory requirements. But if $Z_{j+1}$ is of low rank $r_{j+1}$, then we can approximate it as follows. The rank-revealing $L Q$ decomposition [20, 29] of $M$ yields

$$
\Pi M=\left[\begin{array}{ll}
L_{11} & 0 \\
L_{21} & \Omega
\end{array}\right] V, \quad L_{11} \in \mathbb{R}^{r_{j+1} \times r_{j+1}}, L_{21} \in \mathbb{R}^{\left(n-r_{j+1}\right) \times r_{j+1}},
$$

where $L_{11}$ is lower triangular, $V \in \mathbb{R}^{n \times\left(n+r_{j}\right)}$ is orthogonal, $\Pi$ is a permutation, and $\Omega \approx 0$ can be regarded as negligible. If we partition $V$ in the form

$$
V=\left[\begin{array}{c}
V_{1} \\
V_{2}
\end{array}\right], \quad V_{1} \in \mathbb{R}^{r_{j+1} \times n}
$$

```
Algorithm 3 Newton-Hewer-Smith Method for the DARE
    Input: The coefficient matrices \(A, B, C, E, Q, R, S\) of the DARE (1),
    and a starting guess \(X_{0}\) in terms of its (Cholesky) factor \(T_{0}\), so that \(\sigma(A-\)
    \(\left.K\left(X_{0}\right) B, E\right) \subset \mathcal{D}_{1}(0)\) and \(R+B X_{0} B^{T}>0\).
    Output: An approximate solution \(X_{k+1}\) of the DARE (1) in terms of its
    (Cholesky) factor \(T_{k+1}\).
    Compute the Cholesky factorization \(\left[\begin{array}{cc}Q & S^{T} \\ S & R\end{array}\right]=L L^{T}\).
    for \(k=0,1,2, \ldots\) do
        1. \(K_{k} \leftarrow K\left(X_{k}\right)\) (making use of the fact that \(\left.X_{k}=T_{k} T_{k}^{T}\right)\).
        2. \(A_{k} \leftarrow A-K_{k} B\).
        3. \(U \leftarrow\left[\begin{array}{ll}C & -K_{k}\end{array}\right] L\).
        4. \(Z_{0} \leftarrow 0\).
        for \(j=0,1,2, \ldots\) do
            5. Solve for \(Z_{j+1}\) in \(E Z_{j+1}=\left[\begin{array}{ll}A_{k} Z_{j} & U\end{array}\right]\).
        end for
        6. \(T_{k+1} \leftarrow Z_{j+1}\)
    end for
```

and set $\Omega=0$, we obtain the full-rank approximation

$$
\widetilde{M}=\Pi^{T}\left[\begin{array}{l}
L_{11} \\
L_{21}
\end{array}\right] V_{1}=\Pi \widetilde{L} V_{1}, \quad \widetilde{L} \in \mathbb{R}^{n \times r_{j+1}}
$$

In any unitary invariant norm, we have

$$
\|M-\widetilde{M}\|=\|\Omega\| .
$$

Now we approximate

$$
\left(E Z_{j+1}\right)\left(Z_{j+1}^{T} E^{T}\right) \approx \widetilde{M} \widetilde{M}^{T}=\Pi^{T} \widetilde{L} \widetilde{L}^{T} \Pi .
$$

Hence

$$
E Z_{j+1} \approx \Pi^{T} \widetilde{L}
$$

yields one possible approximate low-rank representation of the solution $Y_{j+1}$ of (14) such that $Z_{j+1} \in \mathbb{R}^{n \times r_{j+1}}$.

Alternatively, as suggested for the Smith iteration in [30], a thin singular value decomposition [29] of $M$ can be employed,

$$
M=U \Sigma V^{T}
$$

where, with $\ell=r_{j}+p+m, V \in \mathbb{R}^{\ell \times \ell}, V^{T} V=I, U \in \mathbb{R}^{n \times \ell}$ has orthonormal columns and $\Sigma=\operatorname{diag}\left(\sigma_{1}, \ldots, \sigma_{\ell}\right) \in \mathbb{R}^{\ell \times \ell}, \sigma_{1} \geq \sigma_{2} \geq \ldots \geq \sigma_{\ell}$. Partition this decomposition as

$$
M=\left[\begin{array}{ll}
U_{1} & U_{2}
\end{array}\right]\left[\begin{array}{ll}
\Sigma_{1} & \\
& \Sigma_{2}
\end{array}\right]\left[\begin{array}{ll}
V_{1}^{T} & V_{2}^{T}
\end{array}\right],
$$

where $\Sigma_{1} \in \mathbb{R}^{r_{j+1} \times r_{j+1}}, V_{1} \in \mathbb{R}^{\ell \times r_{j+1}}, U_{1} \in \mathbb{R}^{n \times r_{j+1}}$. In case $\Sigma_{2}$ is negligible, we can approximate $M$ by

$$
\widehat{M}=U_{1} \Sigma_{1} V_{1}^{T}
$$

By the Schmidt-Eckart-Young-Mirsky-theorem we have

$$
\|M-\widehat{M}\|=\sigma_{r_{j+1}+1} .
$$

Now we approximate

$$
\left(E Z_{j+1}\right)\left(Z_{j+1}^{T} E^{T}\right) \approx \widehat{M} \widehat{M}^{T}=U_{1} \Sigma_{1} V_{1}^{T} V_{1} \Sigma_{1} U_{1}^{T}=U_{1} \Sigma_{1}^{2} U_{1}^{T} .
$$

Hence

$$
E Z_{j+1} \approx U_{1} \Sigma_{1}
$$

yields another possible approximate low-rank representation of the solution $Y_{j+1}$ of (14) such that $Z_{j+1} \in \mathbb{R}^{n \times r_{j+1}}$.
Therefore, instead of computing the solution $Y$ of (9) directly, we only compute an approximation to its low rank factor $Z$ with $Y=Z Z^{T}$. In order to make use of this representation of the solution we have to modify Algorithm 3 by replacing Step 7. appropriately, see Algorithm 4. That is, we simply ignore the trailing part of $L$ which is less than or equal to some given tolerance. While going from the $j$ th to the $(j+1)$ st step, the number of columns of $Z_{j+1}$ generally does not increase. Indeed an increase will only occur if the rank of $Z_{j+1}$ is larger than that of $Z_{j}$. In any case, there can be at most $n$ additional columns added at any step which is the same as in the unmodified version.

Let us assume that $X_{k}=T_{k} T_{k}^{T}$ with $T_{k} \in \mathbb{R}^{n \times s_{k}}$. In order to obtain $Z_{j+1}$ one has to

- compute the (Cholesky) factorization of $\left[\begin{array}{ll}Q & S^{T} \\ S & R\end{array}\right] \in \mathbb{R}^{(p+m) \times(p+m)}$, this requires $\mathcal{O}\left((p+m)^{3} / 3\right)$ flops [29] but only needs to be computed once before the iteration is started;
- evaluate $K_{k}$, this involves five matrix-matrix products of matrices of size $m \times n$ and $n \times s_{k}, m \times s_{k}$ and $s_{k} \times m$ (for $R+\left(B T_{k}\right)\left(T_{k}^{T} B^{T}\right)$, re-using $B T_{k}$ and exploiting symmetry), $s_{k} \times n$ and $n \times n, m \times s_{k}$ and $s_{k} \times n$ (for $\left(B T_{k}\right)\left(T_{k}^{T} A\right)$, re-using $B T_{k}$ from the previous computation), and $m \times p$ and $p \times n$ (for $S C^{T}$, which should be precomputed), and a linear solve with a system matrix of size $m \times m$ and a right hand side of size $m \times n$ which requires $\mathcal{O}\left(2 m^{3} / 3\right)$ flops [29] plus forward and backward substitution;
- compute $U$, this involves one matrix-matrix product of matrices of size $n \times(p+m)$ and $(p+m) \times(p+m)$;
- compute $A_{k} Z_{j}=A Z_{j}-K_{k}\left(B Z_{j}\right)$ for $j=0,1, \ldots$, this involves three matrixmatrix products of matrices of size $n \times n$ and $n \times r_{j}, m \times n$ and $n \times r_{j}$, and $n \times m$ and $m \times r_{j}$ in each iteration step;

```
Algorithm 4 Low Rank Newton-Hewer-Smith Method for the DARE
    Input: The coefficient matrices \(A, B, C, E, Q, R, S\) of the DARE (1),
    and a starting guess \(X_{0}\) in terms of its (Cholesky) factor \(T_{0}\), so that \(\sigma(A-\)
    \(\left.K\left(X_{0}\right) B, E\right) \subset \mathcal{D}_{1}(0)\) and \(R+B X_{0} B^{T}>0\).
    Output: An approximate solution \(X_{k+1}\) of the DARE (1) in terms of its low
    rank (Cholesky) factor \(T_{k+1}\).
    Compute the Cholesky factorization \(\left[\begin{array}{ll}Q & S^{T} \\ S & R\end{array}\right]=L L^{T}\).
    for \(k=0,1,2, \ldots\) do
        1. \(K_{k} \leftarrow K\left(X_{k}\right)\) (making use of the fact that \(\left.X_{k}=T_{k} T_{k}^{T}\right)\).
        2. \(A_{k} \leftarrow A-K_{k} B\).
        3. \(\mathcal{C}_{k} \leftarrow \mathcal{C}\left(X_{k}\right)\).
        4. \(U \leftarrow\left[C-K_{k}\right] L\).
        5. \(Z_{0} \leftarrow 0\).
        for \(j=0,1,2, \ldots\) do
            6. Compute the \(R R L Q\) factorization \(\Pi\left[\begin{array}{ll}A_{k} Z_{j} & U\end{array}\right]=\left[\begin{array}{ll}L_{11} & 0 \\ L_{21} & \Omega\end{array}\right] V\),
                \(L_{11} \in \mathbb{R}^{r_{j+1} \times r_{j+1}}, V=\left[\begin{array}{c}V_{1} \\ V_{2}\end{array}\right], V_{1} \in \mathbb{R}^{r_{j+1} \times n}\)
            7. \(\widetilde{L} \leftarrow\left[\begin{array}{l}L_{11} \\ L_{12}\end{array}\right]\).
            8. Solve for \(Z_{j+1}\) in \(E Z_{j+1}=\Pi \widetilde{L}\).
        end for
        9. \(T_{k+1} \leftarrow Z_{j+1}\)
    end for
```

- compute the rank-revealing $L Q$ decomposition of $M=\left[\begin{array}{ll}A_{k} Z_{j} & U\end{array}\right]$, this requires $4 n\left(m+p+r_{j}\right) r_{j+1}-2\left(n+m+p+r_{j}\right) r_{j+1}^{2}+\frac{4}{3} r_{j+1}^{3}$ flops (see, e.g., [29, Section 5.4.1]) noting that $Q$ need not be accumulated;
- solve $E Z_{j+1}=\Pi \widetilde{L}$ for $Z_{j+1}$ for $j=0,1, \ldots$, this involves a (sparse) linear solve with a system matrix of size $n \times n$ and a right hand side of size $n \times r_{j+1}$. Please note, that the (sparse) $L U$ decomposition of $E$ needs to be computed only once, the rest of the computations are solely forward and backward solves. In the case of a DARE coming from a discretized PDE control problem, $E$ will often be positive definite so that a (sparse) Cholesky decomposition can be employed and the factorization cost can be halved.

Hence, the overall flop count for the computation of $Z_{j+1}$ depends mainly on the choice of the solver for the linear systems in the outer loop as well as on the cost of the rank-revealing $L Q$ decomposition in the inner loop and the pre-computable (Cholesky) factorizations.

The iteration (13) does not allow a factored representation of the solution $Y_{j+1}$ as $\mathcal{R}_{k}$ might be indefinite. But (in contrast to the continuous-time case) $\mathcal{R}_{k}$ can be split explicitly into its positive and negative semidefinite part $\mathcal{P}_{k}$ and $\mathcal{N}_{k}$, resp.,

$$
\mathcal{R}_{k}=\mathcal{P}_{k}-\mathcal{N}_{k}
$$

with

$$
\begin{aligned}
\mathcal{P}_{k} & =C Q C^{T}+A X_{k} A^{T} \\
\mathcal{N}_{k} & =E X_{k} E^{T}+\left(A X_{k} B^{T}+C S^{T}\right)\left(R+B X_{k} B^{T}\right)^{-1}\left(B X_{k} A^{T}+S C^{T}\right)
\end{aligned}
$$

Splitting the iterates into their positive and negative semidefinite parts $Y_{j}^{P}$ and $Y_{j}^{N}$, resp.,

$$
Y_{j}=Y_{j}^{P}-Y_{j}^{N}
$$

we have

$$
E\left(Y_{j+1}^{P}-Y_{j+1}^{N}\right) E^{T}=A_{k}\left(Y_{j}^{P}-Y_{j}^{N}\right) A_{k}^{T}+\left(\mathcal{P}_{k}-\mathcal{N}_{k}\right),
$$

that is, we can iterate on the positive and negative semidefinite parts separately

$$
\begin{align*}
E Y_{j+1}^{P} E^{T} & =A_{k} Y_{j}^{P} A_{k}^{T}+\mathcal{P}_{k}  \tag{17}\\
E Y_{j+1}^{N} E^{T} & =A_{k} Y_{j}^{N} A_{k}^{T}+\mathcal{N}_{k} \tag{18}
\end{align*}
$$

Now, with $Q=L L^{T}$ and $X_{k}=T_{k} T_{k}^{T}, T_{k} \in \mathbb{R}^{n \times s_{k}}$, we have

$$
\mathcal{P}_{k}=C Q C^{T}+A X_{k} A^{T}=C L L^{T} C^{T}+A T_{k} T_{k}^{T} A^{T}=\left[\begin{array}{ll}
C L & A T_{k}
\end{array}\right]\left[\begin{array}{c}
L^{T} C^{T} \\
T_{k}^{T} A^{T}
\end{array}\right] .
$$

Writing $Y_{j}^{P}=Z_{j}^{P}\left(Z_{j}^{P}\right)^{T}$ we obtain

$$
E Z_{j+1}^{P}\left(Z_{j+1}^{P}\right)^{T} E^{T}=\left[\begin{array}{lll}
A_{k} Z_{j}^{P} & C L & A T_{k}
\end{array}\right]\left[\begin{array}{c}
\left(Z_{j}^{P}\right)^{T} A_{k}^{T} \\
L^{T} C^{T} \\
T_{k}^{T} A^{T}
\end{array}\right]
$$

Hence,

$$
E Z_{j+1}^{P}=\left[\begin{array}{lll}
A_{k} Z_{j}^{P} & C L & A T_{k}
\end{array}\right]
$$

is one possible representation of the (Cholesky) factor $Z_{j+1}^{P}$ of the positive semidefinite part $Y_{j+1}^{P}$ of $Y_{j+1}$. If $Z_{j}^{P} \in \mathbb{R}^{n \times r_{j}^{P}}$, then $E Z_{j+1}^{P} \in \mathbb{R}^{n \times\left(r_{j}^{P}+p+s_{k}\right)}$ and, as $E \in \mathbb{R}^{n \times n}, Z_{j+1}^{P} \in \mathbb{R}^{n \times\left(r_{j}^{P}+p+s_{k}\right)}$. The dimension of $Z_{j+1}^{P}$ would increase in each iteration step. But if $Z_{j+1}^{P}$ is of low rank $r_{j+1}^{P}$, then we can approximate it using a rank-revealing $L Q$ factorization as before.

Moreover, with $X_{k}=T_{k} T_{k}^{T}$ and $R+B X_{k} B^{T}=S_{k} S_{k}^{T}$ we have

$$
\begin{aligned}
\mathcal{N}_{k} & =E X_{k} E^{T}+\left(A X_{k} B^{T}+C S^{T}\right)\left(R+B X_{k} B^{T}\right)^{-1}\left(B X_{k} A^{T}+S C^{T}\right) \\
& =E X_{k} E^{T}+\left(A X_{k} B^{T}+C S^{T}\right) S_{k}^{-T} S_{k}^{-1}\left(B X_{k} A^{T}+S C^{T}\right) \\
& =\left[E T_{k}\left(A X_{k} B^{T}+C S^{T}\right) S_{k}^{-T}\right]\left[\begin{array}{c}
T_{k}^{T} E^{T} \\
S_{k}^{-1}\left(B X_{k} A^{T}+S C^{T}\right)
\end{array}\right]
\end{aligned}
$$

From this, we obtain

$$
E Y_{j+1}^{N} E^{T}=\left[\begin{array}{lll}
A_{k} Z_{j}^{N} & E T_{k} & \left(A X_{k} B^{T}+C S^{T}\right) S_{k}^{-T}
\end{array}\right]\left[\begin{array}{c}
\left(Z_{j}^{N}\right)^{T} A_{k}^{T} \\
T_{k}^{T} E^{T} \\
S_{k}^{-1}\left(B X_{k} A^{T}+S C^{T}\right)
\end{array}\right]
$$

Hence, with $Y_{j}^{N}=Z_{j}^{N}\left(Z_{j}^{N}\right)^{T}$

$$
E Z_{j+1}^{N}=\left[\begin{array}{lll}
A_{k} Z_{j}^{N} & E T_{k} & \left(A X_{k} B^{T}+C S^{T}\right) S_{k}^{-T}
\end{array}\right]
$$

is one possible representation of a factor $Z_{j+1}^{N}$ of the negative semidefinite part $Y_{j+1}^{N}$ of $Y_{j+1}$. If $Z_{j}^{N} \in \mathbb{R}^{n \times r_{j}^{N}}$, then $E Z_{j+1}^{N} \in \mathbb{R}^{n \times\left(r_{j}^{N}+m+s_{k}\right)}$ and, as $E \in \mathbb{R}^{n \times n}$, $Z_{j+1}^{N} \in \mathbb{R}^{n \times\left(r_{j}^{N}+m+s_{k}\right)}$. The dimension of $Z_{j+1}^{N}$ would increase in each iteration step. But if $Z_{j+1}^{N}$ is of low rank $r_{j+1}^{N}$, then we can approximate it using a rankrevealing $L Q$ factorization as before.

From $Z_{j+1}^{P}$ and $Z_{j+1}^{N}$, one can determine $Y_{j+1}=Z_{j+1}^{P}\left(Z_{j+1}^{P}\right)^{T}-Z_{j+1}^{N}\left(Z_{j+1}^{N}\right)^{T}$. The overall flop count for this computation of $Y_{j+1}$ is higher than the one for the iteration (14). As for (14), $K_{k}$ and $A_{k}$ have to be set up. In case this is done in a reasonable way, $S_{k}^{-1}\left(B X_{k} A^{T}+S C^{T}\right)$ is computed as by-product when setting up $K_{k}$. Instead of the (Cholesky) factorization of $\left[\begin{array}{ll}Q & S^{T} \\ S & R\end{array}\right] \in \mathbb{R}^{(p+m) \times(p+m)}$, here only the factorization of $Q=L L^{T} \in \mathbb{R}^{p \times p}$ is needed. While before, $A_{k} Z_{j}$
was computed, we now need $A_{k} Z_{j}^{P}, A_{k} Z_{j}^{N}, A T_{k}$ and $E T_{k}$. Moreover, two instead of one rank-revealing $Q L$ factorizations have to be computed and two linear systems with the system matrix $E$ have to be solved (this can be done using one factorization of $E$ ).

As the solution $Y_{j+1}$ is constructed from two different iterations, one should use some kind of defect correction in order to increase the accuracy of the solution and to generate the (Cholesky) factor of $Y_{j+1}$.

The first idea that might come to mind is to use the following result from [45].
Theorem 3.2 Let $X$ be an Hermitian solution of

$$
0=\mathcal{R}(X),
$$

and let $\widetilde{X}$ be an Hermitian approximation to $X$. Let $V=X-\widetilde{X}$. If $\widetilde{R}=$ $R+B \widetilde{X} B^{T}$ and $I+V B^{T} \widetilde{R}^{-1} B$ are nonsingular, then $V$ satisfies the equation

$$
0=\mathcal{R}(\widetilde{X})-E V E^{T}+\widetilde{A} V \widetilde{A}^{T}-\left(\widetilde{A} V B^{T}\right)^{T}\left(\widetilde{R}+B V B^{T}\right)^{-1}\left(\widetilde{A} V B^{T}\right)
$$

where

$$
\widetilde{A}^{T}=-B^{T} \widetilde{R}^{-1} C S^{T}+\left(I-B^{T} \widetilde{R}^{-1} B \widetilde{X}\right) A^{T}
$$

Unfortunately, the resulting algorithm can not be used here, as in each iteration step, one has to solve (1). Using the approach discussed above, this would lead to a solution in terms of (Cholesky) factors of the positive and the negative semidefinite part of the solution and not to a (Cholesky) factor of the solution itself.

A different option might be to make use of the following observation which allows us to solve the DARE in each iteration step so that the desired (Cholesky) factor of the solution itself is computed. A "formidable computation" [39, p.312] shows that the DARE (1) can be rewritten as

$$
\begin{equation*}
E Y E^{T}-\breve{A} Y \breve{A}^{T}=\breve{Q}+\left(C S^{T}-K(Y) R\right) R^{-1}\left(C S^{T}-K(Y) R\right)^{T} \tag{19}
\end{equation*}
$$

with

$$
\begin{aligned}
\breve{A} & =A-K(Y) B \\
\breve{Q} & =C Q C^{T}-C S^{T} R^{-1} S C^{T}=C\left(Q-S^{T} R^{-1} S\right) C^{T}
\end{aligned}
$$

and $K(Y)$ as in (3) as a symmetric Stein equation. As $Q-S^{T} R^{-1} S$ is the Schur complement of $R$ in the positive semidefinite matrix $\left[\begin{array}{ll}Q & S^{T} \\ S & R\end{array}\right][34], \breve{Q}$ is symmetric positive semidefinite. Now, denote the solution of (1) obtained by the

Smith iteration (17)-(18) by $Y^{*}$. Fix the right hand side of (19) by evaluating $K\left(Y^{*}\right)=K^{*}$,

$$
\begin{equation*}
E Y E^{T}-\breve{A} Y \breve{A}^{T}=\breve{Q}+\left(C S^{T}-K^{*} R\right) R^{-1}\left(C S^{T}-K^{*} R\right)^{T} . \tag{20}
\end{equation*}
$$

We will use this as a defect correction and update our approximate solution $Y^{*}$ by $\breve{Y}$ which is the solution of

$$
\begin{equation*}
E \breve{Y} E^{T}=\breve{A} \breve{Y} \breve{A}^{T}+\breve{Q}+\left(C S^{T}-K^{*} R\right) R^{-1}\left(C S^{T}-K^{*} R\right)^{T} . \tag{21}
\end{equation*}
$$

This may be done by the Smith iteration (14) where we substitute $A_{k}$ by $\breve{A}$ and $\mathcal{C}_{k}$ by $\mathcal{C}_{k}^{*}=\breve{Q}+\left(C S^{T}-K^{*} R\right) R^{-1}\left(C S^{T}-K^{*} R\right)^{T}$.

A disadvantage of the Smith iteration is that its convergence is linear and the rate is bounded by $\rho_{k}=\max \left\{|\lambda|, \lambda \in \sigma\left(A_{k}, E\right)\right\}$. As $\rho_{k}$ can be close to 1 , this may be very slow. The ADI iteration discussed next offers an alternative.

### 3.2 ADI Iteration

The alternating direction implicit (ADI) iteration was first introduced [47] to solve linear systems arising from the discretization of elliptic boundary value problems. In general, the ADI iteration is used to solve linear systems of the form

$$
M y=b,
$$

where $M$ is symmetric positive definite and can be split into the sum of two positive definite matrices $M=M_{1}+M_{2}$ for which the following iteration is efficient:

$$
\begin{aligned}
y_{0} & =0 \\
\left(M_{1}+\mu_{j} I\right) y_{j-1 / 2} & =b-\left(M_{2}-\mu_{j} I\right) y_{j-1}, \\
\left(M_{2}+\eta_{j} I\right) y_{j} & =b-\left(M_{1}-\eta_{j} I\right) y_{j-1 / 2},
\end{aligned}
$$

for $j=1,2,3, \ldots$ The ADI shift parameters $\mu_{j}$ and $\eta_{j}$ are determined from spectral bounds on $M_{1}$ and $M_{2}$ to improve the convergence rate.

In [17], the ADI iteration is applied to the equation

$$
\begin{equation*}
X-\widehat{A} X \widehat{B}=\widehat{C}, \quad \widehat{A} \in \mathbb{R}^{s \times s}, \widehat{B} \in \mathbb{R}^{t \times t}, \widehat{C} \in \mathbb{R}^{s \times t} \tag{22}
\end{equation*}
$$

Let $X_{0}$ be an initial approximate solution of this equation. Then the ADI iteration generates new iterates $X_{M+N}$ as follows

$$
\begin{align*}
M:=M+1 ; \quad X_{M+N}\left(I-\mu_{M} \widehat{B}\right) & =\left(\widehat{A}-\mu_{M} I\right) X_{M+N-1} \widehat{B}+\widehat{C}  \tag{23}\\
N:=N+1 ; \quad\left(I-\eta_{N} \widehat{A}\right) X_{M+N} & =\widehat{A} X_{M+N-1}\left(\widehat{B}-\eta_{N} I\right)+\widehat{C} \tag{24}
\end{align*}
$$

where one starts with $M=N=0$. Either (23) or (24) can be used to determine $X_{M+N}$ from $X_{M+N-1}$. When (23) and (24) are used in a strictly alternating fashion, the ADI method obtained is analogous to the "classical" ADI iteration for Sylvester's equation discussed by Wachspress [57]. In [17], strict alternation between the formulas (23) and (24) is not required as, e.g., the structures and sizes of the matrices $\widehat{A}$ and $\widehat{B}$ may make the computation of $X_{M+N}$ in one of the equations, say (23) faster than by the other one. In this case the computational effort to solve (22) may be reduced by applying (23) more often than (24), even if such an approach would result in slightly lower convergence than strict alternation. The convergence of the iteration (23)-(24) for different ratios $M / N$ is analyzed in [17] as well as the choice of the iteration parameters. There is still the problem of storing the usually dense $n \times n$ matrix $X_{M+N}$. This storage can be avoided by observing that, for the problems under consideration, the spectrum of the positive semidefinite matrix $X_{M+N}=Z_{M+N} Z_{M+N}^{T}$ often decays to zero rapidly. Here $Z_{M+N}$ can be considered as a Cholesky factor of $X_{M \pm N}$. We expect that $X_{M+N}$ can be approximated accurately by a factorization $\widetilde{Z} \widetilde{Z}^{T}$ for some $\widetilde{Z} \in \mathbb{R}^{n \times r}$ with $r \ll n$.

Here, we will restrict our discussion to a strict alternation between the two equations. With $\widehat{A}=E^{-1} A_{k}, \widehat{B}=A_{k}^{T} E^{-T}$, and $\widehat{C}=E^{-1} \mathcal{R}_{k} E^{-T}$, (22) is equivalent to (8) and we obtain from (23)

$$
X_{i-\frac{1}{2}}\left(I-\mu_{i} A_{k}^{T} E^{-T}\right)=\left(E^{-1} A_{k}-\mu_{i} I\right) X_{i-1} A_{k}^{T} E^{-T}+E^{-1} \mathcal{R}_{k} E^{-T}
$$

while (24) yields

$$
\left(I-\eta_{i} E^{-1} A_{k}\right) X_{i}=E^{-1} A_{k} X_{i-\frac{1}{2}}\left(A_{k}^{T} E^{-T}-\eta_{i} I\right)+E^{-1} \mathcal{R}_{k} E^{-T}
$$

In order to get rid of the inverses of $E$, one can manipulate both equations in an obvious way to reach

$$
\begin{align*}
E X_{i-\frac{1}{2}}\left(E^{T}-\mu_{i} A_{k}^{T}\right) & =\left(A_{k}-\mu_{i} E\right) X_{i-1} A_{k}^{T}+\mathcal{R}_{k}  \tag{25}\\
\left(E-\eta_{i} A_{k}\right) X_{i} E^{T} & =A_{k} X_{i-\frac{1}{2}}\left(A_{k}^{T}-\eta_{i} E^{T}\right)+\mathcal{R}_{k}
\end{align*}
$$

Similarly, with $\widehat{A}=E^{-1} A_{k}, \widehat{B}=A_{k}^{T} E^{-T}$, and $\widehat{C}=E^{-1} \mathcal{C}_{k} E^{-T}$, (22) is equivalent to (9) and we obtain

$$
\begin{align*}
E X_{i-\frac{1}{2}}\left(E^{T}-\mu_{i} A_{k}^{T}\right) & =\left(A_{k}-\mu_{i} E\right) X_{i-1} A_{k}^{T}+E \widehat{C} E^{T}  \tag{26}\\
\left(E-\eta_{i} A_{k}\right) X_{i} E^{T} & =A_{k} X_{i-\frac{1}{2}}\left(A_{k}^{T}-\eta_{i} E^{T}\right)+E \widehat{C} E^{T} .
\end{align*}
$$

While $\mathcal{R}_{k}$ in (25) may be indefinite, $E^{T} \widehat{C} E$ in (26) is positive semidefinite, see (15).

In order to derive a formulation in terms of (Cholesky) factors, let us rewrite the iteration (26) in a single equation. In particular, we have

$$
\begin{aligned}
X_{i-\frac{1}{2}}\left(I-\mu_{i} A_{k}^{T} E^{-T}\right) & =\left(E^{-1} A_{k}-\mu_{i} I\right) X_{i-1} A_{k}^{T} E^{-T}+\widehat{C} \\
\left(I-\eta_{i} E^{-1} A_{k}\right) X_{i} & =E^{-1} A_{k} X_{i-\frac{1}{2}}\left(A_{k}^{T} E^{-T}-\eta_{i} I\right)+\widehat{C} .
\end{aligned}
$$

Hence, with

$$
X_{i-\frac{1}{2}}=\left(E^{-1} A_{k}-\mu_{i} I\right) X_{i-1} A_{k}^{T} E^{-T}\left(I-\mu_{i} A_{k}^{T} E^{-T}\right)^{-1}+\widehat{C}\left(I-\mu_{i} A_{k}^{T} E^{-T}\right)^{-1}
$$

we have

$$
\begin{aligned}
X_{i}= & \left(I-\eta_{i} E^{-1} A_{k}\right)^{-1}\left(E^{-1} A_{k} X_{i-\frac{1}{2}}\left(A_{k}^{T} E^{-T}-\eta_{i} I\right)+\widehat{C}\right) \\
= & \left(I-\eta_{i} E^{-1} A_{k}\right)^{-1} . \\
& \left(E^{-1} A_{k}\left(E^{-1} A_{k}-\mu_{i} I\right) X_{i-1} A_{k}^{T} E^{-T}\left(I-\mu_{i} A_{k}^{T} E^{-T}\right)^{-1}\left(A_{k}^{T} E^{-T}-\eta_{i} I\right)\right. \\
& \left.\quad+E^{-1} A_{k} \widehat{C}\left(I-\mu_{i} A_{k}^{T} E^{-T}\right)^{-1}\left(A_{k}^{T} E^{-T}-\eta_{i} I\right)+\widehat{C}\right) .
\end{aligned}
$$

As

$$
\left(I-\mu_{i} A_{k}^{T} E^{-T}\right)^{-1}\left(A_{k}^{T} E^{-T}-\eta_{i} I\right)=\left(A_{k}^{T} E^{-T}-\eta_{i} I\right)\left(I-\mu_{i} A_{k}^{T} E^{-T}\right)^{-1}
$$

we have

$$
\begin{aligned}
X_{i}=(I- & \left.\eta_{i} E^{-1} A_{k}\right)^{-1}\left(E^{-1} A_{k}\left(E^{-1} A_{k}-\mu_{i} I\right) X_{i-1} A_{k}^{T} E^{-T}\left(A_{k}^{T} E^{-T}-\eta_{i} I\right)\right. \\
& \left.+E^{-1} A_{k} \widehat{C}\left(A_{k}^{T} E^{-T}-\eta_{i} I\right)+\widehat{C}\left(I-\mu_{i} A_{k}^{T} E^{-T}\right)\right)\left(I-\mu_{i} A_{k}^{T} E^{-T}\right)^{-1}
\end{aligned}
$$

Obviously, in order to obtain a symmetric $X_{i}$ we have to restrict the choice of $\mu_{i}$ and $\eta_{i}$ such that $\eta_{i}=\bar{\mu}_{i}$,

$$
\begin{aligned}
X_{i}=(I- & \left.\bar{\mu}_{i} E^{-1} A_{k}\right)^{-1}\left(E^{-1} A_{k}\left(E^{-1} A_{k}-\mu_{i} I\right) X_{i-1}\left(A_{k}^{T} E^{-T}-\bar{\mu}_{i} I\right) A_{k}^{T} E^{-T}\right. \\
& \left.+E^{-1} A_{k} \widehat{C}\left(A_{k}^{T} E^{-T}-\bar{\mu}_{i} I\right)+\widehat{C}\left(I-\mu_{i} A_{k}^{T} E^{-T}\right)\right)\left(I-\mu_{i} E^{-1} A_{k}\right)^{-T}
\end{aligned}
$$

Observe that

$$
\begin{aligned}
& E^{-1} A_{k} \widehat{C}\left(A_{k}^{T} E^{-T}-\bar{\mu}_{i} I\right)+\widehat{C}\left(I-\mu_{i} A_{k}^{T} E^{-T}\right) \\
&= E^{-1} A_{k} \widehat{C} A_{k}^{T} E^{-T}-\bar{\mu}_{i} E^{-1} A_{k} \widehat{C}+\widehat{C}-\mu_{i} \widehat{C} A_{k}^{T} E^{-T} \\
&= E^{-1} A_{k} \widehat{C} A_{k}^{T} E^{-T}+\left(I-\bar{\mu}_{i} E^{-1} A_{k}\right) \widehat{C}\left(I-\mu_{i} A_{k}^{T} E^{-T}\right) \\
&-\left|\mu_{i}\right|^{2} E^{-1} A_{k} \widehat{C} A_{k}^{T} E^{-T} \\
&=\left(1-\left|\mu_{i}\right|^{2}\right) E^{-1} A_{k} \widehat{C} A_{k}^{T} E^{-T}+\left(I-\bar{\mu}_{i} E^{-1} A_{k}\right) \widehat{C}\left(I-\mu_{i} A_{k}^{T} E^{-T}\right) .
\end{aligned}
$$

```
Algorithm 5 ADI-Newton's Method for the DARE
    Input: The coefficient matrices \(A, B, C, E, Q, R, S\) of the DARE (1), and a
    starting guess \(X_{0}\), so that \(\sigma\left(A-K\left(X_{0}\right) B, E\right) \subset \mathcal{D}_{1}(0)\) and \(R+B X_{0} B^{T}>0\).
    Output: An approximate solution \(X_{k+1}\) of the DARE (1).
    for \(k=0,1,2, \ldots\) do
        1. \(K_{k} \leftarrow K\left(X_{k}\right)\).
        2. \(A_{k} \leftarrow A-K_{k} B\).
        3. \(\mathcal{C}_{k} \leftarrow \mathcal{C}\left(X_{k}\right)\).
        4. \(Y_{0} \leftarrow 0\).
        for \(j=1,2, \ldots\) do
            5. Choose shift \(\mu_{j}\)
            6. \(Y_{j} \leftarrow E^{-1} \mathcal{C}_{k} E^{-T}\).
            7. \(H \leftarrow\left(1-\left|\mu_{j}\right|^{2}\right) E^{-1} A_{k} Y_{j} A_{k}^{T} E^{-T}\).
            8. \(H \leftarrow H+E^{-1} A_{k}\left(E^{-1} A_{k}-\mu_{j} I\right) Y_{j-1}\left(A_{k}^{T} E^{-T}-\bar{\mu}_{j} I\right) A_{k}^{T} E^{-T}\).
            9. \(Y_{j} \leftarrow Y_{j}+\left(I-\bar{\mu}_{j} E^{-1} A_{k}\right)^{-1} H\left(I-\mu_{j} E^{-1} A_{k}\right)^{-T}\).
        end for
        10. \(X_{k} \leftarrow Y_{j}\).
    end for
```

Hence, we obtain the single equation

$$
\begin{align*}
X_{i}= & \widehat{C}+\left(I-\bar{\mu}_{i} E^{-1} A_{k}\right)^{-1}\left(\left(1-\left|\mu_{i}\right|^{2}\right) E^{-1} A_{k} \widehat{C} A_{k}^{T} E^{-T}+\right.  \tag{27}\\
& \left.E^{-1} A_{k}\left(E^{-1} A_{k}-\mu_{i} I\right) X_{i-1}\left(A_{k}^{T} E^{-T}-\bar{\mu}_{i} I\right) A_{k}^{T} E^{-T}\right)\left(I-\mu_{i} E^{-1} A_{k}\right)^{-T} .
\end{align*}
$$

If $X_{i-1}$ is positive semidefinite, then so is $X_{i}$.
As

$$
\left(A_{k}^{T} E^{-T}-\bar{\mu}_{i} I\right) A_{k}^{T} E^{-T}\left(I-\mu_{i} E^{-1} A_{k}\right)^{-T}=\left(A_{k}^{T} E^{-T}-\bar{\mu}_{i} I\right)\left(E^{T} A_{k}^{-T}-\mu_{i} I\right)^{-1},
$$

the spectral radius

$$
\rho_{A D I}=\rho\left(\prod_{j=1}^{\ell}\left(A_{k}^{-1} E-\mu_{i} I\right)^{-1}\left(E^{-1} A_{k}-\bar{\mu}_{i} I\right)\right)
$$

determines the rate of convergence, where $\ell$ is the number of iteration steps and shifts used. The minimization of $\rho_{A D I}$ with respect to the shift parameters $\mu_{i}$ is given by

$$
\begin{equation*}
\min _{\mu_{1}, \ldots, \mu_{\ell}} \max _{\lambda \in \sigma\left(E^{-1} A_{k}\right)} \frac{\left|\left(\lambda-\bar{\mu}_{1}\right) \cdots\left(\lambda-\bar{\mu}_{\ell}\right)\right|}{\left|\left(\frac{1}{\lambda}-\mu_{1}\right) \cdots\left(\frac{1}{\lambda}-\mu_{\ell}\right)\right|} \tag{28}
\end{equation*}
$$

In [17] the choice of shifts for the more general ADI iteration (23), (24) applied to (22) is discussed. Although here we are considering a special case of that
situation, the proposed choice of the shifts can not be carried over directly, as the shifts $\mu_{M}$ and $\eta_{N}$ in (23) and (24) are not related to each other, while we require $\mu_{j}=\overline{\eta_{j}}$ in order to obtain a symmetric solution from a symmetric starting guess.

The choice of the shift parameters $\mu_{i}$ is related to a rational approximation problem. The conventional approach to the computation of these parameters in other ADI settings is to cover the spectrum by a domain $\Omega \in \mathbb{C}$ and to solve the minmax problem with respect to $\Omega$ instead of the spectrum. This approximation theory based approach generally requires the knowledge of certain bounds of the spectrum. Heuristic approaches have also been proposed. See, e.g., $[18,49]$ and the references therein.

In Section 4 we propose to use a quite simple, numerically inexpensive, heuristic algorithm which replaces the spectrum by some approximations to the largest and the smallest eigenvalues in the spectrum $\sigma\left(E^{-1} A_{k}\right)$. Assume that the largest $\ell$ and the smallest $\jmath$ eigenvalues have been approximated, $\mathcal{S}=\left\{\widetilde{\lambda}_{1}, \ldots, \widetilde{\lambda}_{\ell}, \widetilde{\lambda}_{n-\jmath+1}, \ldots, \widetilde{\lambda_{n}}\right\}$. Then the minmax problem is replaced by

$$
\begin{equation*}
\min _{\mu_{1}, \ldots, \mu_{\ell}} \max _{\lambda \in \mathcal{S}} \frac{\left|\left(\lambda-\bar{\mu}_{1}\right) \cdots\left(\lambda-\bar{\mu}_{\ell}\right)\right|}{\left|\left(\frac{1}{\lambda}-\mu_{1}\right) \cdots\left(\frac{1}{\lambda}-\mu_{\ell}\right)\right|} . \tag{29}
\end{equation*}
$$

Next a heuristic optimization method is employed in order to compute suitable shifts. For an efficient solution of large-scale DAREs, the solution of (28) will require further studies in the future.

Using (15) we can write

$$
\widehat{C}=\left(E^{-1} U\right)\left(U^{T} E^{-T}\right)=M M^{T}, \quad M \in \mathbb{R}^{n \times(p+m)} .
$$

Finally, with $\left|\mu_{i}\right| \leq 1, X_{i}=Z_{i} Z_{i}^{T}, Z_{i} \in \mathbb{R}^{n \times r_{j}}$ and

$$
\begin{aligned}
\sqrt{1-\left|\mu_{i}\right|^{2}}\left(I-\bar{\mu}_{i} E^{-1} A_{k}\right)^{-1} E^{-1} A_{k} M & =\widehat{M}, \\
\left(I-\bar{\mu}_{i} E^{-1} A_{k}\right)^{-1} E^{-1} A_{k}\left(E^{-1} A_{k}-\mu_{i} I\right) Z_{i-1} & =\widetilde{M},
\end{aligned}
$$

we have

$$
\begin{aligned}
\sqrt{1-\left|\mu_{i}\right|^{2}}\left(E-\bar{\mu}_{i} A_{k}\right)^{-1} A_{k} M & =\widehat{M} \in \mathbb{R}^{n \times(p+m)}, \\
\left(E-\bar{\mu}_{i} A_{k}\right)^{-1} A_{k} E^{-1}\left(A_{k}-\mu_{i} E\right) Z_{i-1} & =\widetilde{M} \in \mathbb{R}^{n \times r_{j-1}} .
\end{aligned}
$$

With this we can write (27) in terms of a factorization

$$
X_{i}=Z_{i} Z_{i}^{T}=M M^{T}+\widehat{M} \widehat{M}^{T}+\widetilde{M} \widetilde{M}^{T}=\left[\begin{array}{lll}
M & \widehat{M} & \widetilde{M}
\end{array}\right]\left[\begin{array}{c}
M^{T} \\
\widehat{M}^{T} \\
\widetilde{M}^{T}
\end{array}\right] .
$$

Hence,

$$
Z_{i}=\left[\begin{array}{lll}
M & \widehat{M} & \widetilde{M}
\end{array}\right] \in \mathbb{R}^{n \times\left(2(p+m)+r_{j}\right)}
$$

is one possible representation of the (Cholesky) factor of the solution $X_{i}$. The dimension of $Z_{j}$ would increase by $2(p+m)$ in each iteration step. But if $Z_{j}$ is of low rank $r_{j}$, then we can approximate it using a rank-revealing $L Q$ factorization as before. The derivation of the algorithm is straightforward, hence we omit its statement.

## 4 Numerical Examples

All numerical tests were done in Matlab R2009a on a Pentium M notebook with 512 MB main memory. The iterations were stopped in all algorithms as soon as $\left\|Z_{j+1} Z_{j+1}^{T}-Z_{j} Z_{j}^{T}\right\|_{F} /\left\|Z_{j+1} Z_{j+1}^{T}\right\|_{F}$ and $\left\|T_{k+1} T_{k+1}^{T}-T_{k} T_{k}^{T}\right\|_{F} /\left\|T_{k+1} T_{k+1}^{T}\right\|_{F}$ were less than a given tolerance.

The examples considered are optimal control problems of the following form:

$$
\min \mathcal{J}(u), \quad \mathcal{J}(u)=\frac{1}{2} \sum_{k=0}^{\infty} y_{k}^{T} Q y_{k}+u_{k}^{T} R u_{k}+2 y_{k}^{T} S u_{k}
$$

subject to a fully discretized parabolic PDE constraint

$$
\begin{align*}
E x_{k+1} & =A x_{k}+B u_{k},  \tag{30}\\
y_{k} & =C x_{k} . \tag{31}
\end{align*}
$$

In all examples we have $m=p=1$ and choose $R=1, Q=1, S=0$.
In order to compare the computed solutions to an 'exact' one, the Matlab routine dare was used to produce the "exact" solution $X_{\text {exact }}$.

In our examples we used only one shift for the ADI iteration. The minmax problem (28) is replaced by

$$
\begin{equation*}
\min _{\mu} \max _{\lambda \in \mathcal{S}} \frac{|\lambda-\bar{\mu}|}{\left|\frac{1}{\lambda}-\mu\right|}, \tag{32}
\end{equation*}
$$

where $\mathcal{S}$ is either the set of all eigenvalues computed by eig or eigenvalue approximations computed by eigs. Assume that the largest $\ell$ and the smallest $\jmath$ eigenvalues have been approximated so that $\mathcal{S}=\left\{\widetilde{\lambda}_{1}, \ldots, \widetilde{\lambda}_{\ell}, \widetilde{\lambda}_{n-\jmath+1}, \ldots, \widetilde{\lambda}_{n}\right\}$. Next, fminsearch is employed as a heuristic optimization method in order to compute a suitable shift which we simply call $\mu_{*}$ in the following. The standard setting for all examples presented is $\ell=8, j=2$. Moreover, eigs was called with OPTS.tol $=10^{-2}$ as more accuracy in this computation does not increase the accuracy of the computed solution or convergence rate of the algorithm. Note that an optimization with respect to the full spectral set did in no case improve
the convergence of the overall algorithms, i.e., the computed $\ell+j$ Ritz values were sufficient to obtain a suitable shift. Thus, we only report results based on $\mathcal{S}$ being composed of the computed $\ell+j$ Ritz values. When further developing the algorithms proposed in this paper, a multi-shift selection strategy needs to be developed. A first candidate for this would be Penzl's heuristic [49] that often yields good results in the ADI iteration for continuous Lyapunov equations.

Example 4.1 The data of this example come from the autonomous linear-quadratic optimal control problem of one dimensional heat flow and is taken from [1, Example 4.2]. The model equations are

$$
\begin{aligned}
\frac{\partial}{\partial t} x(t, \eta) & =\alpha \frac{\partial^{2}}{\partial \eta^{2}} x(t, \eta)+b(\eta) u(t) & & \eta \in(0,1) ; t>0 \\
x(t, 0) & =0=x(t, 1) & & t>0 \\
x(0, \eta) & =x_{0}(\eta) & & \eta \in(0,1) \\
y(x) & =\int_{0}^{1} c(\eta) x(t, \eta) d \eta, & & x>0,
\end{aligned}
$$

where $\alpha>0$ and $x_{0}, b, c \in L_{2}(0,1)$. Using a standard linear finite element (FE) discretization in space, one obtains an ordinary differential equation

$$
M_{N} \dot{x}(t)=K_{N} x(t)+F u(t), \quad y(t)=C x(t),
$$

where

$$
M_{N}=\frac{1}{6 N}\left[\begin{array}{ccccc}
4 & 1 & & & \\
1 & 4 & 1 & & \\
& \ddots & \ddots & \ddots & \\
& & 1 & 4 & 1 \\
& & 1 & 4
\end{array}\right], K_{N}=-\alpha N\left[\begin{array}{ccccc}
2 & -1 & & & \\
-1 & 2 & -1 & & \\
& \ddots & \ddots & \ddots & \\
& & -1 & 2 & -1 \\
& & -1 & 2
\end{array}\right]
$$

$F=b_{N}$ and $C=c_{N}^{T}, M_{N}, K_{N} \in \mathbb{R}^{N \times N}$. The vectors $b_{N}$ and $c_{N}$ are obtained from the $L_{2}$ inner products of indicator functions for subsets of $(0,1)$ with the FE basis functions $\left\{\phi_{i}^{N}\right\}_{i=1}^{N-1}$ :

$$
\left(b_{N}\right)_{j}=\int_{0}^{1} \beta(s) \phi_{i}^{N}(s) d s, \quad \text { and } \quad\left(c_{N}\right)_{j}=\int_{0}^{1} \gamma(s) \phi_{i}^{N}(s) d s, \quad j=1, \ldots, N-1,
$$

where the functions $\beta, \gamma \in L_{2}(0,1)$ are given by

$$
\beta(s)=\gamma(s)=\left\{\begin{array}{lc}
1, & s \in[0.1,0.5] \\
0, & \text { otherwise }
\end{array}\right.
$$

(Note: $\beta, \gamma$ are parameter-dependent in [1] and may differ - here we chose them equal for simplicity.) Employing the semi-implicit Euler method with stepsize $\Delta t$ yields the difference equation

$$
(M-\Delta t K) x_{k+1}=M x_{k}+\Delta t F u_{k},
$$

that is, $A=M, E=M-\Delta t K$, and $B=\Delta t F$ in (30).
In the computational experiments reported in Tables $1-2$, we set $\alpha=0.05$ (default in [1]). fminsearch was started with different starting guesses $\mu_{0}$. The resulting shift is given as $\mu_{*}$. Different $n$, tol and $h$ have been considered. Here only the results for $n=1000, h=0.1$ and $h=0.01$ as well as $t o l=10^{-8}$ and tol $=10^{-12}$ are presented. Here and in the following, $\rho(A, E)$ denotes the spectral radius of $A-\lambda E$, i.e., $\rho(A, E)=\max _{\lambda \in \sigma(A, E)}\{|\lambda|\}$.

No matter how the parameters were chosen, all methods needed 4 Newton steps. The number of Smith iteration steps was usually quite high, about twice as many as needed for the unshifted ADI iteration. Obviously, neither the Smith iteration nor the unshifted ADI iteration yield feasible methods for solving the DARE with Newton's method in this example. On the other hand, using our fairly simple heuristic for choosing a single shift, the ADI iteration with shift converges up to 100 times faster than the ADI iteration without a shift. The number of iterations needed varied from Newton step to Newton step, being larger at first.

Changing $n$ while keeping the rest of the parameters fixed, does not change the number of iterations needed for convergence or the accuracy achieved.

| $\Delta t=0.1, \rho(A, E) \approx 0.9902$ |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
|  | $X_{\text {dare }}$ | $Z_{\text {smith }}$ | $Z_{\text {adi }}$ | $Z_{\text {adi }}$ |
|  |  |  | $\mu_{*}=0.96328$ | no shift |
| Newton steps |  | 4 | 4 | 4 |
| Smith/ADI steps |  | $711-719$ | $17-18$ | $374-378$ |
| $\mathcal{R}(X)$ | $2.2 \cdot 10^{-12}$ | $7.1 \cdot 10^{-12}$ | $2.1 \cdot 10^{-13}$ | $3.5 \cdot 10^{-12}$ |
| $\left\\|X_{\text {dare }}-X\right\\|_{F} /\\|X\\|_{F}$ |  | $4.9 \cdot 10^{-7}$ | $4.6 \cdot 10^{-9}$ | $2.4 \cdot 10^{-7}$ |
| $\operatorname{rank}\left(Z Z^{T}\right)$ | 13 | 11 | 11 | 11 |


| $\Delta t=0.01, \rho(A, E) \approx 0.999$ |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
|  | $X_{\text {dare }}$ | $Z_{\text {smith }}$ | $Z_{\text {adi }}$ | $Z_{\text {adi }}$ |
|  |  |  | $\mu_{*}=0.99306$ | no shift |
| Newton steps |  | 4 | 4 | 4 |
| Smith/ADI steps |  | $5926-5990$ | $31-35$ | $3138-3172$ |
| $\mathcal{R}(X)$ | $2.3 \cdot 10^{-13}$ | $7.2 \cdot 10^{-11}$ | $1.3 \cdot 10^{-12}$ | $3.6 \cdot 10^{-11}$ |
| $\left\\|X_{\text {dare }}-X\right\\|_{F} /\\|X\\|_{F}$ |  | $5.0 \cdot 10^{-6}$ | $1.1 \cdot 10^{-8}$ | $2.5 \cdot 10^{-6}$ |
| $\operatorname{rank}\left(Z Z^{T}\right)$ | 16 | 16 | 16 | 16 |

Table 1: Example 1, tol $=10^{-8}, n=1000, \mu_{0}=-1$.

Example 4.2 Here, we use again a discretized control problem for the heat equation, this time taken from the SLICOT Model Reduction Benchmark Collection [19]. (The eigenvalues of the DARE solution are shown in Figure 1.) In contrast to Example 4.1, this time a finite differences discretization is used,

| $\Delta t=0.1$ |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
|  | $X_{\text {dare }}$ | $Z_{\text {smith }}$ | $\begin{gathered} Z_{\text {adi }} \\ \mu_{*}=0.96325 \end{gathered}$ | $\begin{gathered} Z_{\text {adi }} \\ \text { no shift } \end{gathered}$ |
| $\begin{gathered} \text { Newton steps } \\ \text { Smith/ADI steps } \\ \mathcal{R}(X) \\ \left\\|X_{\text {dare }}-X\right\\|_{F} /\\|X\\|_{F} \\ \operatorname{rank}\left(Z Z^{T}\right) \\ \hline \end{gathered}$ | $\begin{gathered} 2.2 \cdot 10^{-12} \\ 13 \end{gathered}$ | $\begin{gathered} \hline 4 \\ 1175-1188 \\ 1.1 \cdot 10^{-13} \\ 6.1 \cdot 10^{-11} \\ 11 \end{gathered}$ | 4 26 $1.1 \cdot 10^{-13}$ $1.3 \cdot 10^{-11}$ 11 | $\begin{gathered} \hline 4 \\ 606-613 \\ 1.1 \cdot 10^{-13} \\ 3.6 \cdot 10^{-11} \\ 11 \end{gathered}$ |
| $\Delta t=0.01$ |  |  |  |  |
|  | $X_{\text {dare }}$ | $Z_{\text {smith }}$ | $\begin{gathered} Z_{\text {adi }} \\ \mu_{*}=0.99306 \\ \hline \end{gathered}$ | $\begin{gathered} Z_{\text {adi }} \\ \text { no shift } \end{gathered}$ |
| $\begin{gathered} \text { Newton steps } \\ \text { Smith/ADI steps } \\ \mathcal{R}(X) \\ \left\\|X_{\text {dare }}-X\right\\|_{F} /\\|X\\|_{F} \\ \operatorname{rank}\left(Z Z^{T}\right) \\ \hline \end{gathered}$ | $\begin{gathered} 2.3 \cdot 10^{-13} \\ 16 \end{gathered}$ | $\begin{gathered} 4 \\ 10544-10658 \\ 1.6 \cdot 10^{-14} \\ 4.7 \cdot 10^{-10} \\ 16 \end{gathered}$ | 4 48 $1.5 \cdot 10^{-14}$ $3.2 \cdot 10^{-11}$ 16 | $\begin{gathered} 4 \\ 5447-5506 \\ 1.5 \cdot 10^{-14} \\ 2.2 \cdot 10^{-10} \\ 16 \end{gathered}$ |

Table 2: Example 1, tol $=10^{-12}, n=1000, \mu_{0}=0$.
where the spatial domain is discretized into segments of length $h=\frac{1}{N+1}$. Suppose for example that one wants to heat in a point of the rod located at $1 / 3$ of the length and wants to record the temperature at $2 / 3$ of the length. We obtain the semi-discretized system:

$$
\begin{aligned}
\dot{x}(t) & =K x(t)+F u(t), \quad x(0)=0 \\
y(t) & =C x(t)
\end{aligned}
$$

where
$K=-\frac{\alpha}{h^{2}}\left[\begin{array}{ccccc}2 & -1 & & & \\ -1 & 2 & -1 & & \\ & \ddots & \ddots & \ddots & \\ & & -1 & 2 & -1 \\ & & & -1 & 2\end{array}\right] \in \mathbb{R}^{N \times N}, F=e_{(N+1) / 3} \in \mathbb{R}^{N}, C=e_{2 N / 3} \in \mathbb{R}^{N}$,
and $x(t) \in \mathbb{R}^{N}$ is the solution evaluated at each $x$ value in the discretization for $t$. Here, $e_{k}$ denotes the $k$ th unit vector. Now if we want to completely discretize the system, for example using Crank-Nicholson we obtain:

$$
\begin{aligned}
E x_{k+1} & =A x_{k}+B u_{k}, \quad x_{0}=0 \\
y_{k} & =C x_{k}
\end{aligned}
$$

where $E=I_{N}-\frac{\Delta t}{2} K, A=I_{N}+\frac{\Delta t}{2} K, B=\Delta t F$, and $\Delta t$ is the time step.

The data provided in the benchmark collection [19] yields matrices $K, F$ and $C$ with $N=200, m=p=1$.

As $\rho(A, E) \approx 0.9975$, the iteration for the Stein equation converges quite slowly, for the results see Table 3.

When using the Smith iteration on this example, the iterations converged after 3 Newton iteration, while in each iteration step the Smith iteration took 1343 steps. The rank of $X_{\text {exact }}$ is 31 , the rank of the computed low-rank factorization factor $Z$ is 39 , the rank of $Z Z^{T}$ is 25 . When looking at the singular values of $X_{\text {exact }}$ and $Z Z^{T}$, one can see that this is a numerical accuracy problem, the first about 25 singular values are of the order of $10^{-8}$ to $10^{-14}$, while the rest is even smaller. Without the low-rank reduction of $Z_{j}$ in every step, the factor $Z_{j}$ would increase to the size $\mathbb{R}^{200 \times 2688}$, while due to the low-rank reduction, the factor $Z_{j}$ is at most of size $\mathbb{R}^{200 \times 39}$. The ADI iteration converged much quicker than the Smith iteration, moreover the use of a shift allowed the computation of a more accurate solution.

In order to obtain a better convergence behavior for the Stein equation, we use the same data, but set $K=0.5 \cdot K$. This reduces $\rho(A, E)$ to $\approx 0.4987$. As can be seen from Table 3, the iteration for the Stein equation converges much faster for all algorithms considered. Here, our heuristic to choose the shift does not work very well, the computed shift is close to 0 which is almost the same as using no shift.

Using a smaller tol does give more accurate results, but no significant chance in the number of iterations needed.

Example 4.3 For the final example, we set up

$$
B=C=[v, \ldots, v] \in \mathbb{R}^{100}
$$

with $v=[1,1,0, \ldots,]^{T} \in \mathbb{R}^{10}$ and

$$
A=I, K=\operatorname{tridiag}(D, T, D) \in \mathbb{R}^{100 \times 100}
$$

with $D=\operatorname{diag}(-121, \ldots,-121) \in \mathbb{R}^{10 \times 10}$ and

$$
T=\left[\begin{array}{cccc}
484 & -115.5 & & \\
-126.5 & \ddots & \ddots & \\
& \ddots & \ddots & -115.5 \\
& & -126.5 & 484
\end{array}\right]=\operatorname{tridiag}(-126.5,484,-115.5) \in \mathbb{R}^{10 \times 10}
$$

The $K$ matrix is constructed using the Matlab function fdm2D_matrix from Lyapack [50]. We include this fairly small example, as it shows that the number of Newton steps also may depend on the algorithm used for the inner iteration to solve the Stein equations. This number is slightly larger for the Smith iteration than for the ADI iteration with shift.

| original $K, \rho(A, E) \approx 0.9975$ |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
|  | $X_{\text {dare }}$ | $Z_{\text {smith }}$ | $Z_{\text {adi }}$ | $Z_{\text {adi }}$ <br> no shift |
| Newton steps |  | 3 | $\mu_{*}=-0.99197$ | 3 |
| Smith/ADI steps |  | 1343 | 85 | 737 |
| $\mathcal{R}(X)$ | $1.8 \cdot 10^{-12}$ | $2.6 \cdot 10^{-4}$ | $3.4 \cdot 10^{-10}$ | $1.3 \cdot 10^{-4}$ |
| $\left\\|X_{\text {dare }}-X\right\\|_{F} /\\|X\\|_{F}$ |  | $1.8 \cdot 10^{-6}$ | $3.7 \cdot 10^{-8}$ | $9.4 \cdot 10^{-7}$ |
| $\operatorname{rank}(X) / \operatorname{rank}\left(Z Z^{T}\right)$ | 31 | 25 | 25 | 26 |
| $K=0.5 K, \rho(A, E) \approx 0.4987$ |  |  |  |  |
|  | $X_{\text {dare }}$ | $Z_{\text {smith }}$ | $Z_{\text {adi }}$ | $Z_{\text {adi }}$ |
| Newton steps |  | 3 | $\mu_{*}=-0.45332$ | no shift |
| Smith/ADI steps |  | 12 | 3 | 3 |
| $\mathcal{R}(X)$ | $2.8 \cdot 10^{-12}$ | $1.1 \cdot 10^{-8}$ | $3.3 \cdot 10^{-12}$ | 7 |
| $\left\\|X_{\text {dare }}-X\right\\|_{F} /\\|X\\|_{F}$ |  | $5.5 \cdot 10^{-10}$ | $2.2 \cdot 10^{-10}$ | $1.3 \cdot 10^{-9}$ |
| $\operatorname{rank}(X) / \operatorname{rank}\left(Z Z^{T}\right)$ | 16 | 10 | 10 | 10 |

Table 3: Example 2, $n=200$, tol $=10^{-8}, \mu_{0}=0$.

|  | $X_{\text {dare }}$ | $Z_{\text {smith }}$ | $Z_{\text {adi }}$ |
| :---: | :---: | :---: | :---: |
|  |  |  | $\mu_{*}=-0.85361$ |
| Newton steps |  | 10 | 8 |
| Smith/ADI steps |  | $10-128$ | $8-12$ |
| $\mathcal{R}(X)$ | $1.3 \cdot 10^{-11}$ | $1.1 \cdot 10^{-11}$ | $1.6 \cdot 10^{-12}$ |
| $\left\\|X_{\text {dare }}-X\right\\|_{F} /\\|X\\|_{F}$ |  | $4.3 \cdot 10^{-12}$ | $2.2 \cdot 10^{-14}$ |
| $\operatorname{rank}(X) / \operatorname{rank}\left(Z Z^{T}\right)$ | 6 | 6 | 6 |

Table 4: Example 3, tol $=10^{-12}, n=100, \Delta t=0.01, \mu_{0}=-1$.

## 5 Conclusions

This paper addresses the numerical solution of large, sparse DAREs based on the Newton method. Its primary step involves the solution of large, sparse, stable Stein equations. We have presented two iterative methods which deliver low-rank approximations to the desired solution, a Smith and an ADI iteration. The ADI iteration can be accelerated significantly by introducing suitable shift parameters. We presented a simple heuristic algorithm for determining a set of ADI parameters. Finally, the algorithms are used to numerically solve an optimal control problem for parabolic PDEs.

Future work will include a detailed study of the problem of choosing the ADI parameters as well as the adaptation of the precision to which the Stein equation is solved to that of the Newton recursion as in [24] done for continuous-time
algebraic Riccati equations.

## Acknowledgments

The first author's work was supported by the DFG project Numerische Lösung von Optimalsteuerungsproblemen für instationäre Diffusions-Konvektions- und Diffusions-Reaktionsgleichungen, grant BE3715/1-1.

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