# Implementation of a restarted Krylov subspace method for the evaluation of matrix functions ${ }^{\text {* }}$ 

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Dedicated to Richard S. Varga on the occasion of his 80th birthday.


#### Abstract

A new implementation of restarted Krylov subspace methods for evaluating $f(A) \boldsymbol{b}$ for a function $f$, a matrix $A$ and a vector $\boldsymbol{b}$ is proposed. In contrast to an implementation proposed previously, it requires constant work and constant storage space per restart cycle. The convergence behavior of this scheme is discussed and a new stopping criterion based on an error indicator is given. The performance of the implementation is illustrated for three parabolic initial value problems, requiring the evaluation of $\exp (A) \boldsymbol{b}$. © 2008 Elsevier Inc. All rights reserved.


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## 1. Introduction

The interplay of complex approximation theory and matrix computations has long been, and still is, a recurring theme in the work of Richard Varga. The subject of this paper is an instance where this interplay is fundamental, namely the computation of the vector

$$
\begin{equation*}
f(A) \boldsymbol{b} \tag{1}
\end{equation*}
$$

[^0]given a matrix $A \in \mathbb{C}^{n \times n}$, a vector $\boldsymbol{b} \in \mathbb{C}^{n}$ of unit norm and a function $f$ analytic in a neighborhood of the spectrum $\Lambda(A)$ of $A$. In particular, we shall find that the successful implementation of a technique which two of the authors recently proposed [6] rests on pioneering work of Richard Varga [3] from the late 1960s.

The evaluation of $f(A) \boldsymbol{b}$ is a common task in scientific computing, and the most familiar case is surely that of the exponential function $f(\lambda)=\exp (\lambda)$, which occurs, e.g., in connection with the linear initial value problem

$$
\begin{equation*}
\boldsymbol{u}^{\prime}(t)=A \boldsymbol{u}(t), \quad t>0, \boldsymbol{u}(0)=\boldsymbol{u}_{0} \tag{2}
\end{equation*}
$$

with solution $\exp (t A) \boldsymbol{u}_{0}$. Initial value problems such as (2) result naturally from the method of lines discretization of parabolic partial differential operators. Such discretizations as well as the construction of time-stepping schemes based on Padé or Chebyshev rational approximation were among the subjects of Richard Varga's earlier work [26].

We are concerned with the situation where $A$ is large and either sparse or structured, such that matrix-vector products with $A$ can be carried out inexpensively, whereas first forming $f(A)$ and then multiplying with $\boldsymbol{b}$ cannot. Here Krylov subspace approximations of (1) have become popular (cf. $[4,9,14]$ ) and, with regard to solving initial value problems, have had a large impact on so-called exponential integrators, in which evaluations of the exponential function applied to the Jacobian are incorporated directly into time-stepping schemes (cf. [15]).

Krylov subspace approximations of (1) are based on an Arnoldi-like decomposition

$$
\begin{equation*}
A V_{m}=V_{m+1} \widetilde{H}_{m}=V_{m} H_{m}+\eta_{m+1, m} \boldsymbol{v}_{m+1} \boldsymbol{e}_{m}^{T} \tag{3}
\end{equation*}
$$

of the matrix $A$ in which the columns of $V_{m}$ form an ascending basis of the Krylov space

$$
\mathscr{K}_{m}(A, \boldsymbol{b}):=\operatorname{span}\left\{\boldsymbol{b}, A \boldsymbol{b}, \ldots, A^{m-1} \boldsymbol{b}\right\},
$$

$\widetilde{H}_{m}:=\left[\eta_{i, j}\right]$ is an $(m+1) \times m$ upper Hessenberg matrix, $H_{m}:=\left[I_{m} \mathbf{0}\right] \widetilde{H}_{m}$ and $\boldsymbol{e}_{m}$ denotes the $m$ th unit coordinate vector in $\mathbb{R}^{m}$. In the most common situation, (3) is a proper Arnoldi decomposition, i.e., the basis sequence $\left\{\boldsymbol{v}_{m}\right\}$ consists of orthonormal vectors, as generated by the Arnoldi process, which reduces to the Hermitian Lanczos process when $A$ is Hermitian. In this case $H_{m}$ is Hermitian tridiagonal. More general Arnoldi-like decompositions arise in restarted Arnoldi schemes, in which the basis vectors are only block orthogonal, or when non-orthogonalizing basis generation schemes are used (see the discussion in [6]). For an orthonormal basis sequence the Arnoldi approximation of (1) is then given by

$$
\begin{equation*}
\boldsymbol{f}_{m}:=V_{m} f\left(H_{m}\right) V_{m}^{H} \boldsymbol{b}=V_{m} f\left(H_{m}\right) \boldsymbol{e}_{1} \tag{4}
\end{equation*}
$$

where $\boldsymbol{e}_{1}$ denotes the first unit coordinate vector in $\mathbb{R}^{m}$. There are (at least) three ways of motivating the approximation (4) (see [13]). One is as a subspace approximation of (1), since for a proper Arnoldi decomposition (3) the Hessenberg matrix $H_{m}=V_{m}^{H} A V_{m}$ represents the orthogonal section of $A$ onto $\mathscr{K}_{m}(A, \boldsymbol{b})$, and in this sense $f\left(H_{m}\right)$ an approximation of the action of $f(A)$ on this space. Alternatively, if $\Gamma$ is a contour containing $\Lambda(A)$ in its interior int $\Gamma$ and if $f$ is analytic in the closure int $\bar{\Gamma}$, we have

$$
\begin{aligned}
\boldsymbol{f}_{m}=V_{m} f\left(H_{m}\right) V_{m}^{H} \boldsymbol{b} & =\frac{1}{2 \pi \mathrm{i}} \int_{\Gamma} f(\lambda) V_{m}\left(\lambda I-H_{m}\right)^{-1} V_{m}^{H} \boldsymbol{b} \mathrm{~d} \lambda \\
& \approx \frac{1}{2 \pi \mathrm{i}} \int_{\Gamma} f(\lambda)(\lambda I-A)^{-1} \boldsymbol{b} \mathrm{~d} \lambda=f(A) \boldsymbol{b},
\end{aligned}
$$

which characterizes $\boldsymbol{f}_{m}$ as the Galerkin approximation of the resolvent integral representation of (1). Finally, and this is the interpretation closest in spirit to the method proposed below, it can be shown that (4) can be written as

$$
\boldsymbol{f}_{m}=q(A) \boldsymbol{b}
$$

where $q \in \mathscr{P}_{m-1}$ is the polynomial of degree $m-1$ which interpolates $f$ in the Hermite sense at the eigenvalues of $H_{m}$, i.e., at the Ritz values of $A$ with respect to $\mathscr{K}_{m}(A, \boldsymbol{b})$.

The evaluation of (4) requires that the complete basis of $\mathscr{K}_{m}(A, \boldsymbol{b})$ be available, which can be prohibitive for large $A$ and large values of $m$. One remedy restricted to the Hermitian case, in which the $\boldsymbol{v}_{m}$ can be generated by the three-term Lanczos-recurrence, is to compute the matrix $H_{m}$ in a first pass to subsequently determine the coefficient vector $f\left(H_{m}\right) e_{1}$, followed by a regeneration of the basis vectors $\boldsymbol{v}_{1}, \boldsymbol{v}_{2}, \ldots, \boldsymbol{v}_{m}$ in a second pass to form the linear combination (4), which, while feasible, seems nonetheless unelegant.

Reducing storage and computational requirements was the motivation behind the restarting algorithm proposed in [6], which involves Arnoldi decompositions of a fixed length $m$ and generates a sequence of approximations $\left\{\hat{\boldsymbol{f}}_{k}\right\}_{k \in \mathbb{N}}$ with $\hat{\boldsymbol{f}}_{k} \in \mathscr{K}_{k m}(A, \boldsymbol{b})$. The implementation proposed in [6], however, suffered from the deficiency that, although not requiring the storage of more than $m$ basis vectors at a time, the coordinate calculations required the evaluation of $f$ for a block Hessenberg matrix of size $k m$, resulting in computational work which grows with $k$. In what follows, we introduce a new implementation for which the computational expense as well as the storage requirements are the same for each cycle. In addition, we propose a stopping criterion based on an error indicator, discuss the convergence behavior of the new implementation vs. that proposed in [6], and illustrate their performance in some numerical examples.

## 2. The restarted Krylov subspace algorithm for matrix functions

The restarted Krylov subspace algorithm proposed in [6] proceeds by repeatedly generating a basis of Krylov spaces of fixed dimension $m$, updating the most recent approximation to (1), and then discarding all but the last basis vector, which is subsequently used as the initial vector for the next Krylov space. Although any procedure which generates a nested basis of a Krylov space can be used, we restrict ourselves to the Arnoldi process in the following. ${ }^{1}$ We recall the basic restart step using two Arnoldi decompositions

$$
\begin{align*}
& A V_{1}=V_{1} H_{1}+\eta_{2} \boldsymbol{v}_{m+1} \boldsymbol{e}_{m}^{T}  \tag{5a}\\
& A V_{2}=V_{2} H_{2}+\eta_{3} \boldsymbol{v}_{2 m+1} \boldsymbol{e}_{m}^{T} \tag{5b}
\end{align*}
$$

in which the columns of $V_{1}$ and $V_{2}$ are orthonormal bases of $\mathscr{K}_{m}\left(A, \boldsymbol{v}_{1}\right)$ and $\mathscr{K}_{m}\left(A, \boldsymbol{v}_{m+1}\right)$, respectively, and $H_{1}$ and $H_{2}$ are unreduced upper Hessenberg matrices. The columns of $\widehat{V}_{2}:=\left[\begin{array}{ll}V_{1} & V_{2}\end{array}\right]$ thus form a basis of $\mathscr{K}_{2 m}\left(A, \boldsymbol{v}_{1}\right)$, and we combine the two decompositions in (5) as

$$
\begin{equation*}
A \widehat{V}_{2}=\widehat{V}_{2} \widehat{H}_{2}+\eta_{3} \boldsymbol{v}_{2 m+1} \boldsymbol{e}_{2 m}^{T} \tag{6}
\end{equation*}
$$

where

$$
\widehat{H}_{2}:=\left[\begin{array}{ll}
H_{1} & O \\
E_{2} & H_{2}
\end{array}\right], \quad E_{2}:=\eta_{2} e_{1} \boldsymbol{e}_{m}^{T}
$$

[^1]Note that, since the columns of $\widehat{V}_{2}$ are only blockwise orthonormal, we refer to (6) as an Arnoldi-like decomposition. The restarting scheme computes the Krylov subspace approximation

$$
\hat{\boldsymbol{f}}_{2}:=\widehat{V}_{2} f\left(\widehat{H}_{2}\right) \boldsymbol{e}_{1}
$$

associated with the decomposition (6). ${ }^{2}$
Due to the block lower triangular structure of $\widehat{H}_{2}$, the approximation $\hat{\boldsymbol{f}}_{2}$ has the form

$$
\hat{\boldsymbol{f}}_{2}=V_{1} f\left(H_{1}\right) \boldsymbol{e}_{1}+V_{2} F_{2,1} \boldsymbol{e}_{1}=\hat{\boldsymbol{f}}_{1}+V_{2} F_{2,1} \boldsymbol{e}_{1}
$$

in which $\hat{\boldsymbol{f}}_{1}$ denotes the approximation associated with the first decomposition (5a) and the $m \times m$ matrix $F_{2,1}$ is defined by

$$
f\left(\widehat{H}_{2}\right)=\left[\begin{array}{cc}
f\left(H_{1}\right) & O \\
F_{2,1} & f\left(H_{2}\right)
\end{array}\right] .
$$

If $F_{2,1}$, or rather its first column, can be computed, then only $\hat{\boldsymbol{f}}_{1}$ needs to be stored from the first cycle of the algorithm, and $V_{1}$ can be discarded after computing $\hat{f}_{1}$.

The result of $k$ cycles of this restarting scheme is the Krylov subspace approximation associated with the decomposition

$$
\begin{equation*}
A \widehat{V}_{k}=\widehat{V}_{k} \widehat{H}_{k}+\eta_{k+1} \boldsymbol{v}_{k m+1} \boldsymbol{e}_{k m}^{T} \tag{7}
\end{equation*}
$$

where $\widehat{V}_{k}:=\left[\begin{array}{lll}V_{1} & V_{2} & \cdots V_{k}\end{array}\right] \in \mathbb{C}^{n \times k m}$,

$$
\widehat{H}_{k}:=\left[\begin{array}{cccc}
H_{1} & & & \\
E_{2} & H_{2} & & \\
& \ddots & \ddots & \\
& & E_{k} & H_{k}
\end{array}\right] \in \mathbb{C}^{k m \times k m}, \quad E_{j}:=\eta_{j} \boldsymbol{e}_{1} \boldsymbol{e}_{m}^{T} \in \mathbb{R}^{m \times m}, j=2, \ldots, k,
$$

in which we have collected the quantities of the $k$ Arnoldi decompositions

$$
A V_{j}=V_{j} H_{j}+\eta_{j+1} \boldsymbol{v}_{j m+1} \boldsymbol{e}_{m}^{T}, \quad j=1,2, \ldots, k
$$

Setting

$$
\widehat{F}_{k}:=f\left(\widehat{H}_{k}\right)=\left[\begin{array}{cccc}
F_{1,1} & & & \\
F_{2,1} & F_{2,2} & & \\
\vdots & \vdots & \ddots & \\
F_{k, 1} & F_{k, 2} & \ldots & F_{k, k}
\end{array}\right], \quad \text { where } F_{j, j}=f\left(H_{j}\right), j=1,2, \ldots, k
$$

the approximation after $k$ restart cycles is given by

$$
\begin{equation*}
\hat{\boldsymbol{f}}_{k}:=\widehat{V}_{k} f\left(\widehat{H}_{k}\right) \boldsymbol{e}_{1}=\left[V_{1} V_{2} \cdots V_{k}\right] \widehat{F}_{k} \boldsymbol{e}_{1}=\sum_{j=1}^{k} V_{j} F_{j, 1} \boldsymbol{e}_{1}=\hat{\boldsymbol{f}}_{k-1}+V_{k} F_{k, 1} \boldsymbol{e}_{1} \tag{8}
\end{equation*}
$$

## 3. Implementation

In this section we discuss possible implementations of the basic restarting scheme (8). The crucial issue is the fast and stable computation of the coefficient vector $F_{k, 1} e_{1}$.

[^2]
### 3.1. Previously explored strategies

There are several possible ways of computing the update (8). In view of the interpolation properties of the Krylov subspace approximation, each additional restart cycle interpolates the function $f$ at $m$ additional nodes, which are the Ritz values of $A$ with respect to the most recent Krylov space. A natural approach would therefore be a block Newton type interpolation scheme, which can be carried out by evaluating matrix polynomials of $m \times m$ matrices: It was shown in ([6], Theorem 2.6) that the error of a Krylov subspace approximation (4) with respect to $\mathscr{K}_{m}(A, \boldsymbol{b})$ and an Arnoldi-like decomposition (3) has the representation

$$
\begin{equation*}
f(A) \boldsymbol{b}-\boldsymbol{f}_{m}=\tilde{f}(A) \boldsymbol{v}_{m+1} \tag{9}
\end{equation*}
$$

with a "restart function" $\tilde{f}:=\gamma \Delta_{w_{m}} f$, in which $w_{m}$ denotes the characteristic polynomial of $H_{m}, \gamma$ is the product of the subdiagonal elements of $H_{m}$ and $\eta_{m+1, m}$, and the function $\Delta_{w_{m}} f$ is defined as

$$
\begin{equation*}
\Delta_{w_{m}} f:=\frac{f-I_{w_{m}} f}{w_{m}} \tag{10}
\end{equation*}
$$

where $I_{w_{m}} f$ denotes the polynomial of degree $m-1$ which interpolates $f$ in the Hermite sense at the zeros of the polynomial $w_{m}$, i.e., at the Ritz values of $A$ with respect to $\mathscr{K}_{m}(A, \boldsymbol{b})$. Since the error (9) has exactly the same form as $f(A) \boldsymbol{b}$ with $\tilde{f}$ in place of $f$, one can proceed by computing corrections to $\hat{\boldsymbol{f}}_{1}=\boldsymbol{f}_{m}$ in the form of Arnoldi approximations to $\tilde{f}(A) \boldsymbol{v}_{m+1}$. In view of (9) and (10), this results in a corrected approximation $\hat{\boldsymbol{f}}_{2}=q(A) \boldsymbol{b}$, where now $q \in \mathscr{P}_{2 m-1}$ interpolates $f$ in the Ritz values of $A$ with respect to $\mathscr{K}_{m}(A, \boldsymbol{b})$ as well as those with respect to $\mathscr{K}_{m}\left(A, \boldsymbol{v}_{m+1}\right)$. We note that an alternative expression for the restarted approximations $\hat{\boldsymbol{f}_{m}}$ based on block Newton interpolation for Hermitian $A$ was given in [16].

In [6], the approach based on repeated block Newton interpolation was found to be unstable, and instead it was proposed to evaluate the matrix $f\left(\widehat{H}_{k}\right)$ in each cycle by standard algorithms such as Matlab's function funm, from which the entries required for the update (8) can be extracted. The resulting scheme is summarized below as Algorithm 1.

```
Algorithm 1. Restarted Arnoldi approximation for \(f(A) \boldsymbol{b}\) proposed in [6]
Given: \(A, \boldsymbol{b},\|\boldsymbol{b}\|=1, f\)
\(\boldsymbol{v}_{1}:=\boldsymbol{b}, \hat{f}_{0}:=\mathbf{0}\)
for \(k=1,2, \ldots\) until convergence do
        Compute Arnoldi decomposition \(A V_{k}=V_{k} H_{k}+\eta_{k+1} \boldsymbol{v}_{k m+1} \boldsymbol{e}_{m}^{T}\)
        of \(\mathscr{K}_{m}\left(A, \boldsymbol{v}_{(k-1) m+1}\right)\).
        if \(k=1\) then
        \(\widehat{H}_{k}:=H_{1}\)
    else
        \(\widehat{H}_{k}:=\left[\begin{array}{cc}\widehat{H}_{k-1} & O \\ \eta_{k} \boldsymbol{e}_{1} \boldsymbol{e}_{(k-1) m}^{T} & H_{k}\end{array}\right]\)
    Update the approximation \(\hat{\boldsymbol{f}}_{k}:=\hat{\boldsymbol{f}}_{k-1}+V_{k}\left[f\left(\widehat{H}_{k}\right) \boldsymbol{e}_{1}\right]_{(k-1) m+1: k m}\).
```

Although it allows discarding the basis vectors of previous cycles, Algorithm 1 has the shortcoming that it requires in the $k$ th cycle the evaluation of $f$ for a matrix of size $k m$. Despite the fact that, typically, $k m \ll n$, this can represent substantial computational effort as $k$ gets large. Moreover, it appears wasteful to compute $f\left(\widehat{H}_{k}\right)$ when only the last $m$ entries of its first column are needed.

An alternative approach for computing $f\left(\widehat{H}_{k}\right)$ which promises less work per cycle is to use a recursive scheme [20]: Comparing blocks in the identity

$$
\widehat{F}_{k} \widehat{H}_{k}=\widehat{H}_{k} \widehat{F}_{k}
$$

shows that, for $j>\ell$,

$$
F_{j, \ell} H_{\ell}-H_{j} F_{j, \ell}=E_{j} F_{j-1, \ell}-F_{j, \ell+1} E_{\ell+1}
$$

Since the diagonal blocks are obtained as $F_{k, k}=f\left(H_{k}\right)$, this relation allows us to compute the last block row of $\widehat{F}_{k}$ recursively by solving the Sylvester equations

$$
\begin{equation*}
X H_{k-j}-H_{k} X=E_{k} F_{k-1, k-j}-F_{k, k-j+1} E_{k-j+1}, \quad j=1,2, \ldots, k-1 \tag{11}
\end{equation*}
$$

for $X=F_{k, k-j}$. The Sylvester equation (11) is easy to solve since its coefficients $H_{k-j}$ and $H_{k}$ are upper Hessenberg (see [11]). We still, however, have to store $\widehat{H}_{k}$, i.e., $H_{1}, H_{2}, \ldots, H_{k}$ and $\eta_{2}, \eta_{3}, \ldots, \eta_{k}$. In addition, we need $F_{k-1, k-j}$, more precisely, $E_{k} F_{k-1, k-j}$, i.e., only the last row of $F_{k-1, k-j}(j=0,1, \ldots, k-1)$ has to be saved in the previous cycle. Note further that only the first column of $F_{k, k-j+1}$ enters the equation determining $F_{k, k-j}$, but we still compute $F_{k, j}(j=k, k-1, \ldots, 1)$, although only $F_{k, 1} e_{1}$ is needed, and, most importantly, the above Sylvester equation (11) tends to be severely ill-conditioned since $H_{k-j}$ and $H_{k}$ represent compressions of the same matrix $A$ and thus their spectra are by no means well separated.

### 3.2. Implementation based on a rational approximation of $f$

Our new implementation of restarting Krylov subspace algorithms for approximating (1) is based on evaluating $r\left(\widehat{H}_{k}\right) e_{1} \approx f\left(\widehat{H}_{k}\right) e_{1}$ using a rational approximation (cf. [9,23])

$$
f(\lambda) \approx r(\lambda)=p(\lambda)+\sum_{\ell=1}^{N} \frac{\alpha_{\ell}}{\omega_{\ell}-\lambda}
$$

of $f$ in partial fraction form with polynomial part $p$, coefficients $\alpha_{\ell}$ and poles $\omega_{\ell}$ which we assume to be simple and not contained in the field of values of $A$. In other words, we compute

$$
\begin{equation*}
r\left(\widehat{H}_{k}\right) \boldsymbol{e}_{1}=p\left(\widehat{H}_{k}\right) \boldsymbol{e}_{1}+\sum_{\ell=1}^{N} \alpha_{\ell}\left(\omega_{\ell} I-\widehat{H}_{k}\right)^{-1} \boldsymbol{e}_{1} \tag{12}
\end{equation*}
$$

We note that in the most common application $f(\lambda)=\exp (\lambda), \operatorname{Re} \lambda \leqslant 0$, a polynomial $p$ of degree zero or one usually suffices.

Evaluating $p\left(\widehat{H}_{k}\right) e_{1}$ for a polynomial $p$ of low degree is straightforward: Letting $p(\lambda)=$ $\pi_{1} \lambda+\pi_{0}$, for example, yields

$$
\hat{r_{0}}:=p\left(\widehat{H}_{k}\right) \boldsymbol{e}_{1}=\left[\begin{array}{c}
\left(\pi_{1} H_{1}+\pi_{0} I\right) e_{1} \\
\pi_{1} E_{2} e_{1} \\
\mathbf{0} \\
\vdots \\
\mathbf{0}
\end{array}\right]
$$

and for higher degrees one can proceed analogously. Evaluating the second expression in (12) consists of summing the vectors $\alpha_{\ell}\left(\omega_{\ell} I-\widehat{H}_{k}\right)^{-1} \boldsymbol{e}_{1}, \ell=1,2, \ldots, N$, i.e., solving the linear systems of equations $\left(\omega_{\ell} I-\widehat{H}_{k}\right) \hat{\boldsymbol{\gamma}}_{\ell}=\boldsymbol{e}_{1}$. Due to the sparsity pattern of the right hand side $\boldsymbol{e}_{1}$ and the block lower triangular form of $\widehat{H}_{k}$, this can be carried out recursively as

$$
\left(\omega_{\ell} I-H_{1}\right) \boldsymbol{r}_{\ell, 1}=\boldsymbol{e}_{1}, \quad\left(\omega_{\ell} I-H_{j}\right) \boldsymbol{r}_{\ell, j}=E_{j} \boldsymbol{r}_{\ell, j-1}, j=2, \ldots, k
$$

where we have partitioned $\hat{\boldsymbol{r}}_{\ell}=\left[\boldsymbol{r}_{\ell, 1}^{T}, \boldsymbol{r}_{\ell, 2}^{T}, \ldots, \boldsymbol{r}_{\ell, k}^{T}\right]^{T}$ conformingly with $\widehat{H}_{k}$. Note that these are $k$ Hessenberg systems of size $m$ and thus inexpensive to solve. Moreover, in view of (8), we only require the last block of $r\left(\widehat{H}_{k}\right) e_{1}$, which is obtained as

$$
[O, \ldots, O, I] r\left(\widehat{H}_{k}\right) \boldsymbol{e}_{1}=\boldsymbol{r}_{0, k}+\sum_{\ell=1}^{N} \alpha_{\ell} \boldsymbol{r}_{\ell, k}
$$

where $\boldsymbol{r}_{0, k}$ denotes the last block of $\hat{\boldsymbol{r}}_{0}$.
The resulting algorithm is summarized below as Algorithm 2. (For simplicity, we assume that the polynomial part $p$ of $r$ is the zero polynomial.)

```
Algorithm 2. Restarted Arnoldi approximation for \(f(A) \boldsymbol{b}\) based on rational approximation
Given: \(A, \boldsymbol{b},\|\boldsymbol{b}\|=1\), coefficients and poles \(\left(\alpha_{\ell}, \omega_{\ell}\right)_{\ell=1}^{N}\) of a rational function
    \(r \approx f\)
\(\boldsymbol{v}_{1}:=\boldsymbol{b}, \hat{f}_{0}:=\mathbf{0}\)
for \(k=1,2, \ldots\) until convergence do
    Compute Arnoldi decomposition \(A V_{k}=V_{k} H_{k}+\eta_{k+1} \boldsymbol{v}_{k m+1} \boldsymbol{e}_{m}^{T}\)
    of \(\mathscr{K}_{m}\left(A, \boldsymbol{v}_{(k-1) m+1}\right)\).
    if \(k=1\) then
        for \(\ell=1,2, \ldots, N\) do
            Solve \(\left(\omega_{\ell} I-H_{k}\right) \boldsymbol{r}_{\ell, 1}=\boldsymbol{e}_{1}\)
        else
            for \(\ell=1,2, \ldots, N\) do
            Solve \(\left(\omega_{\ell} I-H_{k}\right) \boldsymbol{r}_{\ell, k}=\eta_{k}\left(\boldsymbol{e}_{m}^{T} \boldsymbol{r}_{\ell, k-1}\right) \boldsymbol{e}_{1}\)
    \(\boldsymbol{h}_{k}:=\sum_{\ell=1}^{N} \alpha_{\ell} \boldsymbol{r}_{\ell, k}\)
    Update the approximation \(\hat{\boldsymbol{f}}_{k}:=\hat{\boldsymbol{f}}_{k-1}+V_{k} \boldsymbol{h}_{k}\).
```

In many applications both $A$ and $\boldsymbol{b}$ are real and $f$ has the property that $f(\bar{\lambda})=\overline{f(\lambda)}$. In this case $f(A) \boldsymbol{b}$ is also real and it is natural to approximate this vector using real arithmetic.

The rational approximation to $f$ is usually also real for real arguments, but its poles $\omega_{\ell}$ and coefficients $\alpha_{\ell}$ appear in complex conjugate pairs, say $\omega_{\ell+1}=\overline{\omega_{\ell}}$ and $\alpha_{\ell+1}=\overline{\alpha_{\ell}}$. Since all other quantities in the equations $\left(\omega_{\ell} I-\widehat{H}_{k}\right) \hat{\boldsymbol{r}}_{\ell}=\boldsymbol{e}_{1}$ are real, we have $\hat{\boldsymbol{r}}_{\ell+1}=\overline{\hat{\boldsymbol{r}}_{\ell}}$ and therefore $\boldsymbol{r}_{\ell+1, j}=\overline{\boldsymbol{r}_{\ell, j}}$ for all $j=1,2, \ldots, k$. For the quantities entering the update of $\hat{\boldsymbol{f}}_{k-1}$, we thus have

$$
\alpha_{\ell} \boldsymbol{r}_{\ell, k}+\alpha_{\ell+1} \boldsymbol{r}_{\ell+1, k}=\alpha_{\ell} \boldsymbol{r}_{\ell, k}+\overline{\alpha_{\ell} \boldsymbol{r}_{\ell, k}}=2 \operatorname{Re}\left(\alpha_{\ell} \boldsymbol{r}_{\ell, k}\right)
$$

and there is no need to solve $\left(\omega_{\ell+1} I-H_{k}\right) \boldsymbol{r}_{\ell+1, k}=E_{k} \boldsymbol{r}_{\ell+1, k-1}$. Setting $\boldsymbol{r}_{\ell, j}=\boldsymbol{r}_{\ell, j}^{(R)}+\mathrm{i} \boldsymbol{r}_{\ell, j}^{(I)}$ and $\omega_{\ell}=\omega_{\ell}^{(R)}+\mathrm{i} \omega_{\ell}^{(I)}$, a straightforward calculation shows that

$$
\begin{aligned}
& \left(\left|\omega_{\ell}\right|^{2} I-2 \omega_{\ell}^{(R)} H_{1}+H_{1}^{2}\right) \boldsymbol{r}_{\ell, 1}^{(R)}=\left(\omega_{\ell}^{(R)} I-H_{1}\right) \boldsymbol{e}_{1} \\
& \boldsymbol{r}_{\ell, 1}^{(I)}=\frac{1}{\omega_{\ell}^{(I)}}\left(\left[\omega_{\ell}^{(R)} I-H_{1}\right] \boldsymbol{r}_{\ell, 1}^{(R)}-\boldsymbol{e}_{1}\right)
\end{aligned}
$$

while for $j=2,3, \ldots, k$,

$$
\begin{aligned}
& \left(\left|\omega_{\ell}\right|^{2} I-2 \omega_{\ell}^{(R)} H_{j}+H_{j}^{2}\right) \boldsymbol{r}_{\ell, j}^{(R)}=\omega_{\ell}^{(I)} E_{j} \boldsymbol{r}_{\ell, j-1}^{(I)}+\left(\omega_{\ell}^{(R)} I-H_{j}\right) E_{j} \boldsymbol{r}_{\ell, j-1}^{(R)}, \\
& \boldsymbol{r}_{\ell, j}^{(I)}=\frac{1}{\omega_{\ell}^{(I)}}\left(\left[\omega_{\ell}^{(R)} I-H_{j}\right] \boldsymbol{r}_{\ell, j}^{(R)}-E_{j} \boldsymbol{r}_{\ell, j-1}^{(R)}\right) .
\end{aligned}
$$

Finally, $\alpha_{\ell} \boldsymbol{r}_{\ell, k}+\alpha_{\ell+1} \boldsymbol{r}_{\ell+1, k}=2 \operatorname{Re}\left(\alpha_{\ell} \boldsymbol{r}_{\ell, k}\right)=2\left[\operatorname{Re}\left(\alpha_{\ell}\right) \boldsymbol{r}_{\ell, k}^{(R)}-\operatorname{Im}\left(\alpha_{\ell}\right) \boldsymbol{r}_{\ell, k}^{(I)}\right]$ and we have avoided complex arithmetic.

To summarize, the two main ingredients of Algorithm 2 for computing $f(A) \boldsymbol{b}$ are the Arnoldi process (there is no difference to Algorithm 1) and (the partial fraction decomposition of) a rational function $r$ such that $r\left(\widehat{H}_{k}\right) \approx f\left(\widehat{H}_{k}\right)$. In the following we restrict our attention to the important special case that $f$ is the exponential function and always choose $r$ as its best uniform rational approximation on $(-\infty, 0]$ of type 16 as derived by Richard Varga and co-workers in [3] and [2]. Alternative rational approximations are described in [25].

Finally, we note that the rational function $r$ could be used right away to approximate

$$
f(A) \boldsymbol{b} \approx r(A) \boldsymbol{b}=p(A) \boldsymbol{b}+\sum_{\ell=1}^{N} \alpha_{\ell}\left(\omega_{\ell} I-A\right)^{-1} \boldsymbol{b}
$$

(see [18]). The evaluation of this approximation requires the solution of $N$ systems of linear equations

$$
\left(\omega_{\ell} I-A\right) \boldsymbol{x}_{\ell}=\boldsymbol{b}, \quad \ell=1,2, \ldots, N
$$

with shifted coefficient matrices (with complex shifts) and constant right-hand side. If $A$ is large and sparse these systems are usually solved by an iterative method, e.g., by a Krylov subspace method. It is well known (cf. [17]) that, due to the shift invariance of Krylov subspaces, it is sufficient to construct only one basis of $\mathscr{K}_{m}(A, \boldsymbol{b})$ to solve all these $N$ systems. In particular, if we use $\operatorname{FOM}(m)$ (restarted full orthogonalization method with restart length $m$ ) as a solver this would lead to the same approximations as Algorithm 2. If $A$ is Hermitian and positive (or negative) definite one could solve the $N$ linear systems simultaneously by the conjugate gradient method (see $[8,7]$ ) where each vector $\boldsymbol{x}_{\ell}$ can be computed by a short recurrence formula for which in total $2 N$ vectors of dimension $n$ need to be stored. In other words, the storage requirements do not grow with the dimension of the Krylov space and hence restarting is not required for reason of limited storage. Note that the cost of generating a Krylov basis of dimension km is essentially the same, both in the conjugate gradient method and in the restarted Hermitian case, where all the $H_{k}$ are tridiagonal. However, we have observed in many cases, that updating the $2 N$ long vectors of the conjugate gradient method (which is done after each matrix vector product, i.e., $k m$ times) is more expensive than the update step in Algorithm 2, $\hat{\boldsymbol{f}}_{k}:=\hat{\boldsymbol{f}}_{k-1}+V_{k} \boldsymbol{h}_{k}$, which is executed only once after $m$ matrix-vector products. Although convergence generally slows down when restarts are applied, i.e., more matrix-vector products are required, this extra cost is often compensated by fewer operations on long vectors. Therefore restart algorithms should be considered as an option even in the Hermitian case, if the delay in the convergence due to restarts is reasonable.

## 4. A stopping criterion based on an error indicator

To obtain a practical stopping criterion for our restarting schemes, we develop an error indicator using an extension of an idea described in [23, Theorem 5.1] (see also [21, Theorem 3.1]). The approach relies on an interpolation expansion of the approximation error $f(A) \boldsymbol{b}-\hat{\boldsymbol{f}}_{k}$ obtained by adding a sequence of auxiliary nodes $\left\{\theta_{j}\right\}$ in addition to the $k m$ Ritz values on which the restarted approximation $\hat{\boldsymbol{f}}_{k}$ is based.

Given $\tilde{m}$ complex numbers $\theta_{1}, \theta_{2}, \ldots, \theta_{\tilde{m}}$ such that $f\left(\theta_{j}\right)$ is defined for each $j$ and the associated nodal polynomials

$$
w_{0}(\lambda):=1, \quad w_{j}(\lambda):=\left(\lambda-\theta_{1}\right)\left(\lambda-\theta_{2}\right) \cdots\left(\lambda-\theta_{j}\right), j=1,2, \ldots, \tilde{m},
$$

we denote the associated divided differences of $f$ by

$$
\phi_{0}(\lambda):=f(\lambda), \quad \phi_{j}(\lambda):=\left[\Delta_{w_{j}} f\right](\lambda), j=1,2, \ldots, \tilde{m} .
$$

From the interpolation identity

$$
I_{w_{j+1}} f=I_{w_{j}} f+w_{j} \Delta_{w_{j}} f, \quad j=0,1, \ldots, \tilde{m}-1,
$$

we see that these obey the recursion

$$
\phi_{j}(\lambda)=\frac{\phi_{j-1}(\lambda)-\phi_{j-1}\left(\theta_{j}\right)}{\lambda-\theta_{j}}, \quad j=1,2, \ldots, \tilde{m} .
$$

Finally, for $1 \leqslant \ell \leqslant \tilde{m}$ and $0 \leqslant j \leqslant \tilde{m}-\ell$ we define $\Delta_{\ell}^{j} f$ to be the $j$ th order divided difference of $f$ with respect to the nodes $\theta_{\ell}, \theta_{\ell+1}, \ldots, \theta_{\ell+j}$ (see [27, Section 3.2]), i.e.,

$$
\Delta_{\ell}^{j} f:=\frac{1}{2 \pi \mathrm{i}} \int_{\Gamma} \frac{f(\lambda)}{\left(\lambda-\theta_{\ell}\right) \cdots\left(\lambda-\theta_{\ell+j}\right)} \mathrm{d} \lambda
$$

Given the Arnoldi-like decomposition (7), we now consider the matrix

$$
\begin{equation*}
W_{\tilde{m}}:=\left[w_{0}(A) \boldsymbol{v}_{k m+1}, w_{1}(A) \boldsymbol{v}_{k m+1}, \ldots, w_{\tilde{m}-1}(A) \boldsymbol{v}_{k m+1}\right] \in \mathbb{C}^{n \times \tilde{m}} \tag{13}
\end{equation*}
$$

and the bidiagonal matrix

$$
B_{\tilde{m}}:=\left[\begin{array}{cccc}
\theta_{1} & & & \\
1 & \theta_{2} & & \\
& \ddots & \ddots & \\
& & 1 & \theta_{\tilde{m}}
\end{array}\right] \in \mathbb{C}^{\tilde{m} \times \tilde{m}}
$$

in terms of which there holds $A W_{\tilde{m}}=W_{\tilde{m}} B_{\tilde{m}}+\left[\mathbf{0}, \ldots, \mathbf{0}, w_{\tilde{m}}(A) \boldsymbol{v}_{k m+1}\right]$. Together with (7) we obtain the Arnoldi-like decomposition

$$
A\left[\widehat{V}_{k} W_{\tilde{m}}\right]=\left[\widehat{V}_{k} W_{\tilde{m}}\right] \widetilde{H}_{k}+w_{\tilde{m}}(A) \boldsymbol{v}_{k m+1} \boldsymbol{e}_{k m+\tilde{m}}^{T}
$$

where

$$
\widetilde{H}_{k}=\left[\begin{array}{c|c}
\widehat{H}_{k} & O \\
\hline \widetilde{E}_{k+1} & B_{\tilde{m}}
\end{array}\right] \in \mathbb{C}^{(k m+\tilde{m}) \times(k m+\tilde{m})} \quad \text { and } \quad \widetilde{E}_{k+1}=\eta_{k+1} \boldsymbol{e}_{1} \boldsymbol{e}_{k m}^{T} \in \mathbb{R}^{\tilde{m} \times k m}
$$

If we now approximate $f(A) \boldsymbol{b}$ by $\tilde{\boldsymbol{f}}_{k}:=\left[\widehat{V}_{k} W_{\tilde{m}}\right] f\left(\widetilde{H}_{k}\right) \boldsymbol{e}_{1}$, then the associated error may be represented as (see [6, Theorem 2.6])

$$
\begin{equation*}
f(A) \boldsymbol{b}-\tilde{\boldsymbol{f}}_{k}=\tilde{f}(A) w_{\tilde{m}}(A) \boldsymbol{v}_{k m+1} \tag{14}
\end{equation*}
$$

where $\tilde{f}:=\gamma \Delta_{\tilde{w}} f$ with $\tilde{w} \in \mathscr{P}_{k m+\tilde{m}}$ the characteristic polynomial of $\tilde{H}_{k}$ and $\gamma$ the product of the subdiagonal entries of $\widetilde{H}_{k}$.

## Lemma 4.1. In terms of the notation introduced above, there holds

$$
f\left(\tilde{H}_{k}\right)=\left[\begin{array}{cc}
f\left(\widehat{H}_{k}\right) & O  \tag{15}\\
\widetilde{F}_{k, \tilde{m}} & f\left(B_{\tilde{m})}\right.
\end{array}\right],
$$

where

$$
f\left(B_{\tilde{m}}\right)=\left[\begin{array}{cccc}
f\left(\theta_{1}\right) & & &  \tag{16}\\
\Delta_{1}^{1} & f\left(\theta_{2}\right) & & \\
\vdots & \vdots & \ddots & \\
\Delta_{1}^{\tilde{\tilde{m}}-1} & \Delta_{2}^{\tilde{m}-2} & \cdots & f\left(\theta_{\tilde{m}}\right)
\end{array}\right] \in \mathbb{C}^{\tilde{m} \times \tilde{m}}
$$

and where $\widetilde{F}_{k, \tilde{m}} \in \mathbb{C}^{\tilde{m} \times k m}$ has the rows

$$
\begin{equation*}
\boldsymbol{e}_{j}^{T} \widetilde{F}_{k, \tilde{m}}=\eta_{k+1} \boldsymbol{e}_{k m}^{T} \phi_{j}\left(\widehat{H}_{k}\right), \quad j=1,2, \ldots, \tilde{m} \tag{17}
\end{equation*}
$$

Proof. The first assertion (16) was proven in [19]. To show (17) we follow the proof of Proposition 3.2 in [21]. Note first that comparing the $(2,1)$ blocks in the identity $f\left(\widetilde{H}_{k}\right) \widetilde{H}_{k}=\widetilde{H}_{k} f\left(\widetilde{H}_{k}\right)$ yields

$$
\begin{equation*}
\widetilde{F}_{k, \tilde{m}} \widehat{H}_{k}-B_{\tilde{m}} \widetilde{F}_{k, \tilde{m}}=\eta_{k+1}\left[\boldsymbol{e}_{1} \boldsymbol{e}_{k m}^{T} f\left(\widehat{H}_{k}\right)-f\left(B_{\tilde{m}}\right) \boldsymbol{e}_{1} \boldsymbol{e}_{k m}^{T}\right] \tag{18}
\end{equation*}
$$

We obtain (17) by induction on $j$. For $j=1$, multiplying by $\boldsymbol{e}_{1}^{T}$ on both sides of (18) yields

$$
\boldsymbol{e}_{1}^{T} \widetilde{F}_{k, \tilde{m}}\left(\widehat{H}_{k}-\theta_{1} I\right)=\eta_{k+1} \boldsymbol{e}_{k m}^{T}\left(f\left(\widehat{H}_{k}\right)-f\left(\theta_{1}\right) I\right),
$$

or

$$
\boldsymbol{e}_{1}^{T} \widetilde{F}_{k, \tilde{m}}=\eta_{k+1} \boldsymbol{e}_{k m}^{T}\left(f\left(\widehat{H}_{k}\right)-f\left(\theta_{1}\right) I\right)\left(\widehat{H}_{k}-\theta_{1} I\right)^{-1}=\eta_{k+1} \boldsymbol{e}_{k m}^{T} \phi_{1}\left(\widehat{H}_{k}\right)
$$

For $j>1$, multiplying (18) from the left by $\boldsymbol{e}_{j}^{T}$ leads to

$$
\boldsymbol{e}_{j}^{T} \widetilde{F}_{k, \tilde{m}}\left(\widehat{H}_{k}-\theta_{j} I\right)=\boldsymbol{e}_{j-1}^{T} \widetilde{F}_{k, \tilde{m}}-\eta_{k+1} \Delta_{1}^{j-1} \boldsymbol{e}_{k m}^{T}=\eta_{k+1} \boldsymbol{e}_{k m}^{T}\left(\phi_{j-1}\left(\widehat{H}_{k}\right)-\phi_{j-1}\left(\theta_{j}\right) I\right),
$$

from which (17) follows after multiplying by $\left(\widehat{H}_{k}-\theta_{j} I\right)^{-1}$. Note that we have tacitly assumed that $\theta_{j} \notin \Lambda\left(\widehat{H}_{k}\right)$ for all $j$. The usual confluent divided difference calculus shows, however, that the assertion is also valid without this restriction.

With the expression for the columns of $\widetilde{F}_{k, \tilde{m}}$ given in Lemma 4.1 together with the definition of $W_{\tilde{m}}$ (13), we find that

$$
\tilde{\boldsymbol{f}}_{k}=\widehat{V}_{k} f\left(\widehat{H}_{k}\right) \boldsymbol{e}_{1}+\eta_{k+1} \sum_{j=1}^{\tilde{m}}\left[\boldsymbol{e}_{k m}^{T} \phi_{j}\left(\widehat{H}_{k}\right) \boldsymbol{e}_{1}\right] w_{j-1}(A) \boldsymbol{v}_{k m+1}
$$

Together with the error representation (14) we have obtained

Theorem 4.2. The error of the restarted Arnoldi approximation (8) can be expanded as

$$
\begin{equation*}
f(A) \boldsymbol{b}-\widehat{V}_{k} f\left(\widehat{H}_{k}\right) \boldsymbol{e}_{1}=\eta_{k+1} \sum_{j=1}^{\tilde{m}}\left[\boldsymbol{e}_{k m}^{T} \phi_{j}\left(\widehat{H}_{k}\right) \boldsymbol{e}_{1}\right] w_{j-1}(A) \boldsymbol{v}_{k m+1}+\tilde{f}(A) w_{\tilde{m}}(A) \boldsymbol{v}_{k m+1}( \tag{19}
\end{equation*}
$$

Remark 4.3. By an obvious modification to a result given in [23, Theorem 5.1] one can show that the remainder term in (19) may also be written

$$
\tilde{f}(A) w_{\tilde{m}}(A) \boldsymbol{v}_{k m+1}=w_{\tilde{m}}(A)\left[\phi_{\tilde{m}}(A) \boldsymbol{b}-\widehat{V}_{k} \phi_{\tilde{m}}\left(\widehat{H}_{k}\right) \boldsymbol{e}_{1}\right] .
$$

The sum in (19) represents the leading $\tilde{m}$ terms of an interpolation series (see [27, Chapter III]) for the function $f-I_{w} f$, where $w \in \mathscr{P}_{k m}$ denotes the nodal polynomial of the eigenvalues of $\widehat{H}_{k}$. Provided the series converges as $\tilde{m} \rightarrow \infty-$ for $f=\exp$ it is sufficient to assume the nodes are bounded - we obtain the expansion

$$
f(A) \boldsymbol{b}-\widehat{V}_{k} f\left(\widehat{H}_{k}\right) \boldsymbol{e}_{1}=\eta_{k+1} \sum_{j=1}^{\infty}\left[\boldsymbol{e}_{k m}^{T} \phi_{j}\left(\widehat{H}_{k}\right) \boldsymbol{e}_{1}\right] w_{j-1}(A) \boldsymbol{v}_{k m+1}
$$

of the error of the restarted Arnoldi approximation. We obtain an error indicator by truncating the series after one or two terms, where we are free to choose the additional nodes $\left\{\theta_{j}\right\}_{j=1}^{\tilde{m}}$. In the numerical experiments presented below, we choose $\theta_{1}=\min \Lambda\left(\widehat{H}_{k}\right)$ for $\tilde{m}=1$ and $\theta_{2}=$ $\max \Lambda\left(\widehat{H}_{k}\right)$ with the same choice for $\theta_{1}$ when $\tilde{m}=2$. Computationally, all that is required for evaluating these error estimates are the coefficients $\boldsymbol{e}_{k m}^{T} \phi_{j}\left(\widehat{H}_{k}\right) \boldsymbol{e}_{1}$ and the vectors $w_{j-1}(A) \boldsymbol{v}_{k m+1}$. The former can be extracted from the first column of $\widetilde{F}_{k, \tilde{m}}$ (cf. (17)) after applying $f$ to $\widetilde{H}_{k}$, a matrix of size $k m+\tilde{m}$ only slightly larger than $\widehat{H}_{k}$, whose evaluation is called for in Algorithms 1 and 2. The latter require $\tilde{m}-1$ additional matrix-vector products of $A$ with $\boldsymbol{v}_{k m+1}$. We emphasize that for $\tilde{m}=1$ this means no additional matrix-vector products are necessary and that for $\tilde{m}=2$ the additional matrix-vector product $A \boldsymbol{v}_{k m+1}$ can be reused in the generation of the Krylov space of the $(k+1)$ st restart cycle.

Note that the choices for $\theta_{1}$ and $\theta_{2}$ given above for $\tilde{m}=1,2$ correspond to a Gauss-Radau and a Gauss-Lobatto quadrature formula for the error (cf. [10]). For certain functions, among these the exponential, it can be shown that these choices result in lower and upper bounds, respectively, for the error. A rigorous derivation of quadrature-based a-posteriori error bounds will be the subject of future work.

## 5. Convergence issues: a case study

With Algorithms 1 and 2, we have described two techniques for implementing a restarted Krylov subspace method for the evaluation of $f(A) \boldsymbol{b}$. These differ in the way the function $f$ is applied to the compressions $\widehat{H}_{k}$ of $A$. Considering specifically the exponential function $f=\exp$, in Algorithm 1 a standard routine such as the built-in Matlab function expm is applied. For the current implementation ${ }^{3}$ of expm, this means that we approximate $\exp \left(\widehat{H}_{k}\right)$ by a rational

[^3]

Fig. 1. Absolute errors of the approximations of $f(A) \boldsymbol{b}$ computed with Algorithms 1 and 2 for restart lengths $m=1,3$ and 10 .
expression of the form $r_{1}\left(\widehat{H}_{k}\right):=r\left(\widehat{H}_{k} / 2^{s}\right)^{2^{s}}$, where $r$ is a $[t / t]$ Padé fraction with $t \leqslant 13$ for the exponential function and $s \in \mathbb{N}_{0}$ depends on $\left\|\widehat{H}_{k}\right\|_{1}$ (see [12]). Note that $s$ and $r$, and thus $r_{1}$, depend on the argument. Therefore, $r_{1}$ is by no means a rational function - but rather it represents a family of rational functions from which one member is chosen depending on the current argument: If this is a scalar $\lambda$ then $r_{1}(\lambda)$ is an accurate approximation to $\exp (\lambda)$ regardless of where in the complex plane $\lambda$ is located. In floating point arithmetic, $r_{1}$ is for all practical purposes indistinguishable from exp; in particular, $r_{1}$ has no finite poles.

By contrast, Algorithm 2 approximates $\exp \left(\widehat{H}_{k}\right)$ by $r_{2}\left(\widehat{H}_{k}\right)$, where $r_{2}$ is a fixed rational function, namely the best uniform rational approximation of type 16 to the exponential function on $(-\infty, 0]$. In sharp contrast to $r_{1}$ above, $r_{2}$ approximates exp well only in a neighborhood of the negative real axis and therefore lacks the universal approximation property of $r_{1}$. Moreover, $r_{2}$ has finite poles. When approximating $\exp (A)$ for matrices with real nonpositive eigenvalues, as arise e.g. in connection with parabolic initial value problems, this may appear as a somewhat academic issue; but, as we shall see, this leads to a tremendous difference in numerical behavior between the two approaches if the restart length $m$ is sufficiently small. The main distinction is that a restarted Krylov approximation to $f(A) \boldsymbol{b}$ converges superlinearly if $f$ is an entire function (such as exp, see [6, Theorem 4.2]), whereas only a linear asymptotic rate of convergence results if $f$ has finite singularities (such as $r_{2}$ ).

We demonstrate this using the following model problem. Let

$$
A=\operatorname{diag}(-100,-99, \ldots, 0) \in \mathbb{R}^{101 \times 101} \quad \text { and } \quad \boldsymbol{b}=[1,1, \ldots, 1]^{T} / \sqrt{101} \in \mathbb{R}^{101}
$$

(We choose such a simple example in order that the quantities entering our analysis be explicitly known.) Applying Algorithms 1 and 2 yields approximations to $\exp (A) \boldsymbol{b}$ which we shall denote by $\hat{\boldsymbol{f}}_{k}^{(1)}$ and $\hat{\boldsymbol{f}}_{k}^{(2)}$, respectively. We consider restart cycles consisting of $m=1, m=3$ and $m=10$ steps. Fig. 1 shows the corresponding absolute errors $\left\|\hat{\boldsymbol{f}}_{k}^{(\mu)}-\exp (A) \boldsymbol{b}\right\|$ in the Euclidian norm ( $\mu=1,2$ ).

For restart length $m=1$, we observe that Algorithm 1 converges while Algorithm 2 does not. To explain why we recall that both algorithms are based on interpolation processes. The interpolation nodes are the eigenvalues of the Hessenberg matrices $\widehat{H}_{k}$ which, for restart length
$m=1$, are bidiagonal. It can be shown ${ }^{4}$ that, in our example, the diagonal entries (and thus the eigenvalues) of $\widehat{H}_{k}$ are all equal to $\vartheta=-50$. Thus for restart length $m=1$, both algorithms are based on (truncated) Taylor expansions of $f=\exp$ and $f=r_{2}$, respectively, about $\vartheta=-50$ : We have

$$
\begin{aligned}
& \hat{\boldsymbol{f}}_{k}^{(1)}=\widehat{V}_{k} \exp \left(\widehat{H}_{k}\right) \boldsymbol{e}_{1}=q_{k-1}^{(1)}(A) \boldsymbol{b}=\sum_{j=0}^{k-1} \frac{\exp (-50)}{j!}(A+50 I)^{j} \boldsymbol{b} \quad \text { and } \\
& \hat{\boldsymbol{f}}_{k}^{(2)}=\widehat{V}_{k} r_{2}\left(\widehat{H}_{k}\right) \boldsymbol{e}_{1}=q_{k-1}^{(2)}(A) \boldsymbol{b}=\sum_{j=0}^{k-1} \frac{r_{2}^{(j)}(-50)}{j!}(A+50 I)^{j} \boldsymbol{b} .
\end{aligned}
$$

The Taylor polynomials $\left\{q_{k-1}^{(1)}\right\}_{k \geqslant 1}$ converge (albeit slowly) to exp uniformly on compact subsets of $\mathbb{C}$. Fig. 2 shows the errors $\left\|q_{k-1}^{(1)}-\exp \right\|_{\infty, A(A)}:=\max _{\lambda \in \Lambda(A)}\left|q_{k-1}^{(1)}(\lambda)-\exp (\lambda)\right|$ of these Taylor polynomials, and these are seen to agree perfectly with the errors of Algorithm 1. The Taylor polynomials $\left\{q_{k-1}^{(2)}\right\}_{k \geqslant 1}$, on the other hand, converge to $r_{2}$ in a disk with center $\vartheta=-50$ and radius $\min _{\omega}|\omega+50|$, where $\omega$ runs over all poles of $r_{2}$, and they diverge outside this disk. The poles of $r_{2}$ closest to $\vartheta=-50$ are $\omega \approx-11 \pm 19 \mathrm{i}$ with $|\omega-\vartheta| \approx 44<50$. In other words, the Taylor series of $r_{2}$ (with expansion point $\vartheta=-50$ ) has $\tau \approx 44$ as its radius of convergence. Since some of the eigenvalues of $A$ lie outside the convergence disk, the sequence $\left\{\hat{\boldsymbol{f}}_{k}^{(2)}\right\}_{k \geqslant 1}$ must ultimately diverge like $\left[\max _{\lambda \in \Lambda(A)}|\lambda+50| / 44\right]^{k}=(50 / 44)^{k} \approx 1.14^{k}$ (cf. the dotted line in Fig. 2). Moreover, Fig. 2 shows that also $\left\|q_{k-1}^{(2)}-\exp \right\|_{\infty, \Lambda(A)}$ is in perfect agreement with the errors of Algorithm 2. (Note that this error curve cannot be distinguished from $\left\|q_{k-1}^{(2)}-r_{2}\right\|_{\infty, A(A)}$ because $\left\|r_{2}-\exp \right\|_{\infty, A(A)}<10^{-15}$.)

The matrices $\widehat{H}_{k}$ are highly nonnormal $(\vartheta=-50$ is their only eigenvalue, with algebraic multiplicity $k$ but geometric multiplicity 1 ). This raises the question of whether it is justified to base our analysis on this eigenvalue of $\widehat{H}_{k}$, which is extremely sensitive to perturbations, or whether an approach using pseudo-eigenvalues would not be more appropriate. Our answer is that it does not matter. This is a consequence of a theorem due to Walsh [27, Theorem 7.1] on the overconvergence of differences of interpolating polynomials (which by the way, is another area where Richard Varga has made significant contributions). As mentioned above, the diagonal entries of $\widehat{H}_{k}$ are all equal to -50 . We observe that its subdiagonal entries $\eta_{j, j-1}$ converge to 50 (see [1] for a proof), i.e., asymptotically $\widehat{H}_{k}$ resembles a Toeplitz matrix with symbol 50( $\lambda-1$ ), and thus its pseudo-spectra (asymptotically) are disks with center -50 (cf. [22]). In Fig. 2, the eigenvalues of $\widehat{H}_{100}+E$ are shown for 10 random perturbations of norm $\|E\|=\varepsilon=2^{-52}$. The large majority of these pseudo-eigenvalues lies on a circle with center -50 and some radius $\delta>0$ (here $\delta \approx 32$ ). If we now consider $\hat{\boldsymbol{f}}_{k}^{(3)}=q_{k-1}^{(3)}(A) \boldsymbol{b}$, where $q_{k-1}$ interpolates $r_{2}$ at the nodes $-50+$

[^4]

Fig. 2. Left: $\operatorname{errors}^{\max } \boldsymbol{\lambda}_{\lambda \in \Lambda(A)}\left|q_{k-1}^{(\mu)}(\lambda)-\exp (\lambda)\right|$ for the truncated Taylor series $q_{k-1}^{(\mu)}$ of $\exp (\mu=1)$ and $r_{2}(\mu=2)$. Right: eigenvalues of $\widehat{H}_{100}+E$ for 10 random matrices of norm $\|E\|=2^{-52}$.
$\delta \exp (2 \pi i j / k)(j=0,1, \ldots, k-1)$, then Walsh's overconvergence result tells us that, as long as $\delta<44$, the difference of the interpolating polynomials $\left|q_{k-1}^{(2)}(\lambda)-q_{k-1}^{(3)}(\lambda)\right|$ tends to zero for $\mid \lambda+$ $50 \mid \leqslant 44^{2} / \delta$. This convergence is linear, more precisely $\lim \sup _{k \rightarrow \infty} \max _{|\lambda+50| \leqslant \tau} \mid q_{k-1}^{(2)}(\lambda)-$ $\left.q_{k-1}^{(3)}(\lambda)\right|^{1 / k}=(\delta \tau) / 44^{2}$. We set $\delta=32$ and observe that $\lim \sup _{k \rightarrow \infty}\left\|q_{k-1}^{(2)}-q_{k-1}^{(3)}\right\|_{\infty, \Lambda(A)}^{1 / k} \leqslant$ $(50 \cdot 32) /\left(44^{2}\right) \approx 0.83$. In other words, $\left\|\hat{\boldsymbol{f}}_{k}^{(2)}-\hat{\boldsymbol{f}}_{k}^{(3)}\right\|$ tends to zero like $0.83^{k}$ which is negligible compared to the size of $\left\|\hat{\boldsymbol{f}}_{k}^{(2)}-\exp (A) \boldsymbol{b}\right\|$.

The behavior of Algorithms 1 and 2 for the restart lengths $m=3$ and $m=10$ (cf. Fig. 1) can be explained along the same lines: Algorithm 1 relies on an interpolation process for the exponential function whereas in Algorithm 2 its uniform best rational approximation $r_{2}$ is interpolated. The nodes are again the eigenvalues $\vartheta_{1}, \vartheta_{2}, \ldots$ of the block bidiagonal matrices $\widehat{H}_{k}$, i.e., the eigenvalues of its diagonal blocks $H_{1}, H_{2}, \ldots, H_{k}$ which are symmetric tridiagonal matrices of dimension $m=3$ and $m=10$, respectively. Now the following observation is crucial: There holds

$$
\lim _{j \rightarrow \infty} H_{2 j-1}=: \widetilde{H}_{1} \quad \text { and } \quad \lim _{j \rightarrow \infty} H_{2 j}=: \widetilde{H}_{2}
$$

(see [1] for a theoretical justification). The nodal sequence $\vartheta_{1}, \vartheta_{2}, \vartheta_{3}, \ldots$ therefore has the property

$$
\lim _{j \rightarrow \infty} \vartheta_{2 m j+v}=: \tilde{\vartheta}_{v} \quad \text { for } v=1,2, \ldots, 2 m
$$

Asymptotically, we interpolate exp and $r_{2}$ at $2 m$ nodes which are repeated cyclically. In our example, these nodes are given by

$$
\tilde{\vartheta}_{1} \approx-87, \quad \tilde{\vartheta}_{2}=-50, \quad \tilde{\vartheta}_{3} \approx-13, \quad \tilde{\vartheta}_{4} \approx-92, \quad \tilde{\vartheta}_{5}=-50, \quad \tilde{\vartheta}_{6} \approx-8
$$

if $m=3$. The convergence properties of the corresponding interpolation polynomials can be described in terms of the lemniscates


Fig. 3. Lemniscates $L_{\tau}$ for $\tau=\tau_{s}$ and $\tau=\tau_{A}$ governing the convergence rate of Algorithm 2 for restart lengths $m=3$ (left) and $m=10$ (right). The squares mark the nodes $\vartheta_{j}$, the dots the poles of the rational approximation $r_{2}$.

$$
L_{\tau}:=\left\{\lambda \in \mathbb{C}:\left|w_{2 m}(\lambda)\right|=\tau^{2 m}\right\}, \quad \tau>0, \quad w_{2 m}(\lambda)=\prod_{\nu=1}^{2 m}\left(\lambda-\tilde{\vartheta}_{\nu}\right)
$$

From a theorem of Walsh [27, Theorem 3.6] it follows that

$$
\limsup _{k \rightarrow \infty}\left\|\hat{\boldsymbol{f}}_{k}^{(2)}-\exp (A) \boldsymbol{b}\right\|^{1 / k m}=\frac{\tau_{A}}{\tau_{s}}
$$

where
$\tau_{s}=\max \left\{\tau: r_{2}\right.$ is analytic in the interior of $\left.L_{\tau}\right\}$,
i.e., $L_{\tau_{s}}$ is the largest lemniscate with foci $\tilde{\vartheta}_{v}(v=1,2, \ldots, 2 m)$ such that $r_{2}$ is analytic in the interior of $L_{\tau_{s}}$, and where
$\tau_{A}=\min \left\{\tau: \Lambda(A)\right.$ is contained in the closed interior of $\left.L_{\tau}\right\}$,
i.e., $L_{\tau_{A}}$ is the smallest lemniscate with foci $\tilde{\vartheta}_{v}$ such that all eigenvalues of $A$ are contained in the closed interior of $L_{\tau_{A}}$. Fig. 3 shows these lemniscates for our example.

To summarize, if we replace the exponential function by its rational best approximation then the convergence of the restarted Lanczos method shows two phases. Initially, we observe the error behavior of a polynomial approximation to an entire function, i.e., after a start-up phase, where the error is not reduced, the polynomial converges superlinearly. But there is a point from where on the poles become visible and then we have slower linear convergence or even linear divergence. This point is fairly independent of the restart length, whereas the linear rate of convergence/divergence depends on it (of course, it also depends on the eigenvalues of $A$ ). The aim is to choose the restart length large enough such that at the point of transition the desired accuracy is reached or nearly reached.

## 6. Numerical examples

In this section we illustrate the performance of the two restart algorithms for some initialboundary value problems. All computations were carried out in Matlab Release 2007a on an Intel Xeon 5160 at 3 GHz with 16 GB RAM running SuSE Linux Enterprise Server (SLES) Version 10.

Table 1 Execution times, number of matrix-vector products and final accuracy for the solution of the heat equation ( $n=50$ ) for different restart lengths $m$

| $m$ | Algorithm 1 |  |  | Algorithm 2 |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | Time (s) | Mvp | Acc. | Time (s) | Mvp | Acc. |
| $\infty(I)$ | 5.6 | 282 | 5e-14 | 3.9 | 282 | 5e-12 |
| $\infty(I I)$ | 9.7 | 564 | $5 \mathrm{e}-14$ | 7.6 | 564 | 5e-12 |
| 50 | 4.8 | 350 | $3 \mathrm{e}-14$ | 4.1 | 300 | 6e-12 |
| 30 | 7.9 | 360 | 2e-14 | 6.0 | 330 | $5 \mathrm{e}-12$ |
| 20 | 9.5 | 380 | 5e-15 | 7.2 | 400 | 6e-12 |
| 10 | 17.4 | 430 | $9 \mathrm{e}-15$ | - | - | - |

Here, $m=\infty(I)$ refers to the standard (unrestarted) Lanczos method, while $m=\infty(I I)$ stands for the (unrestarted) two-pass Lanczos algorithm.

### 6.1. The heat equation

Our first numerical experiment is based on a standard example in this area (see, e.g. [9,6]): We consider the initial-boundary value problem

$$
\begin{align*}
& \partial_{t} u-\Delta u=0 \quad \text { in } \Omega=(0,1)^{3}, \quad t>0,  \tag{20a}\\
& u(x, t)=0 \quad \text { on } \Gamma=\partial \Omega, \quad t>0  \tag{20b}\\
& u(x, 0)=u_{0}(x) \quad \text { in } \Omega . \tag{20c}
\end{align*}
$$

When the Laplacian is discretized by the usual seven-point stencil on a uniform grid involving $n_{1}$ interior grid points in each Cartesian direction, problem (20) reduces to the initial value problem

$$
\begin{aligned}
& \boldsymbol{u}^{\prime}(t)=A \boldsymbol{u}(t), \quad t>0 \\
& \boldsymbol{u}(0)=\boldsymbol{u}_{0}
\end{aligned}
$$

with an $n \times n$ matrix $A\left(n=n_{1}^{3}\right)$ and an initial vector $\boldsymbol{u}_{0}$ consisting of the values $u_{0}(x)$ at the grid points $x$, the solution of which is given by

$$
\begin{equation*}
\boldsymbol{u}(t)=\exp (t A) \boldsymbol{u}_{0} \tag{21}
\end{equation*}
$$

As in [9], we give the initial vector in terms of its expansion in eigenfunctions of the discrete Laplacian as

$$
\boldsymbol{u}_{0}^{i, j, k}=\sum_{i^{\prime}, j^{\prime}, k^{\prime}} \frac{1}{i^{\prime}+j^{\prime}+k^{\prime}} \sin \left(i i^{\prime} \pi h\right) \sin \left(j j^{\prime} \pi h\right) \sin \left(k k^{\prime} \pi h\right) .
$$

Here $h=1 /\left(n_{1}+1\right)$ is the mesh size and the triple indexing is relative to the lexicographic ordering of the mesh points in the unit cube.

We consider this problem for a discretization with $n_{1}=50$, resulting in a matrix of dimension $n=125,000$. We apply the restarted Lanczos approximation with restart lengths $m=10,20,30$, and 50 as well as the full Lanczos algorithm ( $m=\infty(I)$ ) and the two-pass Lanczos method $(m=\infty(I I))$. Each iteration is run until the accuracy no longer improves. Note that, in the full methods $(m=\infty)$, the evaluation of $\exp \left(H_{k}\right)$ is only performed once, when final accuracy is reached. The resulting execution times are shown in Table 1.

For $m=10$, Algorithm 2 converges so slowly that the final accuracy was not reached even after 2000 matrix-vector multiplications. The fastest method is clearly the full Lanczos algorithm


Fig. 4. The error indicators described in Section 4 applied to the heat equation for the restart lengths $m=20$ (left) and $m=50$ (right) and for Algorithm 1 (dashed line) and 2 (dotted line).
( $m=\infty(I)$ ), but note that it requires storing 282 vectors of dimension $n=125,000$ ( 282 MB ). To our surprise, the restarted method (for $m=30$ and $m=50$ ) is faster than the two-pass Lanczos algorithm. One also observes that Algorithm 2 is generally faster than Algorithm 1, but that the final accuracy of Algorithm 1 is higher. This loss of accuracy for Algorithm 2 stems from the eight linear systems we have to solve in each cycle. For $m=50$ the condition numbers of their coefficient matrices vary between $2 \times 10^{2}$ and $4 \times 10^{2}$ (see Table 2).

We conclude this example by applying the error indicators described in Section 4 for the restart lengths $m=20$ and $m=50$ (see Fig. 4). We obtain lower ( $\tilde{m}=1$ ) and upper ( $\tilde{m}=2$ ) bounds for our two restart algorithms. By their construction it is clear that the error indicators cannot detect when the method begins to stagnate because final accuracy has been reached (but this can easily be detected by other means). However, for Algorithm 2 the point where the convergence behavior changes from superlinear to linear is located precisely (even in cases where this occurs long after final accuracy is reached).

### 6.2. Maxwell's equations

We next consider a problem which occurs in geoelectrical exploration and for which Lanczosbased Krylov subspace approximations have been very successful, see [5] and the references given there. In the absence of impressed source currents the time-evolution of an electric field $\boldsymbol{E}=\boldsymbol{E}(x, t)$ from a given initial state $\boldsymbol{E}_{0}=\boldsymbol{E}_{0}(x)$ at time $t_{0}$ is the solution of the initial value problem

$$
\begin{align*}
& \partial_{t}(\sigma \boldsymbol{E})+\nabla \times\left(\mu^{-1} \nabla \times \boldsymbol{E}\right)=\mathbf{0} \quad \text { in } \Omega, \quad t>t_{0},  \tag{22a}\\
& \boldsymbol{n} \times \boldsymbol{E}=\mathbf{0} \quad \text { on } \partial \Omega,  \tag{22b}\\
& \boldsymbol{E}\left(x, t_{0}\right)=\boldsymbol{E}_{0}(x) \quad \text { in } \Omega \tag{22c}
\end{align*}
$$

for the quasi-static Maxwell's equations. Here the domain is a cube $\Omega=(-L, L)^{3}$, where $L=$ 2000 and the magnetic permeability $\mu$ and electric conductivity $\sigma$ are taken to have the constant values $\mu=4 \pi \times 10^{-7}$ and $\sigma=0.1$. The initial data $\boldsymbol{E}_{0}$ is the field at time $t_{0}=10^{-6}$ of a vertical magnetic dipole of unit strength located at the origin for which an analytic expression is known (see [28]).


Fig. 5. Computational domain and tensor product grid lines shown on the boundary for problem (22). The axis indicates the orientation of the magnetic dipole and the circular lines indicate the direction of the electric field.


Fig. 6. Absolute errors of the approximations of $f(A) \boldsymbol{b}$ for different restart lengths $m$ in the case of Maxwell's equations.
We discretize the operator $\sigma^{-1} \nabla \times\left(\mu^{-1} \nabla \times \cdot\right)$ in space using the Yee finite difference scheme [29] on a graded tensor product mesh (see Fig. 5). After symmetrization (cf. [5]) this yields a symmetric matrix $A \in \mathbb{R}^{n \times n}$ and a vector $e_{0} \in \mathbb{R}^{n}$ which is a sampled version of the initial electric field. The semi-discretized system (22) thus reduces to a linear linear system of ordinary differential equations with constant coefficients

$$
\boldsymbol{e}^{\prime}(t)=-A \boldsymbol{e}(t), \quad \boldsymbol{e}\left(t_{0}\right)=\boldsymbol{e}_{0}
$$

with solution

$$
\boldsymbol{e}(t)=\exp \left(-\left(t-t_{0}\right) A\right) \boldsymbol{e}_{0}
$$

The eigenvalues of $A$ are contained in $\left[0, \lambda_{\max }\right], \lambda_{\max } \leqslant 13 /\left(h_{\min }^{2} \sigma \mu\right)$, where $h_{\min }$ is the minimal distance of two adjacent grid-points.

In our example (cf. Table 2 and Fig. 6) the matrix $A$ is of size $n=565,326$ and $\lambda_{\max } \approx 10^{8}$. We approximate $\boldsymbol{E}\left(t_{1}\right), t_{1}=10^{-3}$, using the Lanczos algorithm (I), the two-pass Lanczos algorithm

Table 2 Performance comparison of several variants of Lanczos and restarted Lanczos methods for Maxwell's equations

| $m$ | Algorithm 1 |  |  | Algorithm 2 |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | Time (s) | Mvp | Acc. | Time (s) | Mvp | Acc. |
| $\infty(I)$ | 118 | 1072 | $9.93 \mathrm{e}-13$ | 86 | 1072 | $9.93 \mathrm{e}-13$ |
| $\infty(I I)$ | 176 | 2144 | $9.93 \mathrm{e}-13$ | 144 | 2144 | $9.93 \mathrm{e}-13$ |
| 90 | 273 | 1350 | $1.92 \mathrm{e}-13$ | 118 | 1350 | $2.01 \mathrm{e}-13$ |
| 70 | 339 | 1400 | $3.28 \mathrm{e}-13$ | 112 | 1400 | $9.13 \mathrm{e}-13$ |
| 50 | 613 | 1600 | $2.10 \mathrm{e}-13$ | Slow convergence |  |  |
| 30 | 2014 | 2040 | 5.64e-13 | Divergen |  |  |

(II), as well as Algorithm 1 and Algorithm 2 given above. As a reference solution $\boldsymbol{e}\left(t_{1}\right)$ we used the approximation obtained by the Lanczos method when it stagnates at final accuracy.

In Algorithm 2 we solve eight tridiagonal linear systems in each cycle instead of calling Matlab's expm which is much more expensive. As a consequence the execution time of Algorithm 2 is dominated by the number of matrix-vector products. This explains why the restarted Krylov methods perform slightly faster than the full two-pass Lanczos method if Algorithm 2 is applied. We terminated our timing measurements for all algorithms when the absolute error to $\boldsymbol{e}\left(t_{1}\right)$ fell below $10^{-12}$.

### 6.3. The advection-diffusion equation

We consider the initial value problem

$$
\begin{align*}
& \partial_{t} u-\frac{1}{\mathrm{Pe}} \Delta u+\boldsymbol{a} \cdot \nabla u=0 \quad \text { in } \Omega=(-1,1) \times(0,1),  \tag{23a}\\
& u=1-\tanh (\mathrm{Pe}) \quad \text { on } \Gamma_{0},  \tag{23b}\\
& u=1+\tanh ((2 x+1) \mathrm{Pe}) \quad \text { on } \Gamma_{\mathrm{in}},  \tag{23c}\\
& \frac{\partial u}{\partial n}=0 \quad \text { on } \Gamma_{\mathrm{out}},  \tag{23d}\\
& u(x, 0)=u_{0}(x) \quad \text { in } \Omega \tag{23e}
\end{align*}
$$

for the advection-diffusion equation, which is a popular benchmark problem for discretizations of advection-dominated problems, see [24]. The convective field is given as

$$
\boldsymbol{a}(x, y)=\left[\begin{array}{c}
2 y\left(1-x^{2}\right) \\
-2 x\left(1-y^{2}\right)
\end{array}\right], \quad(x, y) \in \Omega
$$

and the boundary $\Gamma=\partial \Omega$ is divided into the inflow boundary $\Gamma_{\mathrm{in}}:=[-1,0] \times\{0\}$, the outflow boundary $\Gamma_{\text {out }}:=[0,1] \times\{0\}$ and the remaining portion $\Gamma_{0}$ (cf. Fig. 7, left). The Péclet number Pe is a nondimensional parameter describing the strength of advection relative to diffusion and therefore also how far the discrete operators are from symmetric.

We discretize the advection-diffusion operator for $\mathrm{Pe}=10$ in space using linear finite elements on a triangulation generated by the adaptive mesh generation facility in the Comsol Multiphysics finite element software (version 3.3a). In the resulting mesh, shown on the right of Fig. 7, one can observe refinements near the sharp transition in the inflow profile and near the origin, around which the advection field rotates. After multiplying from the left by the square root of the (lumped) mass matrix, ${ }^{5}$ the semi-discretized problem is a system of ODEs

[^5]

Fig. 7. Geometry and streamlines for the advection-diffusion problem (23).


Fig. 8. Relative errors of the full and restarted Arnoldi approximation for the advection-diffusion problem.

$$
\boldsymbol{u}^{\prime}(t)=A \boldsymbol{u}(t)+\boldsymbol{g}, \quad \boldsymbol{u}(0)=\boldsymbol{u}_{0}
$$

where $A$ is of size $n=2157$ and the constant inhomogeneous term $g$ results from the inhomogeneous Dirichlet boundary condition. We then approximate the matrix exponential part of the solution

$$
\boldsymbol{u}(t)=\exp (t A)\left(\boldsymbol{u}_{0}+A^{-1} \boldsymbol{g}\right)-A^{-1} \boldsymbol{g}
$$

at time $t=6$, at which the flow hat reached a steady state, starting from rest $\boldsymbol{u}_{0}=\mathbf{0}$, using the unrestarted Arnoldi approximation as well as the restarted schemes of Algorithms 1 and 2.

The error curves for restart lengths $m=30,60$ and 90 are shown in Fig. 8. We observe that, despite the rather small system size, even the unrestarted Arnoldi method requires roughly 500 steps to reach an accuracy of around $10^{-12}$. This slow convergence as well as the nonmonotonic, somewhat more erratic convergence curve indicate that this is a harder problem than the preceding two. We also observe that the methods of restart lengths $m=60$ and $m=90$ converge at nearly the same rate, and it is noticeable that Algorithm 2 reaches a final accuracy of only around $10^{-12}$ in all cases. The main reason for this loss of accuracy lies in the fact that the spectrum of $A$ extends into the complex plane, and on the spectrum we have $\left\|r_{2}-\exp \right\|_{\infty, \Lambda(A)} \approx 10^{-14}$ for the best uniform rational approximation $r_{2}$ on $(-\infty, 0]$. Moreover, since $A$ is highly nonnormal, the difference $r_{2}(A)-\exp (A)$ is not determined by the eigenvalues of $A$ alone and on the field of values $W(A)$ of $A$ we have $\left\|r_{2}-\exp \right\|_{\infty, W(A)} \approx 10^{-6}$.

Table 3 gives execution times for this example. We observe that the efficiency of Algorithm 2 over Algorithm 1 is most pronounced here. The fact that both restarted variants took longer

Table 3 Execution times, number of matrix-vector products and final accuracy for the solution of the advection-diffusion problem for different restart lengths $m$

| $m$ | Algorithm 1 |  |  | Algorithm 2 |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | Time (s) | Mvp | Acc. | Time (s) | Mvp | Acc. |
| $\infty$ | 8.6 | 410 | 9e-13 | 6.7 | 420 | 9e-13 |
| 90 | 90.1 | 1170 | $8 \mathrm{e}-14$ | 12.0 | 1170 | $3 \mathrm{e}-13$ |
| 60 | 121.0 | 1140 | 6e-13 | 11.8 | 1140 | $8 \mathrm{e}-13$ |
| 30 | 685.7 | 1560 | 4e-13 | 15.3 | 1560 | $6 \mathrm{e}-13$ |

than the full versions is attributed to the small size of the problem. For large problems, where the orthogonalization effort of full Arnoldi becomes more noticeable, we expect the timings to increasingly favor the restarted versions. Note that the advantage of requiring less storage is present also for these small problems.

## 7. Conclusions

Restarting Krylov subspace approximations of $f(A) \boldsymbol{b}$ is of interest because short recurrences for Krylov basis vectors do not translate to short recurrences for the quantity being approximated. We have introduced an efficient implementation for a restarted Krylov subspace approximation and compared it against the approach introduced in [6], both in terms of execution time and convergence properties. The new approach, Algorithm 2, is faster, as it solves a fixed number of linear systems of equations the size of the restart length $m$ instead of, as is the case with Algorithm 1, evaluating a function of a matrix of increasing size $k m$ in the $k$ th cycle. Algorithm 2, while sometimes considerably faster, has the disadvantage that the restart length may need to be chosen somewhat larger to ensure convergence. Moreover, the solution of the linear equations in each cycle of Algorithm 2 can introduce some ill-conditioning which limits the final attainable accuracy. For the examples considered here involving the exponential propagation of semi-discretized partial differential operators subject to discretization errors, the requirements on final accuracy are usually sufficiently low that this is not a severe limitation. We have further introduced an error indicator which allows the termination of the iteration once sufficient or final accuracy has been reached. Moreover, we have pointed out the fundamental limitations of using a fixed rational approximation of $f$ to evaluate $f\left(\widehat{H}_{k}\right)$ when the poles of the former approach the spectrum of $A$.

For the case of the exponential function, the new method was seen to be competitive with established methods for two symmetric problems, and a viable solution approach for a difficult non-Hermitian problem. In summary, Algorithm 2 is an attractive scheme, particularly when memory is limited so that the unrestarted Arnoldi method is not an option.

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[^1]:    ${ }^{1}$ When we use terminology such as Arnoldi decomposition, Arnoldi approximation, and Arnoldi algorithm in the context of a Hermitian matrix $A$, we tacitly assume that computations are carried out with the Hermitian Lanczos process.

[^2]:    ${ }^{2}$ From now on $e_{1}$ denotes a first unit coordinate vector whose dimension is dictated by the context, whereas, for $j \geqslant 2, e_{j} \in \mathbb{R}^{j}$ shall denote the $j$ th unit coordinate vector.

[^3]:    ${ }^{3}$ Release 2007a.

[^4]:    ${ }^{4}$ Let $p_{k}(\lambda)=\lambda^{k} \in \mathscr{P}_{k}$ denote the $k$ th monomial and define the vectors $\boldsymbol{w}_{k}=\left[p_{k}(-50), p_{k}(-49), \ldots, p_{k}(50)\right]^{T}$. Then $\quad \sum_{j=-50}^{50} j p_{k}(j)^{2}=2 \sum_{j=0}^{50}(j-j) p_{k}(j)^{2}=0 \quad$ and $\quad$ therefore, $\quad \rho_{k}:=\left(\boldsymbol{w}_{k}^{T} A \boldsymbol{w}_{k}\right) /\left(\boldsymbol{w}_{k}^{T} \boldsymbol{w}_{k}\right)=$ $\left(\sum_{j=-50}^{50}(j-50) p_{k}(j)^{2}\right) /\left(\sum_{j=-50}^{50} p_{k}(j)^{2}\right)=-50$. Consider the vector $A \boldsymbol{w}_{k}-\rho_{k} \boldsymbol{w}_{k}$ whose $j$ th component (indexed by $j \in\{-50,-49, \ldots, 50\}$ ) is $(j-50) p_{k}(j)+50 p_{k}(j)=j p_{k}(j)=p_{k+1}(j)$, i.e., $A \boldsymbol{w}_{k}-\rho_{k} \boldsymbol{w}_{k}=\boldsymbol{w}_{k+1}$. Since for the first Arnoldi vector $\boldsymbol{v}_{1}$, there holds $\boldsymbol{v}_{1}=\boldsymbol{b}=\boldsymbol{w}_{0} /\left\|\boldsymbol{w}_{0}\right\|$ and since the Arnoldi vectors satisfy the recursion $\boldsymbol{v}_{k+1}=\left(A \boldsymbol{v}_{k}-\left(\boldsymbol{v}_{k}^{T} A \boldsymbol{v}_{k}\right) \boldsymbol{v}_{k}\right) /\left\|A \boldsymbol{v}_{k}-\left(\boldsymbol{v}_{k}^{T} A \boldsymbol{v}_{k}\right) \boldsymbol{v}_{k}\right\|$, we see that $\boldsymbol{v}_{k}=\boldsymbol{w}_{k} /\left\|\boldsymbol{w}_{k}\right\|$ for all $k$. Thus, $\boldsymbol{v}_{k}^{T} A \boldsymbol{v}_{k}=\left(\boldsymbol{w}_{k}^{T} A \boldsymbol{w}_{k}\right) /\left(\boldsymbol{w}_{k}^{T} \boldsymbol{w}_{k}\right)=-50$, but these are the diagonal entries of $\widehat{H}_{k}$.

[^5]:    ${ }^{5}$ This makes the Euclidean norm of the transformed vectors coincide with the $L^{2}$-norm on the finite element space.

