# Multilevel preconditioning and low rank tensor iteration for space-time simultaneous discretizations of parabolic PDEs 

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# MULTILEVEL PRECONDITIONING AND LOW RANK TENSOR ITERATION FOR SPACE-TIME SIMULTANEOUS DISCRETIZATIONS OF PARABOLIC PDES 

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#### Abstract

This paper addresses the solution of instationary parabolic evolution equations simultaneously in space and time. As a model problem we consider the heat equation posed on the unit cube in Euclidean space of moderately high dimension. An a priori stable minimal residual Petrov-Galerkin variational formulation of the heat equation in space-time results in a generalized least squares problem. This formulation admits a unique, quasi-optimal solution in the natural space-time Hilbert space and serves as a basis for the development of space-time compressive algorithms.

The solution of the heat equation is obtained by applying the conjugate gradient method to the normal equations of the generalized least squares problem. Starting from the well-known BPX preconditioner, multilevel space-time preconditioners for the normal equations are derived. The resulting "parabolic BPX preconditioners" render the normal equations well-conditioned uniformly in the discretization level.

In order to reduce the complexity of the full space-time problem, all computations are performed in a compressed or sparse format called the hierarchical Tucker format supposing that the input data is available in this format. In order to maintain sparsity, compression of the iterates within the hierarchical Tucker format is performed in each conjugate gradient iteration. Its application to vectors in the hierarchical Tucker format is detailed.

Finally, numerical results in up to five spatial dimensions based on the recently developed htucker toolbox for Matlab are presented.


## 1. Introduction

Parabolic evolution equations arise in many applications ranging from engineering to the social sciences. The standard and versatile numerical approach to solving such equations is the method of lines, or time stepping, which either reduces the problem to a system of coupled ordinary differential equations by means of a semidiscretization in space, or to a set of elliptic problems to be solved sequentially by means of a semidiscretization in time, see e.g. [35]. However, two fundamental issues, a practical and a theoretical one, motivate a simultaneous space-time discretization and solution of parabolic equations.
The practical issue is the parallelization of the solution process. An intrinsic limitation to the exploitation of growing parallel computer architectures in time stepping methods is the sequential dependence of the computed solution at intermediate time points on the previous ones. Several methods have been devised to cope with this limitation, we refer to [14] for example. Still, some applications, such as optimal control problems with parabolic PDE constraints, may require the knowledge of the solution to a parabolic problem over the whole time horizon. Since the storage

[^0]of the full solution in space and time may quickly become prohibitive for problems posed in several dimensions, adaptivity that can be performed in parallel and simultaneously in space and time becomes essential.
The theoretical issue is the optimality of the obtained approximate solution. While convergence can be routinely established for time-stepping methods (and the derived methods mentioned in [14]), optimality seems to be exclusive to space-time methods. The reason may be sought in the hidden elliptic character of the parabolic equation which is revealed once it is formulated as a well-posed operator equation in suitable Bochner spaces, see [28, Cha. 3, Sects. 4-7]. Recent numerical methods which build upon this fact are the adaptive wavelet method [34] and the a priori stable (nonadaptive) minimal residual Petrov-Galerkin discretization [1].
The adaptive wavelet method of [34] adopts the abstract operator equation perspective and, after a choice of suitable Riesz bases, reformulates the parabolic equation as an equivalent matrix-vector equation, where the vectors are elements of $\ell^{2}(\mathbb{N})$ and the parabolic operator is represented by a bi-infinite matrix. If the matrix can be suitably approximated by matrices with only finitely many nonzero entries this is the case for many relevant parabolic equations - then the adaptive wavelet methods of $[8,9]$ yield optimal rates of convergence.
While the optimality of the adaptive wavelet method is understood, several of its ingredients are difficult to obtain in practice. This is particularly true for parabolic problems, cf. $[7,6]$. The purpose of this paper is therefore to discuss a practical space-time discretization algorithm based on [1] for the particular example of the heat equation on the unit cube in several dimensions. While parallels to [34] can be identified, the first important and distinctive feature of the present algorithm is a pair of fixed finite-dimensional space-time trial and test spaces of tensor product type. Unlike in the adaptive wavelet method, these are determined a priori and are shown to be stable, i.e., the parabolic operator satisfies an inf-sup condition on those spaces uniformly in the discretization level. A second ingredient is a pair of numerically accessible operators on the test and trial spaces which generate the "correct" norms. From these, a well-posed finite system of least squares type is obtained. The solution satisfies a quasi-optimality estimate in the natural spaces for the continuous equation (in no way dependent on the mesh), akin to the one known as Céa's lemma. This is the second important feature of the present algorithm. The norm generating operators proposed and tested here are derived from the wellknown BPX preconditioner [5, 37, 38], several copies of which are combined into one "parabolic BPX preconditioner" acting in space-time. In this way, the intricate construction of suitable wavelets required in [34] is bypassed. This contributes to the practicality and the novelty of the algorithm.
Adaptivity is crucial in applications with solutions featuring singularities and/or posed in high dimension. As already indicated, we compute in spaces of tensor product type. These are constructed by tensorization of univariate finite-dimensional spaces and the dimension of the trial and test spaces increases rapidly as we only consider uniform mesh refinement. However, at no point in the computation do we require the storage or the computation of the full vector of this size. This presupposes that the input data and the solution may be well-approximated in a low rank tensor format. For the solution, such an approximation is performed adaptively during the iterative solution process. Such formats include the "classical" CP and Tucker formats, see [25] for a review. We choose the more recently developed hierarchical Tucker format [19, 17], since it admits efficient basic operations such as addition, inner product and multiplication by a matrix, as well as approximation of a low rank tensor by a lower-rank tensor.

Methods for the solution of linear systems in the hierarchical Tucker format, as well as in the special case of the Tensor Train format [29], have been the subject of recent work $[3,12,21,23,24,26]$. The first main approach is via developing low rank variants of iterative methods for the solution of linear systems, replacing the iterates by low rank tensors and using the basic operations described above. As the ranks of the iterates tend to grow dramatically during such an iteration, a recurrent approximation of the iterates by tensors of lower rank is employed. A good preconditioner is usually critical to the performance. The second approach are the DMRG or (M)ALS methods, which work directly within the structure of the low rank tensor format, typically the hierarchical Tucker or the Tensor Train format. These methods are based on repeatedly projecting the high-dimensional linear system onto subspaces of lower dimension, and solving the resulting smaller systems. We follow the first route: the proposed BPX preconditioner renders the linear system well conditioned, while the system matrix does not readily lend itself to the projections necessary in the DMRG or (M)ALS methods. In [22, Section 4.3], [15] other low rank tensor approaches for the solution of evolution equations are proposed. These are based on time-stepping or on the explicit representation of the solution using the exponential of the generator.
The outline of the paper is as follows. Section 2 introduces the model heat equation and the space-time variational formulation in certain Bochner spaces $\mathcal{X}$ and $\mathcal{Y}$. Then, a so-called minimal residual Petrov-Galerkin formulation is introduced which allows a stable space-time discretization of the heat equation in finite dimensional spaces. It relies on the availability of certain operators which generate equivalent norms on $\mathcal{X}$ and $\mathcal{Y}$. Following the methodology of operator preconditioning [20], these operators provide preconditioners for the resulting linear system, which has the form of generalized least squares equations, i.e., least squares with minimization w.r.t. specific norms. An example of such operators based on the BPX operator [5] is then constructed. Finite dimensional tensor product space-time test and trial spaces which conform with the minimal residual Petrov-Galerkin formulation are defined. Section 3 introduces the hierarchical Tucker format [17, 19], and describes the Kronecker product structure of the system matrix B resulting from the minimal residual Petrov-Galerkin discretization. The matrix B can be efficiently applied to a vector $\mathbf{x}$ stored as a low rank tensor in the hierarchical Tucker format. Approximate preconditioners are derived that can be applied to such a tensor efficiently. This allows the formulation of a variant of the preconditioned conjugate gradient method for the generalized least squares equations based on the hierarchical Tucker format. Section 4 discusses our numerical examples which comprise isotropic and anisotropic diffusion in the cube $(-1,1)^{d}$ of dimension $d=1, \ldots, 5$. See Section 5 for conclusions and outlook.
Let us comment on the notation used throughout. The symbol $\otimes$ denotes the tensor product of Hilbert spaces [32], as well as the Kronecker product of matrices. A disjoint union of two sets $t$ and $s$ is denoted by $t \dot{\cup} s$. An element $\mathbf{x} \in \mathbb{R}^{n_{0} \times \cdots \times n_{d}}$ is called a tensor. We identify $\mathbf{x}$ with the vector of length $n_{0} \cdots n_{d}$ where convenient. The norm $\|\mathbf{x}\|$ of a tensor $\mathbf{x}$ is defined as the Euclidean norm of the corresponding vector. As a rule, vectors are denoted by lowercase letters and matrices by uppercase letters. For a symmetric positive definite (s.p.d.) matrix M we use the notation $M^{1 / 2}$ to denote a matrix such that $M=M^{\top / 2} M^{1 / 2}$, e.g. the Cholesky factor, where $(\cdot)^{\top}$ denotes transposition.

## 2. Parabolic PDEs

2.1. The model problem. In this section we introduce our model parabolic evolution equation which is the instationary diffusion in a cube. Let $0<T_{\text {final }}<\infty$
and set $J:=\left(0, T_{\text {final }}\right) \subset \mathbb{R}^{1}$. Let $D \subset \mathbb{R}^{d}$, where $d \in \mathbb{N}$ is fixed, be an open "spatial" domain with a Lipschitz boundary. In Section 2.4.2 we will specialize to the case $D=(-1,1)^{d}$. We consider the evolution equation

$$
\begin{align*}
\partial_{t} x(t, \xi)-\operatorname{div}(q(t, \xi) \operatorname{grad} x(t, \xi)) & =p(t, \xi), & (t, \xi) & \in J \times D  \tag{2.1}\\
x(0, \xi) & =h(\xi), & \xi & \in D  \tag{2.2}\\
x(t, \xi) & =0, & (t, \xi) & \in J \times \partial D, \tag{2.3}
\end{align*}
$$

where $q \in L^{\infty}(J \times D)$ is a space- and time-dependent coefficient and the differential operators "div" and "grad" are w.r.t. the spatial variable $\xi \in D$. We assume

$$
\begin{equation*}
0<a_{\min }:=\underset{J \times D}{\operatorname{ess} \inf } q \leq \underset{J \times D}{\operatorname{ess} \sup } q=: a_{\max }<\infty \tag{2.4}
\end{equation*}
$$

We introduce the standard Sobolev spaces $V:=H_{0}^{1}(D)$ and $H:=L^{2}(D)$. The norm on $V$ is given by the $H^{1}(D)$ seminorm. The embedding $V \hookrightarrow H$ is continuous and dense. We denote the scalar product on $H$ by $\langle\cdot, \cdot\rangle_{H}$ and the duality pairing on $V^{\prime} \times V$ by $\langle\cdot, \cdot\rangle_{V^{\prime} \times V}$, and similarly for other Hilbert spaces. We identify $H$ with its dual $H^{\prime}$ via the scalar product $\langle\cdot, \cdot\rangle_{H}$. Thus, one can think of $\langle\cdot, \cdot\rangle_{V^{\prime} \times V}$ as being the unique continuous extension of $\langle\cdot, \cdot\rangle_{H}: V \times V \rightarrow \mathbb{R}$. For (a.e.) $t \in J$ we define $a(t ; \cdot, \cdot): V \times V \rightarrow \mathbb{R}$ by

$$
\begin{equation*}
a(t ; \nu, \bar{\nu})=\int_{D} q(t, \xi) \operatorname{grad} \nu(\xi) \cdot \operatorname{grad} \bar{\nu}(\xi) d \xi \quad \text { for all } \quad \nu, \bar{\nu} \in V \tag{2.5}
\end{equation*}
$$

By (2.4), the symmetric bilinear form $a(t ; \cdot, \cdot)$ is continuous and coercive with constants uniform in (a.e.) $t \in J$. Further, it follows from [13, Lemma 4.4.1] that for all $\nu, \bar{\nu} \in V$, the map $J \ni t \mapsto a(t ; \nu, \bar{\nu})$ is measurable.
In order to obtain a well-posed space-time variational formulation of (2.1)-(2.3) we follow [34] and introduce the Bochner spaces

$$
\begin{aligned}
& \mathcal{X}:=L^{2}(J, V) \cap H^{1}\left(J, V^{\prime}\right) \cong\left(L^{2}(J) \otimes V\right) \cap\left(H^{1}(J) \otimes V^{\prime}\right) \\
& \mathcal{Y}:=\mathcal{Y}_{1} \times \mathcal{Y}_{2}:=L^{2}(J, V) \times H \cong\left(L^{2}(J) \otimes V\right) \times H
\end{aligned}
$$

with norms $\|\cdot\|_{\mathcal{X}}$ and $\|\cdot\|_{\mathcal{Y}}$ given by

$$
\begin{equation*}
\|x\|_{\mathcal{X}}^{2}:=\|x\|_{L^{2}(J, V)}^{2}+\left\|\partial_{t} x\right\|_{L^{2}\left(J, V^{\prime}\right)}^{2}, \quad\|y\|_{\mathcal{Y}}^{2}:=\left\|y_{1}\right\|_{L^{2}(J, V)}^{2}+\left\|y_{2}\right\|_{H}^{2} \tag{2.6}
\end{equation*}
$$

for all $x \in \mathcal{X}$ and $y=\left(y_{1}, y_{2}\right) \in \mathcal{Y}$. One can view $x \in \mathcal{X}$ as either (an equivalence class of functions) $x: J \times D \rightarrow \mathbb{R}$ or as $x: J \rightarrow V^{\prime}$, and we will frequently switch between these two interpretations. Let $\bar{J}$ denote the closure of $J$. We note that

$$
\begin{equation*}
\sup _{t \in \bar{J}} \sup _{x \in \mathcal{X} \backslash\{0\}} \frac{\|x(t)\|_{H}}{\|x\|_{\mathcal{X}}}<\infty \tag{2.7}
\end{equation*}
$$

i.e., $\mathcal{X} \hookrightarrow C^{0}(\bar{J}, H)$ continuously, and in particular, the trace map $\left.(\cdot)\right|_{t=0}: \mathcal{X} \rightarrow H$, $\left.x \mapsto x\right|_{t=0}=x(0) \in H$ is well-defined and continuous [28, Chapter 1]. We define the linear operator $\mathcal{B}: \mathcal{X} \rightarrow \mathcal{Y}^{\prime}$ by

$$
\begin{equation*}
(\mathcal{B} x)(y)=\int_{J}\left\{\left\langle\partial_{t} x(t), y_{1}(t)\right\rangle_{V^{\prime} \times V}+a\left(t ; x(t), y_{1}(t)\right)\right\} d t+\left\langle x(0), y_{2}\right\rangle_{H} \tag{2.8}
\end{equation*}
$$

for all $x \in \mathcal{X}$ and $y=\left(y_{1}, y_{2}\right) \in \mathcal{Y}$, as well as the load functional $b: \mathcal{Y} \rightarrow \mathbb{R}$ by

$$
\begin{equation*}
b(y)=\int_{J}\left\langle p(t), y_{1}(t)\right\rangle_{V^{\prime} \times V} d t+\left\langle h, y_{2}\right\rangle_{H} \tag{2.9}
\end{equation*}
$$

for all $y=\left(y_{1}, y_{2}\right) \in \mathcal{Y}$. It is easy to check that $\mathcal{B} \in \mathcal{L}\left(\mathcal{X}, \mathcal{Y}^{\prime}\right)$ and $b \in \mathcal{Y}^{\prime}$. The space-time variational formulation of (2.1)-(2.3) now reads:

$$
\begin{equation*}
\text { find } \quad x \in \mathcal{X} \quad \text { s.t. } \quad(\mathcal{B} x)(y)=b(y) \quad \forall y \in \mathcal{Y} \tag{2.10}
\end{equation*}
$$

This variational formulation is well-posed by [34, Theorem 5.1], which we recall here for future reference.
Theorem 2.1. The operator $\mathcal{B} \in \mathcal{L}\left(\mathcal{X}, \mathcal{Y}^{\prime}\right)$ is an isomorphism.
2.2. Minimal residual Petrov-Galerkin discretization. In this section we derive a discrete space-time variational formulation of (2.1)-(2.3) from (2.10). It is based on operators $\mathcal{M}$ and $\mathcal{N}$ which generate norms on $\mathcal{X}$ and $\mathcal{Y}$, resp., and provide preconditioners for the discrete system.
Let $\mathcal{M} \in \mathcal{L}\left(\mathcal{X}, \mathcal{X}^{\prime}\right)$ be an operator inducing a scalar product $\langle\cdot, \cdot\rangle_{\mathcal{M}}:=\langle\mathcal{M} \cdot, \cdot\rangle_{\mathcal{X}^{\prime} \times \mathcal{X}}$ on $\mathcal{X}$, and similarly $\mathcal{N} \in \mathcal{L}\left(\mathcal{Y}, \mathcal{Y}^{\prime}\right)$ on $\mathcal{Y}$. The corresponding induced norms on $\mathcal{X}$ and $\mathcal{Y}$ are denoted by $\|\cdot\|_{\mathcal{M}}$ and $\|\cdot\|_{\mathcal{N}}$. We assume the norm equivalences $\|\cdot\|_{\mathcal{M}} \sim$ $\|\cdot\|_{\mathcal{X}}$ and $\|\cdot\|_{\mathcal{N}} \sim\|\cdot\|_{\mathcal{Y}}$. One immediate example is given by the Riesz operators, e.g. if $\mathcal{M}$ is defined by $\langle\mathcal{M} x, \bar{x}\rangle_{\mathcal{X}^{\prime} \times \mathcal{X}}:=\langle x, \bar{x}\rangle_{\mathcal{X}}$ for all $x, \bar{x} \in \mathcal{X}$. Another example based on the well-known BPX operator will be given in Section 2.3 below.
Theorem 2.2. Let $\mathcal{U} \subseteq \mathcal{X}$ and $\mathcal{V} \subseteq \mathcal{Y}$ be closed subspaces. Let $\mathcal{B}$ be given by (2.8). Assume that the discrete inf-sup condition

$$
\begin{equation*}
\gamma_{\mathcal{U}, \mathcal{V}}:=\inf _{u \in \mathcal{U} \backslash\{0\}} \sup _{v \in \mathcal{V} \backslash\{0\}} \frac{(\mathcal{B} u)(v)}{\|u\|_{\mathcal{X}}\|v\|_{\mathcal{Y}}}>0 \tag{2.11}
\end{equation*}
$$

holds. Then there exists a unique $u \in \mathcal{U}$ satisfying

$$
\begin{equation*}
u=\underset{w \in \mathcal{U}}{\arg \min } \sup _{v \in \mathcal{V} \backslash\{0\}} \frac{|(\mathcal{B} w-b)(v)|}{\|v\|_{\mathcal{N}}} . \tag{2.12}
\end{equation*}
$$

Moreover, with $x:=\mathcal{B}^{-1} b$, the quasi-optimality estimate

$$
\begin{equation*}
\|x-u\|_{\mathcal{X}} \leq C \inf _{w \in \mathcal{U}}\|x-w\|_{\mathcal{X}} \quad \text { where } \quad C=\frac{\|\mathcal{B}\|_{\mathcal{L}(\mathcal{X}, \mathcal{Y})}}{\gamma_{\mathcal{U}, \mathcal{V}}} \frac{C_{\mathcal{N}}}{c_{\mathcal{N}}} \geq 0 \tag{2.13}
\end{equation*}
$$

holds, with the constants of the norm equivalence $c_{\mathcal{N}}\|\cdot\|_{\mathcal{Y}} \leq\|\cdot\|_{\mathcal{N}} \leq C_{\mathcal{N}}\|\cdot\|_{\mathcal{Y}}$.
Proof. See [1, Theorem 3.1].
Assume that we are given finite-dimensional subspaces $\mathcal{U} \subseteq \mathcal{X}$ and $\mathcal{V} \subseteq \mathcal{Y}$, as well as bases $\Phi \subset \mathcal{X}$ for $\mathcal{U}$ and $\Psi \subset \mathcal{Y}$ for $\mathcal{V}$. Assume further that the pair $(\mathcal{U}, \mathcal{V})$ satisfies the inf-sup condition (2.11). Set $\mathbb{U}:=\mathbb{R}^{\operatorname{dim} \mathcal{U}}$ and $\mathbb{V}:=\mathbb{R}^{\operatorname{dim} \mathcal{V}}$. Define the matrices $N \in \mathbb{V} \times \mathbb{V}, B \in \mathbb{V} \times \mathbb{U}$ and $M \in \mathbb{U} \times \mathbb{U}$ by

$$
\begin{equation*}
\mathrm{N}:=\langle\mathcal{N} \Psi, \Psi\rangle_{\mathcal{Y}^{\prime} \times \mathcal{Y}}, \quad \mathrm{B}:=\langle\mathcal{B} \Phi, \Psi\rangle_{\mathcal{Y}^{\prime} \times \mathcal{Y}}, \quad \text { and } \quad \mathrm{M}:=\langle\mathcal{M} \Phi, \Phi\rangle_{\mathcal{X}^{\prime} \times \mathcal{X}} \tag{2.14}
\end{equation*}
$$

(i.e., the component $\mathrm{B}_{j i}$ is given by $\left\langle\mathcal{B} \phi_{i}, \psi_{j}\right\rangle_{\mathcal{Y}^{\prime} \times \mathcal{Y}}$ where $\phi_{i} \in \Phi, \psi_{j} \in \Psi$ ) and $\mathbf{b} \in \mathbb{V}$ as the column load vector by $\mathbf{b}:=\langle b, \Psi\rangle_{\mathcal{Y}^{\prime} \times \mathcal{Y}}$. Note that M and N are s.p.d. matrices. We set $\widetilde{B}:=N^{-\top / 2} \mathrm{BM}^{-1 / 2}$ and $\widetilde{\mathbf{b}}:=\mathrm{N}^{-\top / 2} \mathbf{b}$.

Theorem 2.3. With the above definitions the following hold.
(1) The condition number $\kappa_{2}\left(\widetilde{\mathrm{~B}}^{\top} \widetilde{\mathrm{B}}\right)$ w.r.t. the Euclidean norm satisfies

$$
\begin{equation*}
\kappa_{2}\left(\widetilde{\mathrm{~B}}^{\top} \widetilde{\mathrm{B}}\right) \leq C \tag{2.15}
\end{equation*}
$$

where $C \geq 0$ is a monotonic function of $\gamma_{\mathcal{U}, \mathcal{V}}^{-1}\|\mathcal{B}\|_{\mathcal{L}(\mathcal{X}, \mathcal{Y})}$ and the constants in the norm equivalences $\|\cdot\|_{\mathcal{X}} \sim\|\cdot\|_{\mathcal{M}}$ and $\|\cdot\|_{\mathcal{Y}} \sim\|\cdot\|_{\mathcal{N}}$ only.
(2) There exists a unique $\mathbf{u} \in \mathbb{U}$ satisfying

$$
\mathbf{u}=\underset{\mathbf{w} \in \mathbb{U}}{\arg \min }\|B \mathbf{w}-\mathbf{b}\|_{\mathbf{N}^{-1}},
$$

or, equivalently,

$$
\begin{equation*}
\widetilde{\mathrm{B}}^{\top} \widetilde{\mathrm{B}} \widetilde{\mathbf{u}}=\widetilde{\mathrm{B}}^{\top} \widetilde{\mathbf{b}} \text { with } \widetilde{\mathbf{u}}:=\mathrm{M}^{1 / 2} \mathbf{u} . \tag{2.17}
\end{equation*}
$$

(3) The function $u:=\Phi^{\top} \mathbf{u} \in \mathcal{U}$ is characterized by (2.12).

Proof. See [1, Proposition 3.2-3.3].
Motivated by the bound (2.15) on the condition number of the "system matrix" $\widetilde{B}^{\top} \widetilde{B}$ we will apply a version of the conjugate gradient method to the preconditioned normal equations (2.17). This will require an efficient method for the computation of (the action of) the inverses of N and M , given in (2.14); we discuss a particular setting where this is possible in the next section.
2.3. The parabolic BPX preconditioner. The popular BPX preconditioner [5] was developed for second order elliptic partial differential equations, in particular also for anisotropic problems [18], and has been put into greater perspective, see [38, Section 5.4] and references therein. In the optimality of the preconditioner, norm equivalences of the type

$$
\begin{equation*}
\|\nu\|_{V}^{2} \sim \sum_{\ell \in \mathbb{N}_{0}} 2^{2 \ell}\left\|Q_{\ell} \nu\right\|_{H}^{2} \quad \forall \nu \in V, \tag{2.18}
\end{equation*}
$$

where $Q_{\ell}: H \rightarrow H, \ell \in \mathbb{N}_{0}$, are suitable projectors, play an important role. Starting with such norm equivalences in the spatial domain and in the temporal domain, the preconditioner may be adapted to the parabolic operator, which involves different orders of differentiability. We briefly describe the construction of this "parabolic BPX preconditioner" here and refer to [2] for more details. The first set of requirements is given in the following.
(1) In the temporal domain, there exist closed nested subspaces $E_{k} \subseteq E_{k+1} \subseteq$ $H^{1}(J), k \in \mathbb{N}_{0}$, and linear (not necessarily surjective) projectors $P_{k}$ : $H^{1}(J) \rightarrow E_{k}, k \in \mathbb{N}_{0}$, satisfying

$$
\begin{equation*}
\|e\|_{L^{2}(J)}^{2} \sim \sum_{k \in \mathbb{N}_{0}}\left\|P_{k} e\right\|_{L^{2}(J)}^{2} \quad \forall e \in L^{2}(J) \tag{2.19}
\end{equation*}
$$

and

$$
\begin{equation*}
\|e\|_{H^{1}(J)}^{2} \sim \sum_{k \in \mathbb{N}_{0}} 2^{2 k}\left\|P_{k} e\right\|_{L^{2}(J)}^{2} \quad \forall e \in H^{1}(J) . \tag{2.20}
\end{equation*}
$$

Further, we require that $P_{k} P_{k^{\prime}}=0$ for all nonnegative integers $k \neq k^{\prime}$.
(2) In the spatial domain, there exist closed nested subspaces $V_{\ell} \subset V, \ell \in \mathbb{N}_{0}$, with $\{0\}=: V_{-1} \subseteq V_{\ell} \subseteq V_{\ell+1} \subseteq V$ and such that $\bigcup_{\ell \in \mathbb{N}_{0}} V_{\ell}$ is dense in $V$. Further, $H$-orthogonal projectors $Q_{\underline{\ell}}: V \rightarrow H$ are needed, with $\underline{\ell} \in\{\underline{\ell}\}=$ a suitable index set of countable cardinality, such that

$$
\begin{equation*}
d_{V}\|\nu\|_{V}^{2} \leq \sum_{\underline{\ell} \in\{\underline{\ell}\}} q_{\underline{\ell}}^{2}\left\|Q_{\underline{\ell}} \nu\right\|_{H}^{2} \leq D_{V}\|\nu\|_{V}^{2} \quad \forall \nu \in V \tag{2.21}
\end{equation*}
$$

holds with some constants $0<d_{V} \leq D_{V}$ and some $q_{\underline{\ell}} \in \mathbb{R}$, and furthermore, $Q_{\ell} Q_{\ell^{\prime}}=0$ for all $\underline{\ell} \neq \underline{\ell}^{\prime}$. We think of the indices $\underline{\ell}$ as vectors. A possible connection between $Q_{\ell}$ and the subspaces $V_{\ell}$ will be given in Section 2.4.2; it is not needed for the theoretical considerations in this section.
For the following recall that, since the duality pairing $\langle\cdot, \cdot\rangle_{V^{\prime} \times V}$ and the scalar product $\langle\cdot, \cdot\rangle_{H}$ are compatible, i.e., agree on the set where both are defined, so are $\langle\cdot, \cdot\rangle_{\mathcal{X}^{\prime} \times \mathcal{X}}$ and $\langle\cdot, \cdot\rangle_{L^{2}(J ; H)}$, as well as $\langle\cdot, \cdot\rangle_{\mathcal{Y}^{\prime} \times \mathcal{Y}}$ and $\langle\cdot, \cdot\rangle_{\left[L^{2}(J ; H) \times H\right]}$.
An important observation is the fact that the norm equivalence (2.21) in $V$ and $H$-orthogonality of the projectors $Q_{\ell}$ imply a similar norm equivalence for the dual $V^{\prime}$.

Lemma 2.4. With (2.21) we also have

$$
\begin{equation*}
D_{V}^{-1}\|\nu\|_{V^{\prime}}^{2} \leq \sum_{\underline{\ell} \in\{\underline{\ell}\}} q_{\underline{\ell}}^{-2}\left\|Q_{\underline{\ell}} \nu\right\|_{H}^{2} \leq d_{V}^{-1}\|\nu\|_{V^{\prime}}^{2} \quad \forall \nu \in V . \tag{2.22}
\end{equation*}
$$

Proof. Verbatim as in [31, Lemma 1].
For any $e \otimes \nu \in E_{k} \otimes V_{\ell}, k, \ell \in \mathbb{N}_{0}$, we define the operators $\mathcal{M}_{+}$and $\mathcal{M}_{-}$by

$$
\begin{equation*}
\mathcal{M}_{ \pm}(e \otimes \nu):=\sum_{k \in \mathbb{N}_{0}} \sum_{\underline{\ell} \in\{\underline{\ell}\}} g_{k, \underline{\ell}}^{ \pm 1}\left(P_{k} e \otimes Q_{\underline{\ell}} \nu\right), \quad \text { where } \quad g_{k, \underline{\ell}}:=q_{\underline{\ell}}^{2}+2^{2 k} q_{\underline{\ell}}^{-2} \tag{2.23}
\end{equation*}
$$

and their extension to $\bigcup_{k, \ell \in \mathbb{N}_{0}} E_{k} \otimes V_{\ell}$ by linearity. Using (2.19)-(2.20) and (2.21) it can be shown that

$$
\begin{equation*}
\left\langle\mathcal{M}_{+} x, x\right\rangle_{\mathcal{X}^{\prime} \times \mathcal{X}} \sim\|x\|_{\mathcal{X}}^{2} \quad \forall x \in \bigcup_{k, \ell \in \mathbb{N}_{0}} E_{k} \otimes V_{\ell} \tag{2.24}
\end{equation*}
$$

holds with constants uniform in $x$. Moreover, we have $\mathcal{M}_{-} \mathcal{M}_{+} x=x$. Since $\sum_{k, \ell \in \mathbb{N}_{0}} E_{k} \otimes V_{\ell}$ is dense in $\mathcal{X}$, where the sum is direct and orthogonal in $\mathcal{X}$, the operator $\mathcal{M}_{+}$extends uniquely by linearity and continuity to an operator (still denoted by) $\mathcal{M}_{+} \in \mathcal{L}\left(\mathcal{X}, \mathcal{X}^{\prime}\right)$, and $\mathcal{M}_{-}$extends to the inverse thereof, $\mathcal{M}_{-}=$ $\mathcal{M}_{+}^{-1} \in \mathcal{L}\left(\mathcal{X}^{\prime}, \mathcal{X}\right)$.
In order to obtain a pair of subspaces $\mathcal{U}$ and $\mathcal{V}$ satisfying (2.11), we further require closed subspaces $F_{k} \subset L^{2}(J), k \in \mathbb{N}_{0}$, such that $E_{k} \subseteq F_{k}, k \in \mathbb{N}_{0}$, and

$$
\begin{equation*}
\inf _{k \in \mathbb{N}_{0}} \inf _{e^{\prime} \in E_{k}^{\prime} \backslash\{0\}} \sup _{f \in F_{k} \backslash\{0\}} \frac{\left\langle e^{\prime}, f\right\rangle_{L^{2}(J)}}{\left\|e^{\prime}\right\|_{L^{2}(J)}\|f\|_{L^{2}(J)}}>0 \tag{2.25}
\end{equation*}
$$

where $E_{k}^{\prime}:=\left\{e^{\prime}: e \in E_{k}\right\}$. For any $(f \otimes \nu, h) \in\left[F_{k} \otimes V_{\ell}\right] \times H, k, \ell \in \mathbb{N}_{0}$, we then define the operator $\mathcal{N}_{ \pm}$by

$$
\begin{equation*}
\mathcal{N}_{ \pm}(f \otimes \nu, h):=\left(f \otimes \sum_{\underline{\ell} \in\{\underline{\ell}\}} q_{\underline{\ell}}^{ \pm 2} Q_{\underline{\ell}} \nu, h\right) \tag{2.26}
\end{equation*}
$$

and the extension to $\bigcup_{k, \ell \in \mathbb{N}_{0}}\left[E_{k} \otimes V_{\ell}\right] \times V_{\ell}$ by linearity. As above, we have

$$
\begin{equation*}
\left\langle\mathcal{N}_{+} y, y\right\rangle_{\mathcal{Y}^{\prime} \times \mathcal{Y}} \sim\|y\|_{\mathcal{Y}}^{2} \quad \forall y=\left(y_{1}, y_{2}\right) \in \bigcup_{k, \ell \in \mathbb{N}_{0}}\left[E_{k} \otimes V_{\ell}\right] \times V_{\ell} \tag{2.27}
\end{equation*}
$$

uniformly in $y$, and further, $\mathcal{N}_{+}$extends uniquely by linearity and continuity to an operator (still denoted by) $\mathcal{N}_{+} \in \mathcal{L}\left(\mathcal{Y}, \mathcal{Y}^{\prime}\right)$, and $\mathcal{N}_{-}$extends to the inverse thereof, $\mathcal{N}_{-}=\mathcal{N}_{+}^{-1} \in \mathcal{L}\left(\mathcal{Y}^{\prime}, \mathcal{Y}\right)$.
Using the notation from Section 2.2, we define the matrices $\mathrm{M}_{ \pm} \in \mathbb{U} \times \mathbb{U}$ and $\mathrm{N}_{ \pm} \in \mathbb{V} \times \mathbb{V}$ by

$$
\begin{equation*}
\mathrm{M}_{ \pm}:=\left\langle\mathcal{M}_{ \pm} \Phi, \Phi\right\rangle_{\mathcal{X}^{\prime} \times \mathcal{X}} \quad \text { and } \quad \mathrm{N}_{ \pm}:=\left\langle\mathcal{N}_{ \pm} \Psi, \Psi\right\rangle_{\mathcal{Y}^{\prime} \times \mathcal{Y}}, \tag{2.28}
\end{equation*}
$$

as well as the mass matrices $\mathrm{M}_{0} \in \mathbb{U} \times \mathbb{U}$ and $\mathrm{N}_{0} \in \mathbb{V} \times \mathbb{V}$ by

$$
\begin{equation*}
\mathrm{M}_{0}:=\langle\Phi, \Phi\rangle_{\mathcal{X}^{\prime} \times \mathcal{X}} \quad \text { and } \quad \mathrm{N}_{0}:=\langle\Psi, \Psi\rangle_{\mathcal{Y}^{\prime} \times \mathcal{Y}} \tag{2.29}
\end{equation*}
$$

We will apply the results of Section 2.2 with $\mathcal{M}:=\mathcal{M}_{+}$and $\mathcal{N}:=\mathcal{N}_{+} ;$consequently, $\mathrm{M} \equiv \mathrm{M}_{+}$and $\mathrm{N} \equiv \mathrm{N}_{+}$. The following observation will therefore be important for the efficient application of the inverse of the matrices $M_{+}$and $N_{+}$required in the resolution of (2.17).

Proposition 2.5. It holds

$$
\begin{equation*}
\mathrm{M}_{+}^{-1}=\mathrm{M}_{0}^{-1} \mathrm{M}_{-} \mathrm{M}_{0}^{-1} \quad \text { and } \quad \mathrm{N}_{+}^{-1}=\mathrm{N}_{0}^{-1} \mathrm{~N}_{-} \mathrm{N}_{0}^{-1} . \tag{2.30}
\end{equation*}
$$

Proof. We show that $\mathrm{M}_{0}=\mathrm{M}_{-} \mathrm{M}_{0}^{-1} \mathrm{M}_{+}$. Indeed, observe that

$$
\begin{align*}
\overline{\mathbf{x}}^{\top} \mathrm{M}_{0} \mathbf{x}=\langle\bar{x}, x\rangle_{\mathcal{X}^{\prime} \times \mathcal{X}} & =\left\langle\mathcal{M}_{+}^{-1} \bar{x}, \mathcal{M}_{+} x\right\rangle_{\mathcal{X}^{\prime} \times \mathcal{X}}=\left\langle\mathcal{M}_{-} \bar{x}, \mathcal{M}_{+} x\right\rangle_{\mathcal{X}^{\prime} \times \mathcal{X}}  \tag{2.31}\\
& =\left(\mathrm{M}_{0}^{-1} \mathrm{M}_{-} \overline{\mathbf{x}}\right)^{\top} \mathrm{M}_{0}\left(\mathrm{M}_{0}^{-1} \mathrm{M}_{+} \mathbf{x}\right)=\overline{\mathbf{x}}^{\top} \mathrm{M}_{-} \mathrm{M}_{0}^{-1} \mathrm{M}_{+} \mathbf{x} \tag{2.32}
\end{align*}
$$

holds for all $\bar{x}=\Phi^{\top} \overline{\mathbf{x}}, x=\Phi^{\top} \mathbf{x} \in \mathcal{U}$. Analogously for $\mathbf{N}_{+}^{-1}$.
2.4. Space-time test and trial spaces. In order to apply the framework of Section 2.2 we construct finite dimensional subspaces $\mathcal{U} \subseteq \mathcal{X}$ and $\mathcal{V} \subseteq \mathcal{Y}$ for which the inf-sup condition (2.11) can be verified. These are constructed as a tensor product of univariate spaces.
2.4.1. Discretization in time. We define the temporal mesh as

$$
\begin{equation*}
\mathcal{T}_{k}:=\left\{i 2^{-(k+1)} T_{\text {final }}: i=0, \ldots, 2^{k+1}\right\} . \tag{2.33}
\end{equation*}
$$

The space $E_{k} \subset H^{1}(J)$ is defined as the space of continuous, $\mathcal{T}_{k}$-piecewise linear functions on $J$ with the convention $E_{-1}:=\{0\}$. Note that $\mathcal{T}_{k} \subset \mathcal{T}_{k+1}, k \in \mathbb{N}_{0}$, hence the nestedness property $E_{k} \subset E_{k+1}, k \in \mathbb{N}_{0}$. Due to $\sup _{t \in J} \operatorname{dist}\left(t, \mathcal{T}_{k}\right)=\mathcal{O}\left(2^{-k}\right)$, $k \rightarrow \infty$, the Jackson inequality

$$
\begin{equation*}
\inf _{e \in E_{k}}\|f-e\|_{L^{2}(J)} \lesssim 2^{-k \tau}\|f\|_{H^{\tau}(J)} \quad \forall f \in H^{\tau}(J) \quad \forall k \in \mathbb{N}_{0} \tag{2.34}
\end{equation*}
$$

and the Bernstein inequality

$$
\begin{equation*}
\|e\|_{H^{\tau}(J)} \lesssim 2^{k \tau}\|e\|_{L^{2}(J)} \quad \forall e \in E_{k} \quad \forall k \in \mathbb{N}_{0} \tag{2.35}
\end{equation*}
$$

hold for $\tau=0,1,2$ with the implied constant independent of $f$ and $k$, resp. $e$ and $k$, as in the case of the uniform mesh $\mathcal{T}_{k}^{1}$. Thus, [11, Theorem 3.2] implies that (2.19)-(2.20) hold with $P_{k}: H^{1}(J) \rightarrow E_{k} \cap\left(E_{k-1}\right)^{\perp_{L^{2}(J)}}$ being the $L^{2}(J)$-orthogonal projector, which we assume for $P_{k}$ from now on.
Starting from the sequence $E_{k}, k \in \mathbb{N}_{0}$, we define the spaces $F_{k}, k \in \mathbb{N}_{0}$, required for the construction of the operator $\mathcal{N}$ in (2.27), as $F_{k}:=E_{k+1}$. This choice is motivated by the following result, see [1, Proposition 6.1] for the proof.

Proposition 2.6. With $E_{k}$ and $F_{k}$ as above, (2.25) holds.
2.4.2. Discretization in space. We specialize the discussion to the case

$$
\begin{equation*}
V:=H_{0}^{1}(D) \quad \text { for } \quad D:=(-1,1)^{d} \tag{2.36}
\end{equation*}
$$

where $d \in \mathbb{N}$. For each $\mu=1, \ldots, d$ and $\ell \in \mathbb{N}_{0}$, we take $V_{\ell}^{(\mu)} \subset H_{0}^{1}(-1,1)$ as the standard conforming finite element space defined as the space of all continuous, piecewise linear functions w.r.t. the uniform partition of the interval $(-1,1) \subset \mathbb{R}^{1}$ into $2^{\ell+1}$ subintervals. For $\underline{\ell}_{\mu} \in \mathbb{N}_{0}, \mu=1, \ldots, d$, we define $V_{\underline{\ell}}:=V_{\ell_{1}}^{(1)} \otimes \cdots \otimes V_{\ell_{d}}^{(d)} \subset$ $V$. Further, we set $V_{L}:=V_{(L, \ldots, L)}, L \in \mathbb{N}_{0}$. For each $\underline{\ell} \in \mathbb{N}_{0}^{d}$ we define $Q_{\underline{\ell}}: V \rightarrow H$ as the $H$-orthogonal projector

$$
Q_{\underline{\ell}}: V \rightarrow V_{\underline{\ell}} \cap\left(\sum_{\underline{\ell} \neq \underline{\ell}^{\prime} \leq \underline{\ell}} V_{\underline{\ell}^{\prime}}\right)^{\perp_{H}}
$$

where the sum runs over all $\underline{\ell}_{\mu}^{\prime} \in \mathbb{N}_{0}, \mu=1, \ldots, d$, satisfying the constraint, and an empty sum evaluates to the trivial vector space $\{0\} \subset V$. Since (2.21) holds for $d=1$, see [10, Theorem 5.8], from

$$
\left\|\nu_{1} \otimes \cdots \otimes \nu_{d}\right\|_{V}^{2}=\sum_{\mu=1}^{d}\left(\left\|\nu_{\mu}\right\|_{H_{0}^{1}(-1,1)}^{2} \prod_{\mu \neq \mu^{\prime}=1}^{d}\left\|\nu_{\mu^{\prime}}\right\|_{L^{2}(-1,1)}^{2}\right)
$$

it follows (cf. [18, Theorem 3]) that (2.21) holds for $d \geq 1$ with the same constants $0<d_{V} \leq D_{V}<\infty$ for

$$
\begin{equation*}
q_{\underline{\ell}}:=\sqrt{\sum_{\mu=1}^{d} 2^{2 \underline{\ell}_{\mu}}} . \tag{2.37}
\end{equation*}
$$

2.4.3. Tensor product space-time spaces. One possible construction of finite element spaces $\mathcal{U} \subset \mathcal{X}$ and $\mathcal{V} \subset \mathcal{Y}$ satisfying the inf-sup condition (2.11) can be obtained from the following result.

Theorem 2.7. With the notation introduced above, there exists $\gamma>0$ such that the pair

$$
\begin{equation*}
\mathcal{U}:=E_{K} \otimes V_{L} \quad \text { and } \quad \mathcal{V}:=\left(F_{K} \otimes V_{L}\right) \times V_{L} \tag{2.38}
\end{equation*}
$$

satisfies $\gamma_{\mathcal{U}, \mathcal{V}} \geq \gamma>0$ for all $K \in \mathbb{N}_{0}$ and $L \in \mathbb{N}_{0}$.
Proof. The claim follows from (2.25) and the following corresponding property of the spaces $V_{L}$

$$
\begin{equation*}
\inf _{L \in \mathbb{N}_{0}} \inf _{\nu^{\prime} \in V_{L} \backslash\{0\}} \sup _{\nu \in V_{L} \backslash\{0\}} \frac{\left\langle\nu^{\prime}, \nu\right\rangle_{V^{\prime} \times V}}{\left\|\nu^{\prime}\right\|_{V^{\prime}}\|\nu\|_{V}}>0 . \tag{2.39}
\end{equation*}
$$

For details we refer to [1, Proof of Theorem 6.3].

## 3. Tensor format

With the space-time discretization proposed in Section 2.4, the solution to the parabolic equation is approximated by a vector $\mathbf{x}$, which can be interpreted as a high-dimensional array, or a tensor. The tensor $\mathbf{x} \in \mathbb{R}^{N_{t} \times N_{\xi} \times \cdots \times N_{\xi}}$ is of order $d+1$, where $N_{t}$ represents the number of basis functions in time, and $N_{\xi}$ the number of basis functions in space. Hence, the storage cost for $\mathbf{x}$ increases exponentially in $d$. We reduce this cost by approximating $\mathbf{x}$ in a low rank tensor format, described in Section 3.1. We show in Section 3.3 that the matrices $B, \mathrm{M}_{+}^{-1}$ and $\mathrm{N}_{+}^{-1}$ can be efficiently applied to a tensor in such a format. Using a variant of the conjugate gradient method, an approximation of $\mathbf{x}$ in the low rank tensor format can be computed as discussed in Section 3.4.
We now consider a general tensor $\mathbf{x} \in \mathbb{R}^{n_{1} \times \cdots \times n_{d}}$ of order $d \in \mathbb{N}$. Low rank formats provide an approximation of $\mathbf{x}$, similar to the truncated SVD for $d=2$. Consider the case of a matrix $\mathbf{x} \in \mathbb{R}^{n_{1} \times n_{2}}$. If the matrix $\mathbf{x}$ has a steep singular value decay, it can be approximated by a low rank matrix $U V^{\top} \approx \mathrm{x}$ with $U \in \mathbb{R}^{n_{1} \times r}$ and $V \in \mathbb{R}^{n_{2} \times r}$. Basic operations such as addition, multiplication by a matrix or a scalar and inner product of two low rank matrices $\mathbf{x}=U_{\mathbf{x}} V_{\mathbf{x}}^{\top}$ and $\mathbf{y}=U_{\mathbf{y}} V_{\mathbf{y}}^{\top}$ can be performed efficiently while preserving the low rank structure:

$$
\begin{align*}
\mathbf{x}+\mathbf{y} & =U_{\mathbf{x}} V_{\mathbf{x}}^{\top}+U_{\mathbf{y}} V_{\mathbf{y}}^{\top}=\left[U_{\mathbf{x}} U_{\mathbf{y}}\right]\left[V_{\mathbf{x}} V_{\mathbf{y}}\right]^{\top}  \tag{3.1}\\
A \mathbf{x} & =A\left(U V^{\top}\right)=(A U) V^{\top}, \quad \mathbf{x} B^{\top}=U V^{\top} B^{\top}=U(B V)^{\top}  \tag{3.2}\\
\langle\mathbf{x}, \mathbf{y}\rangle & :=\sum_{i_{1}=1}^{n_{1}} \sum_{i_{2}=1}^{n_{2}} \mathbf{x}_{i_{1}, i_{2}} \mathbf{y}_{i_{1}, i_{2}}=\left\langle U_{\mathbf{x}} V_{\mathbf{x}}^{\top}, U_{\mathbf{y}} V_{\mathbf{y}}^{\top}\right\rangle=\left\langle U_{\mathbf{x}}^{\top} U_{\mathbf{y}}, V_{\mathbf{x}}^{\top} V_{\mathbf{y}}\right\rangle . \tag{3.3}
\end{align*}
$$

Moreover, for a matrix $U V^{\top}$ of rank $r$, an approximation with lower rank $\widetilde{r} \leq r$ can be efficiently computed using the truncated SVD:

$$
U=Q_{U} R_{U}, V=Q_{V} R_{V} \quad \Rightarrow \quad U V^{\top}=Q_{U}\left(R_{U} R_{V}^{\top}\right) Q_{V}^{\top}=\left(Q_{U} X\right) \Sigma\left(Q_{V} Y\right)^{\top}
$$

where the singular value decomposition $R_{U} R_{V}^{\top}=X \Sigma Y^{\top}$ was used in the last equation. The best approximation at rank $\widetilde{r} \leq r$ w.r.t. the Frobenius norm is obtained by setting the diagonal entries of $\Sigma$ at the locations $\widetilde{r}+1, \ldots, r$ to zero. These operations allow for iterative solution algorithms, such as the Richardson or the conjugate gradient method, to be performed in low rank formats (see Section 3.4).

Let us now consider the general case $d \geq 2$, that is, approximating a tensor $\mathbf{x} \in$ $\mathbb{R}^{n_{1} \times \cdots \times n_{d}}$ in a low rank format. The CP and Tucker decompositions are wellknown low rank tensor formats, see [25] for a review. Neither of these formats
is well suited for our setting. For example, in the Tucker decomposition storage requirements grow exponentially with the order $d$ of the tensor. Therefore, we choose the $\mathcal{H}$-Tucker format [17, 19]. A special case of this format, called Tensor Train format, was independently proposed in [29].
3.1. The hierarchical Tucker format. Here we give a brief overview of the hierarchical Tucker ( $\mathcal{H}$-Tucker) format, and refer the reader to [17, 19] for a more detailed explanation.
In order to describe the $\mathcal{H}$-Tucker format, we introduce the concept of matricization of a tensor. Consider a splitting of the dimensions into two disjoint sets: $t \dot{\cup} s=$ $\{1, \ldots, d\}$ with $t=\left\{t_{1}, \ldots, t_{k}\right\}$ and $s=\left\{s_{1}, \ldots, s_{d-k}\right\}$. The matricization $X^{(t)}$ of a tensor $\mathbf{x}$ with respect to $t$ is obtained by merging the first set $t$ of modes into row indices and the second set $s$ into column indices:

$$
\left.X^{(t)} \in \mathbb{R}^{\left(n_{t_{1}} \cdots n_{t_{k}}\right.}\right) \times\left(n_{s_{1}} \cdots n_{s_{d-k}}\right) \quad \text { with } \quad\left(X^{(t)}\right)_{\left(i_{t_{1}}, \ldots, i_{t_{k}}\right),\left(i_{s_{1}}, \ldots, i_{s_{d-k}}\right)}:=\mathbf{x}_{i_{1}, \ldots, i_{d}}
$$

for any multiindex $\left(i_{1}, \ldots, i_{d}\right) \in\left\{1, \ldots, n_{1}\right\} \times \cdots \times\left\{1, \ldots, n_{d}\right\}$. In this notation, the vectorization $\operatorname{vec}(\mathbf{x})$ is just the matricization with $t=\{1, \ldots, d\}$ and $s=$ Ø. Consider a collection $T \subseteq 2^{\{1, \ldots, d\}}$ of subsets $t \subset\{1, \ldots, d\}$. We define the hierarchical rank $r_{t} \in \mathbb{N}_{0}$ for all $t \in T$, and the corresponding set of $\mathcal{H}$-Tucker tensors as

$$
\begin{equation*}
\mathcal{H}-\operatorname{Tucker}\left(\left(r_{t}\right)_{t \in T}\right):=\left\{\mathbf{x} \in \mathbb{R}^{n_{1} \times \cdots \times n_{d}}: \operatorname{rank}\left(X^{(t)}\right) \leq r_{t} \forall t \in T\right\} \tag{3.4}
\end{equation*}
$$

For each set $t \in T$, there exist matrices $U_{t} \in \mathbb{R}^{\left(n_{t_{1}} \cdots n_{t_{k}}\right) \times r_{t}}$ and $V_{t} \in \mathbb{R}^{\left(n_{s_{1}} \cdots n_{s_{d-k}}\right) \times r_{t}}$ such that $X^{(t)}=U_{t} V_{t}^{\top}$. The nestedness property Range $\left(U_{t}\right) \subseteq$ Range $\left(U_{t_{r}} \otimes U_{t_{l}}\right)$ holds for each disjoint splitting $t_{l} \cup t_{r}=t$ [16, Lemma 17], which implies that there exists a matrix $B_{t}$ such that $U_{t}=\left(U_{t_{r}} \otimes U_{t_{l}}\right) B_{t}$. Consequently, it is sufficient to store $U_{t_{l}}, U_{t_{r}}$ and $B_{t}$ to be able to represent $U_{t}$. This property, applied recursively, allows a storage-efficient representation of $\mathcal{H}$-Tucker tensors. Consider the example of $d=4$ :

$$
\begin{align*}
\operatorname{vec}(\mathbf{x})=X^{(1234)} & =\left(U_{34} \otimes U_{12}\right) B_{1234}  \tag{3.5}\\
U_{12} & =\left(U_{2} \otimes U_{1}\right) B_{12}  \tag{3.6}\\
U_{34} & =\left(U_{4} \otimes U_{3}\right) B_{34}  \tag{3.7}\\
\Rightarrow \quad \operatorname{vec}(\mathbf{x}) & =\left(U_{4} \otimes U_{3} \otimes U_{2} \otimes U_{1}\right)\left(B_{34} \otimes B_{12}\right) B_{1234} . \tag{3.8}
\end{align*}
$$

In the $\mathcal{H}$-Tucker format, $T$ is allowed to be any binary tree with the root node $t_{\text {root }}=\{1, \ldots, d\}$, leaf nodes $t_{\text {leaf }}$ containing only one element, and all other nodes $t$ having exactly two children $t_{l}, t_{r}$ with $t_{l} \dot{\cup} t_{r}=t$.
The storage requirements are $\mathcal{O}\left(d n r+d r^{3}\right)$, for $n:=\max _{\mu=1, \ldots, d} n_{\mu}$ and $r:=$ $\max _{t \in T} r_{t}$. Addition, multiplication by a matrix and inner product of two tensors in $\mathcal{H}$-Tucker format can be performed efficiently, we refer the reader to $[27,36]$ for details on the implementation. The truncation of a tensor $\mathbf{x}$ in $\mathcal{H}$-Tucker format is the approximation by a tensor $\widetilde{\mathbf{x}}$ with given lower hierarchical ranks $\left(\widetilde{r}_{t}\right)_{t \in T}$, where $\widetilde{r}_{t} \leq r_{t}, t \in T$. An implementation of this operation that requires $\mathcal{O}\left(d n r^{2}+d r^{4}\right)$ floating point operations, and satisfies the quasi-optimality property [17]

$$
\begin{equation*}
\|\mathbf{x}-\widetilde{\mathbf{x}}\| \leq \sqrt{2 d-3} \inf \left\{\|\mathbf{x}-\mathbf{y}\|: \mathbf{y} \in \mathcal{H}-\operatorname{Tucker}\left(\left(\widetilde{r}_{t}\right)_{t \in T}\right)\right\} \tag{3.9}
\end{equation*}
$$

is possible. Given two parameters, a relative truncation accuracy rel_eps $>0$ and a maximal truncation rank max_rank $\in \mathbb{N}$, we define the truncation $\mathcal{T}$ which chooses the ranks $\left(\widetilde{r}_{t}\right)_{t \in T}$ adaptively such that $\|\mathbf{x}-\widetilde{\mathbf{x}}\| \leq$ rel_eps $\|\mathbf{x}\|$ if this is possible with $\widetilde{r}_{t} \leq$ max_rank for all $t \in T$; otherwise, some of the ranks are set to max_rank, and the relative accuracy requirement may not be satisfied. For details we refer to $[17,27,36]$. The numerical experiments below are based on the publicly available
htucker toolbox [27, 36] in Matlab which provides a representation of a tensor in $\mathcal{H}$-Tucker format, as well as the operations described above.
3.2. Discretized generalized linear least squares problem. We assume that the functions $q, p, h$ are separable, with

$$
q(t, \xi)=q_{0}(t) \prod_{\mu=1}^{d} q_{\mu}\left(\xi_{\mu}\right), \quad p(t, \xi)=p_{0}(t) \prod_{\mu=1}^{d} p_{\mu}\left(\xi_{\mu}\right) \quad \text { and } \quad h(\xi)=\prod_{\mu=1}^{d} h_{\mu}\left(\xi_{\mu}\right) .
$$

The extension to finite sums of separable functions is straightforward. Let $K \in \mathbb{N}_{0}$ and $L \in \mathbb{N}_{0}$ be fixed and consider the spaces $\mathcal{U}$ and $\mathcal{V}$ defined in (2.38). Let $b_{i}^{\mu}, i=1, \ldots, 2^{L}-1$, be the standard hat functions spanning the space $V_{L}^{(\mu)}$, $\mu=1, \ldots, d$. Let $e_{i}, i=1, \ldots, 2^{K}+1$, be the standard hat functions spanning the space $E_{K}$, and similarly $f_{j}, j=1, \ldots, 2^{K+1}+1$ for $F_{K}$. Then the basis $\Phi$ for $\mathcal{U}$ is given by the functions $e_{i_{0}}^{0} \otimes b_{i_{1}}^{1} \otimes \cdots \otimes b_{i_{d}}^{d}$, and the basis $\Psi$ for $\mathcal{V}$ consists of $f_{j_{0}}^{0} \otimes b_{i_{1}}^{1} \otimes \cdots \otimes b_{i_{d}}^{d}$, where $i_{0}=1, \ldots, 2^{K}+1, j_{0}=1, \ldots, 2^{K+1}+1$ and $i_{\mu}=1, \ldots, 2^{L}-1, \mu=1, \ldots, d$. As described above in Section 2.2, spacetime minimal residual Petrov-Galerkin discretization w.r.t. those bases leads to the generalized linear least-squares problem $\arg \min _{\mathbf{u}}\|B \mathbf{B}-\mathbf{b}\|_{N^{-1}}$, where the matrix $B$ has the form

$$
\binom{C^{(0)} \otimes M^{(1)} \otimes \cdots \otimes M^{(d)}+\sum_{\mu=1}^{d} \widehat{M}^{(0)} \otimes \widehat{M}^{(1)} \otimes \cdots \otimes A^{(\mu)} \otimes \cdots \otimes \widehat{M}^{(d)}}{\left(c^{(0)}\right)^{\top} \otimes M^{(1)} \otimes \cdots \otimes M^{(d)}}
$$

and

$$
\mathbf{b}=\binom{p^{(0)} \otimes p^{(1)} \otimes \cdots \otimes p^{(d)}}{1 \otimes h^{(1)} \otimes \cdots \otimes h^{(d)}} .
$$

Here, and in the following, the superscript in vectors and matrices refers to the mode $\mu=0, \ldots, d$. These equations can be combined into one Kronecker product structure by using concatenation in the first mode:

$$
\begin{align*}
& \mathrm{B}=\widetilde{C}^{(0)} \otimes M^{(1)} \otimes \cdots \otimes M^{(d)}+\sum_{\mu=1}^{d} \widetilde{M}^{(0)} \otimes \widehat{M}^{(1)} \otimes \cdots \otimes A^{(\mu)} \otimes \cdots \otimes \widehat{M}^{(d)}  \tag{3.10}\\
& \mathbf{b}=\widetilde{p}^{(0)} \otimes p^{(1)} \otimes \cdots \otimes p^{(d)}+\widetilde{h}^{(0)} \otimes h^{(1)} \otimes \cdots \otimes h^{(d)} \tag{3.11}
\end{align*}
$$

with block matrices $\widetilde{C}^{(0)}, \widetilde{M}^{(0)}$ and block vectors $\widetilde{p}^{(0)}, \widetilde{h}^{(0)}$, given by

$$
\widetilde{C}^{(0)}=\binom{C^{(0)}}{\left(c^{(0)}\right)^{\top}}, \quad \widetilde{M}^{(0)}=\binom{\widehat{M}^{(0)}}{0}, \widetilde{p}^{(0)}=\binom{p^{(0)}}{0}, \widetilde{h}^{(0)}=\binom{0}{1} .
$$

The matrices for the spatial domain are defined as follows for $\mu=1, \ldots, d$ :

$$
\begin{gather*}
A_{i j}^{(\mu)}=\int_{-1}^{1} q^{(\mu)}\left(\xi_{\mu}\right) \operatorname{grad} \nu_{i}^{(\mu)}\left(\xi_{\mu}\right) \cdot \operatorname{grad} \nu_{j}^{(\mu)}\left(\xi_{\mu}\right) \mathrm{d} \xi_{\mu},  \tag{3.12}\\
M_{i j}^{(\mu)}=\int_{-1}^{1} \nu_{i}^{(\mu)}\left(\xi_{\mu}\right) \nu_{j}^{(\mu)}\left(\xi_{\mu}\right) \mathrm{d} \xi_{\mu}, \quad \widehat{M}_{i j}^{(\mu)}=\int_{-1}^{1} q_{\mu}\left(\xi_{\mu}\right) \nu_{i}^{(\mu)}\left(\xi_{\mu}\right) \nu_{j}^{(\mu)}\left(\xi_{\mu}\right) \mathrm{d} \xi_{\mu},  \tag{3.13}\\
p_{i}^{(\mu)}=\int_{-1}^{1} p_{\mu}\left(\xi_{\mu}\right) \nu_{i}^{(\mu)}\left(\xi_{\mu}\right) \mathrm{d} \xi_{\mu}, \quad h_{i}^{(\mu)}=\int_{-1}^{1} h_{\mu}\left(\xi_{\mu}\right) \nu_{i}^{(\mu)}\left(\xi_{\mu}\right) \mathrm{d} \xi_{\mu} . \tag{3.14}
\end{gather*}
$$

For the temporal domain, the following matrices are needed:

$$
\begin{gather*}
C_{i j}^{(0)}=\int_{0}^{T} f_{i}(t) \frac{\mathrm{d} e_{j}}{\mathrm{~d} t}(t) \mathrm{d} t, \quad \widehat{M}_{i j}^{(0)}=\int_{0}^{T} q^{(0)}(t) f_{i}(t) e_{j}(t) \mathrm{d} t  \tag{3.15}\\
c_{j}^{(0)}=e_{j}(0), \quad p_{i}^{(0)}=\int_{0}^{T} p^{(0)}(t) f_{i}(t) \mathrm{d} t  \tag{3.16}\\
\left(M_{e}^{(0)}\right)_{i j}=\int_{0}^{T} e_{i}(t) e_{j}(t) \mathrm{d} t, \quad\left(M_{f}^{(0)}\right)_{i j}=\int_{0}^{T} f_{i}(t) f_{j}(t) \mathrm{d} t \tag{3.17}
\end{gather*}
$$

Observe that there are $2^{K+1}+1$ rows, and $2^{K}+1$ columns in $\widetilde{C}^{(0)}$ and $\widetilde{M}^{(0)}$. With this notation, the mass matrices $\mathrm{M}_{0}$ and $\mathrm{N}_{0}$ are given by:

$$
\begin{align*}
& \mathrm{M}_{0}=M_{e}^{(0)} \otimes M^{(1)} \otimes \cdots \otimes M^{(d)},  \tag{3.18}\\
& \mathrm{N}_{0}=\left(\begin{array}{cc}
M_{f}^{(0)} & 0 \\
0 & 1
\end{array}\right) \otimes M^{(1)} \otimes \cdots \otimes M^{(d)}=:\left(\begin{array}{cc}
\mathrm{N}_{0,1} & 0 \\
0 & \mathrm{~N}_{0,2}
\end{array}\right) . \tag{3.19}
\end{align*}
$$

3.3. Application of matrices to a tensor in low rank format. Note that the application of $A_{0} \otimes A_{1} \otimes \cdots \otimes A_{d}$ to a tensor $\mathbf{x}$ in $\mathcal{H}$-Tucker format preserves its hierarchical ranks. Therefore, given a tensor $\mathbf{x}$ in $\mathcal{H}$-Tucker format with hierarchical rank $r_{t}$, a straightforward application of B as defined in (3.10) to $\mathbf{x}$ will result in a tensor with hierarchical ranks at most $(d+1) r_{t}$. This increases the storage cost by a factor of $(d+1)^{3}$. However, this problem can be alleviated by writing B in the following form (as proposed for the TT format in [29]):

$$
\begin{equation*}
\sum_{j_{0}=1}^{J_{0}} \cdots \sum_{j_{d}=1}^{J_{d}} \mathbf{h}_{j_{0}, \ldots, j_{d}}\left(A_{j_{0}}^{(0)} \otimes \cdots \otimes A_{j_{d}}^{(d)}\right) \tag{3.20}
\end{equation*}
$$

where $\mathbf{h}$ is a tensor in $\mathcal{H}$-Tucker format with hierarchical ranks $s_{t}$. The application of such a matrix to the tensor $\mathbf{x}$ results in a tensor of ranks $k_{t} \leq r_{t} s_{t}, t \in T$ [27].
A matrix of "Laplacian" structure, such as the second term in (3.10), can be represented as in (3.20) with hierarchical ranks $s_{t}=2$ [29]. Therefore, we can apply B to $\mathbf{x}$ in such a way that the hierarchical ranks increase only by a factor of 3 . Incidentally, if the coefficient $q(t, \xi)$ is independent of $\xi$, i.e., $q(t, \xi)=q(t)$, the matrices $\widehat{M}^{(\mu)}$ and $M^{(\mu)}$ are equal for all $\mu$. It follows that B as a whole has "Laplacian" structure, and can thus be represented with ranks $s_{t}=2$.
It is an important observation that the parabolic BPX preconditioners can be written in the form (3.20). To show this, in the following we write $L_{0}:=K$ and $L_{\mu}:=L$ for $\mu=1, \ldots, d$. We focus on M , as it is the more involved case. All considerations will apply similarly to N . For $s \in \mathbb{R}$, let us first define

$$
\begin{equation*}
\mathrm{M}_{s}=\sum_{\ell_{0}=0}^{L_{0}} \cdots \sum_{\ell_{d}=0}^{L_{d}}\left(\mathbf{g}_{\ell_{0}, \ldots, \ell_{d}}\right)^{s}\left(P_{\ell_{0}}^{(0)} \otimes \cdots \otimes P_{\ell_{d}}^{(d)}\right) \tag{3.21}
\end{equation*}
$$

where the tensor $\mathbf{g}$ is given by

$$
\begin{equation*}
\mathbf{g}_{\ell_{0}, \ldots, \ell_{d}}:=\left(\sum_{\mu=1}^{d} 2^{2 \ell_{\mu}}\right)+2^{2 \ell_{0}}\left(\sum_{\mu=1}^{d} 2^{2 \ell_{\mu}}\right)^{-1} \tag{3.22}
\end{equation*}
$$

and the projection matrix $P_{\ell}^{(\mu)}$ is defined for each $\mu=0, \ldots, d, \ell=1, \ldots, L_{\mu}$ by (we omit the superscripts $(\cdot)^{(\mu)}$ on matrices in the right hand side)

$$
\begin{align*}
& P_{\ell}^{(\mu)}:=M_{L_{\mu}}\left(S_{L_{\mu} \backslash \ell} M_{\ell}^{-1} S_{\ell \downarrow L_{\mu}}-S_{L_{\mu} \backslash \ell-1} M_{\ell-1}^{-1} S_{\ell-1 \downarrow L_{\mu}}\right) M_{L_{\mu}},  \tag{3.23}\\
& P_{0}^{(\mu)}:=M_{L_{\mu}}\left(S_{L_{\mu} \backslash 0} M_{0}^{-1} S_{0 \downarrow L_{\mu}}\right) M_{L_{\mu}} . \tag{3.24}
\end{align*}
$$

Here, $S_{\ell_{2} \backslash \ell_{1}}^{(\mu)}$ denotes the prolongation from level $\ell_{1}$ to level $\ell_{2}$ if $\ell_{1} \leq \ell_{2}$, and $S_{\ell_{1} \downarrow \ell_{2}}^{(\mu)}:=\left(S_{\ell_{2} \backslash \ell_{1}}^{(\mu)}\right)^{\top}$ the reverse operation of restriction, and, furthermore,

$$
\begin{equation*}
M_{\ell}^{(\mu)}:=S_{\ell \not L_{\mu}}^{(\mu)} M^{(\mu)} S_{L_{\mu} \backslash \ell}^{(\mu)} \quad \mu=0, \ldots, d, \quad \ell=0,1, \ldots, L_{\mu}, \tag{3.25}
\end{equation*}
$$

where $M^{(\mu)}$ stands for $M_{e}^{(0)}$ in the case $\mu=0$. Note that, by definition, for each $\mu$, the symmetric matrices $P_{\ell}^{(\mu)}$ satisfy

$$
\begin{equation*}
P_{\ell^{\prime}}^{(\mu)}\left(M_{L_{\mu}}^{(\mu)}\right)^{-1} P_{\ell}^{(\mu)}=\delta_{\ell^{\prime} \ell} P_{\ell}^{(\mu)} \quad \text { for all } \quad \ell^{\prime}, \ell=0, \ldots, L_{\mu} \tag{3.26}
\end{equation*}
$$

and

$$
\begin{equation*}
\sum_{\ell=0}^{L_{\mu}} P_{\ell}^{(\mu)}=M_{L_{\mu}}^{(\mu)} \tag{3.27}
\end{equation*}
$$

Consequently, owing to (2.23) we have $\mathrm{M}_{ \pm}=\mathrm{M}_{ \pm 1}$ and, moreover, the space-time mass matrix $\mathrm{M}_{0}$ defined in (3.21) coincides with $\mathrm{M}_{s}$ for $s=0$.
The block diagonal matrix $\mathrm{N}_{-}$is defined analogously, where the first block has the form

$$
\begin{equation*}
\sum_{\ell_{1}=0}^{L_{1}} \cdots \sum_{\ell_{d}=0}^{L_{d}}\left(\sum_{\mu=1}^{d} 2^{2 \ell_{\mu}}\right) \operatorname{Id} \otimes P_{\ell_{1}}^{(1)} \otimes \cdots \otimes P_{\ell_{d}}^{(d)} \tag{3.28}
\end{equation*}
$$

and the second block is given by $\mathrm{N}_{0,2}^{-1}$.
Recall that we need to compute the action of $\mathrm{M}_{+}^{-1}=\mathrm{M}_{0}^{-1} \mathrm{M}_{-} \mathrm{M}_{0}^{-1}$ and $\mathrm{N}_{+}^{-1}=$ $\mathrm{N}_{0}^{-1} \mathrm{~N}_{-} \mathrm{N}_{0}^{-1}$ for preconditioning. In order to approximate the action of $\mathrm{M}_{-}$, we substitute the entrywise reciprocal tensor $\mathbf{h}$ of $\mathbf{g}$ by a tensor $\widehat{\mathbf{h}}$ in $\mathcal{H}$-Tucker format, i.e., $\widehat{\mathbf{h}}_{\ell_{0}, \ldots, \ell_{d}} \approx \mathbf{h}_{\ell_{0}, \ldots, \ell_{d}}:=\mathbf{g}_{\ell_{0}, \ldots, \ell_{d}}^{-1}$. In our implementation we obtain $\widehat{\mathbf{h}}$ by truncating the full tensor $\mathbf{h}$ of size $(K+1)(L+1)^{d}$. This results in the operators $\widehat{\mathrm{M}}_{-}$and $\widehat{\mathrm{N}}_{-}$. Let us write $\widehat{\mathrm{M}}_{+}:=\mathrm{M}_{0}^{-1} \widehat{\mathrm{M}}_{-} \mathrm{M}_{0}^{-1}$, and analogously for $\widehat{\mathrm{N}}_{+}$. As in the case of B , the hierarchical ranks $s_{t}$ of $\widehat{\mathbf{h}}$ determine the computational cost of applying these matrices to $\mathbf{x}$ in $\mathcal{H}$-Tucker format. Using the variational characterization of the singular values (see (3.13)-(3.14) in [1]), it is easy to check that the operators $\widetilde{B}:=N_{+}^{-\top / 2} \mathrm{BM}_{+}^{-1 / 2}$ and $\widehat{B}:=\widehat{\mathrm{N}}_{+}^{-\top / 2} B \widehat{M}_{+}^{-1 / 2}$ satisfy

$$
\begin{equation*}
\kappa_{2}\left(\widehat{\mathrm{~B}}^{\top} \widehat{\mathrm{B}}\right) \leq \kappa_{2}\left(\widehat{\mathrm{M}}_{+}^{-\top / 2} \mathrm{M}_{+} \widehat{\mathrm{M}}_{+}^{-1 / 2}\right) \kappa_{2}\left(\widehat{\mathrm{~N}}_{+}^{-\top / 2} \mathrm{~N}_{+} \widehat{\mathrm{N}}_{+}^{-1 / 2}\right) \kappa_{2}\left(\widetilde{\mathrm{~B}}^{\top} \widetilde{\mathrm{B}}\right) \tag{3.29}
\end{equation*}
$$

Recall from (2.15) that $\kappa_{2}\left(\widetilde{\mathrm{~B}}^{\top} \widetilde{\mathrm{B}}\right)$ is bounded independently of $K$ and $L$, hence it is enough to choose $\widehat{\mathrm{M}}_{+}^{-1}, \widehat{\mathrm{~N}}_{+}^{-1}$ to be good approximations of $\mathrm{M}_{+}^{-1}, \mathrm{~N}_{+}^{-1}$, respectively. One can show that

$$
\begin{equation*}
J(\mathbf{h}, \widehat{\mathbf{h}}):=\frac{\max _{\ell_{0}, \ldots, \ell_{d}} \mathbf{h}_{\ell_{0}, \ldots, \ell_{d}} / \widehat{\mathbf{h}}_{\ell_{0}, \ldots, \ell_{d}}}{\min _{\ell_{0}, \ldots, \ell_{d}} \mathbf{h}_{\ell_{0}, \ldots, \ell_{d}} / \widehat{\mathbf{h}}_{\ell_{0}, \ldots, \ell_{d}}}=\kappa_{2}\left(\widehat{\mathrm{M}}_{+}^{-\top / 2} \mathrm{M}_{+} \widehat{\mathrm{M}}_{+}^{-1 / 2}\right), \tag{3.30}
\end{equation*}
$$

thus we need to find a tensor $\widehat{\mathbf{h}}$ in $\mathcal{H}$-Tucker format of small hierarchical ranks $s_{t}$ such that $J(\mathbf{h}, \widehat{\mathbf{h}})$ is small. Unfortunately, we only have access to the quasibest approximation in Euclidean norm, i.e., $\|\mathbf{h}-\widehat{\mathbf{h}}\|$, which typically leads to large relative error in entries that are small in absolute value. This is due to the fact that the relative difference in absolute value of the entries is quite pronounced, and grows exponentially in $L$. Table 3.1 displays the smallest ranks $s_{t}$ of $\widehat{\mathbf{h}}$ at which we achieved $J(\mathbf{h}, \widehat{\mathbf{h}}) \leq 100$, for different choices of $d$ and $L$. Note that these are not the minimal ranks needed, but only the minimal ranks available to us while using the quasi-best approximation in the norm $\|\cdot\|$. The ranks go up to 13 , which is unacceptably large for our purpose.

|  |  |  |  | $L$ |  |  |  |  |  |  |  |
| ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: |
|  |  | 1 | 2 | 3 | 4 | 5 | 6 | 7 | 8 | 9 | 10 |
| $d$ | 1 | 1 | 1 | 2 | 3 | 5 | 5 | 5 | 5 | 5 | 5 |
|  | 2 | 1 | 3 | 3 | 6 | 7 | 7 | 7 | 7 | 8 | 8 |
|  | 3 | 1 | 2 | 6 | 7 | 9 | 8 | 8 | 8 | 9 | 10 |
|  | 4 | 2 | 2 | 5 | 8 | 9 | 8 | 10 | 10 | 10 | 11 |
|  | 5 | 1 | 2 | 7 | 8 | 8 | 10 | 9 | 11 | 12 | 13 |

Table 3.1. Hierarchical rank of $\mathcal{H}$-Tucker approximation $\widehat{\mathrm{M}}_{-}$to the parabolic BPX preconditioner $\mathrm{M}_{-}$found to yield $J(\mathbf{h}, \widehat{\mathbf{h}}) \leq$ 100, for discretization levels $L$ and spatial dimensions $d$.

|  | $L$ |  |  |  |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | 1 | 2 | 3 | 4 | 5 | 6 | 7 | 8 | 9 | 10 |
| 1 | 1 | 1 | 2 | 2 | 3 | 3 | 3 | 3 | 3 | 3 |
| 2 | 1 | 1 | 2 | 3 | 3 | 3 | 3 | 3 | 3 | 4 |
| d 3 | 1 | 1 | 2 | 3 | 3 | 3 | 3 | 3 | 4 | 4 |
| 4 | 1 | 1 | 2 | 3 | 3 | 3 | 4 | 4 | 4 | 4 |
| 5 | 1 | 1 | 2 | 3 | 3 | 3 | 4 | 4 | 4 | 4 |

Table 3.2. Hierarchical rank of $\mathcal{H}$-Tucker approximation $\widehat{\mathrm{M}}_{-1 / n}$ to the parabolic BPX preconditioner $\mathrm{M}_{-1 / n}$ found to yield $J\left(\mathbf{h}^{1 / n}, \widehat{\mathbf{h}^{1 / n}}\right)^{n} \leq 100$, for discretization levels $L$ and spatial dimensions $d$ with $n=4$.

In order to alleviate this problem we resort to the approximation

$$
\begin{equation*}
\mathrm{M}_{-}=\mathrm{M}_{0}\left(\mathrm{M}_{0}^{-1} \mathrm{M}_{-1 / n}\right)^{n} \approx \mathrm{M}_{0}\left(\mathcal{T} \circ \mathrm{M}_{0}^{-1} \mathrm{M}_{-1 / n}\right)^{n} \approx \mathrm{M}_{0}\left(\mathcal{T} \circ \mathrm{M}_{0}^{-1} \widehat{\mathrm{M}}_{-1 / n}\right)^{n} \tag{3.31}
\end{equation*}
$$

where $\mathcal{T}$ represents the truncation to lower hierarchical ranks and $\widehat{\mathrm{M}}_{-1 / n}$ results from replacing the full tensor $\mathbf{h}^{1 / n}=\mathbf{g}^{-1 / n}$ (entrywise power) by a suitable $\mathcal{H}$ Tucker approximation $\widehat{\mathbf{h}^{1 / n}}$. The matrix $\widehat{\mathrm{M}}_{-1 / n}$ thus depends on the $\mathcal{H}$-Tucker approximation of $\mathbf{h}^{1 / n}$, which now has a less pronounced relative difference in entry sizes. Indeed, Table 3.2 shows the hierarchical ranks as described for Table 3.1, for the condition $J\left(\mathbf{h}^{1 / n}, \widehat{\mathbf{h}^{1 / n}}\right)^{n} \leq 100$, for $n=4$. The hierarchical ranks $s_{t}$ are now significantly smaller. Since these enter as $n\left(s_{t}\right)^{4}$ into the computation cost, this is a dramatic improvement. Numerical experiments in Section 4 show that this approach is indeed profitable.
3.4. The preconditioned conjugate gradient method. Low-rank tensor variants of classical iterative methods for linear systems have been recently proposed in $[3,24,26]$. A classical iterative method (e.g. Richardson, conjugate gradient) is formulated with low rank tensors as its iterates. Due to repeated addition and application of matrices to the iterates, the ranks will grow rapidly throughout the iteration. This growth is limited by truncation $\mathcal{T}: \mathbf{u} \mapsto \widetilde{\mathbf{u}}$ of the iterates with certain relative accuracy rel_eps and maximal truncation rank max_rank.
We apply the conjugate gradient (CG) method to $B^{\top} N^{-1} B u=B^{\top} N^{-1} \mathbf{b}$ with the preconditioner $\mathrm{M}^{-1}$. We reformulate the $C G$ method for the preconditioned normal equations to resemble the CGNR method [33]. The matrices $\mathrm{N}^{-1}, \mathrm{M}^{-1}$ in Algorithm 3.4 denote the (approximate) preconditioners described in Section
3.4. We monitor the following residuals of the $k$-th iterate: the $C G$ residual $\left\|\mathrm{B}^{\top} \mathrm{N}^{-1} \mathrm{~B} \mathbf{u}_{k}-\mathrm{B}^{\top} \mathrm{N}^{-1} \mathbf{b}\right\|_{\mathrm{M}^{-1}}$ and the least squares residual $\left\|\mathrm{B} \mathbf{u}_{k}-\mathbf{b}\right\|_{\mathrm{N}^{-1}}$.
Typically, the hierarchical ranks grow in the transient phase of the iteration, and decrease as the iterates approach the least squares minimizer. It has been observed in [26] that a good preconditioner is essential for such low rank variants of classical iterative solvers, in order to keep the ranks of the iterates moderate without significant loss of accuracy.

```
Algorithm 1 Variant of low rank CGNR method in \(\mathcal{H}\)-Tucker format
Input: Functions applying \(\mathcal{T} \circ \mathrm{B}, \mathcal{T} \circ \mathrm{M}^{-1}, \mathcal{T} \circ \mathrm{~N}^{-1}\) to a tensor in \(\mathcal{H}\)-Tucker format,
    right-hand side \(\mathbf{b}\) in \(\mathcal{H}\)-Tucker format. Truncation operator \(\mathcal{T}\) with rel. accuracy \(\epsilon_{\text {rel }}\).
Output: Tensor \(\mathbf{u}\) fulfilling \(\|\mathrm{Bu}-\mathbf{b}\|_{\mathrm{N}-1} \leq\) tol.
    \(\mathbf{u}_{0}=0, \mathbf{r}_{0}=\mathbf{b}, \mathbf{s}_{0}=\mathrm{B}^{\top} \mathbf{N}^{-1} \mathbf{r}_{0}, \mathbf{p}_{0}=\mathbf{s}_{0}, \gamma_{0}=\left\langle\mathrm{B} \mathbf{p}_{0}, \mathrm{~B} \mathbf{p}_{0}\right\rangle_{\mathrm{N}^{-1}}, k=0\)
    while \(\left\|\mathbf{r}_{k}\right\|_{N^{-1}}>\) tol do
        \(\alpha_{k}=\left\langle\mathbf{s}_{k}, \mathbf{p}_{k}\right\rangle / \gamma_{k}\)
        \(\mathbf{u}_{k+1}=\mathbf{u}_{k}+\alpha_{k} \mathbf{p}_{k} \quad \mathbf{u}_{k+1} \leftarrow \mathcal{T}\left(\mathbf{u}_{k+1}\right)\)
        \(\mathbf{r}_{k+1}=\mathbf{b}-\mathrm{Bu}_{k+1} \quad \mathbf{r}_{k+1} \leftarrow \mathcal{T}\left(\mathbf{r}_{k+1}\right)\)
        \(\mathbf{s}_{k+1}=\mathrm{B}^{\top} \mathbf{N}^{-1} \mathbf{r}_{k+1}\)
        \(\mathbf{z}_{k+1}=\mathbf{M}^{-1} \mathbf{s}_{k+1}\)
        \(\beta_{k+1}=-\left\langle\mathbf{B z}_{k+1}, \mathbf{B p}_{k}\right\rangle_{\mathbf{N}^{-1}} / \gamma_{k}\)
        \(\mathbf{p}_{k+1}=\mathbf{z}_{k+1}+\beta_{k+1} \mathbf{p}_{k} \quad \mathbf{p}_{k+1} \leftarrow \mathcal{T}\left(\mathbf{p}_{k+1}\right)\)
        \(\gamma_{k+1}=\left\langle\mathbf{B p}_{k+1}, \mathbf{B} \mathbf{p}_{k+1}\right\rangle_{N^{-1}}\)
        \(k=k+1\)
    end while
    \(\mathbf{u}=\mathbf{u}_{k}\)
```


## 4. Numerical experiments

In our numerical experiments, we set $K=8$ and $L=8$ for the discretizaion levels in time and space, resp., cf. Section 2.4, unless specified otherwise. Hence, the trial space is spanned by $513 \times 511^{d}$ functions which are tensor products of standard univariate hat functions. In addition to the isotropic case $q: J \times D \rightarrow \mathbb{R}$ discussed so far, we investigate the anisotropic case $q: J \times D \rightarrow \mathbb{R}^{d \times d}$ with $q$ taking values in the set of diagonal and positive definite matrices such that $0<a_{\min } \leq q_{i i}(t, \xi) \leq$ $a_{\max }<\infty$ still holds, cf. (2.4). We set $T_{\text {final }}=2, D=(-1,1)^{d}$ with $\bar{d}=1, \ldots, 5$, as well as $h \equiv 0$ and

$$
p(t, \xi)=\left(1+\sin \frac{\pi t}{2}\right) \prod_{j=1}^{d}\left(1-\xi_{j}^{2}\right) e^{\xi_{j}}, \quad t \in J, \quad \xi \in D .
$$

In all cases, the preconditioners are approximated with $n=4$ and $s_{t}=3$, see Section 3.3. All numerical experiments have been performed in Matlab version 7.7.0.471, on an Intel Xeon DP X5450 with 3 GHz and $2 \times 6 \mathrm{MB}$ L2 Cache.
4.1. Isotropic diffusion. Throughout this section we set $q \equiv 1$. Figure 4.1 shows the typical convergence behavior of the CG method. Here, $d=2, K=L=5$, max_rank $=50$ and we compare relative truncation accuracy rel_eps $=10^{-2}$ and rel_eps $=10^{-6}$. Evidently, for rel_eps $=10^{-2}$, the hierarchical ranks remain very low (below 5), and the truncation error does not allow the CG residual to decrease below $10^{-4}$. Therefore, the least squares residual also stagnates. For rel_eps $=10^{-6}$, the convergence follows the typical pattern: the least squares residual clearly stagnates while the CG residual still decreases; the stagnation is therefore not caused by the truncation (but rather by the insufficiently exact discretization). The ranks (depicted is $\max _{t \in T} r_{t}$ of the iterate $\mathbf{u}_{k}$ ) grow in the transient phase of the CG

| $d$ | 1 | 2 | 3 | 4 | 5 |
| :---: | :--- | :--- | :--- | :--- | :--- |
| full | 2 MB | 1 GB | 510 GB | 254 TB | 127 PB |
| sparse | 32 KB | 152 KB | 524 KB | 1.5 MB | 3.8 MB |

TABLE 4.1. Storage space required for full and sparse grid approximation with $K=L=8$ and $d=1, \ldots, 5$.


Figure 4.1. Convergence plot of CG for $d=2, L=5$. Truncation tolerance: $10^{-2}$ (left), $10^{-6}$ (right).
method but eventually decrease to about 8 . A good preconditioner is therefore required to limit the number of iterations and consequently the intermediate ranks and the computational effort.
Now we set $K=L=8$. Figure 4.2 shows the convergence history of the least squares residual as a function of the iteration count and of the computation time. We use rel_eps $=10^{-5}$ and max_rank $=30$ for $d=1, \ldots, 5$. The CG convergence rate deteriorates with the number of dimensions $d$. Figure 4.3 shows the (wallclock) time until the least squares residual falls below $2 \times 10^{-3}$ and the average time per iteration as a function of the dimension $d=1, \ldots, 5$.
Figure 4.4 shows the approximate storage space in kilobytes needed for the result of the above iteration with different values of rel_eps between $10^{-1}$ and $10^{-10}$, and max_rank $=30$. For each run, the prolongation of the solution $\mathbf{u}$ onto level $K=L=11$, denoted by $\widehat{\mathbf{u}}$, and the least squares residual $\|\mathrm{B} \widehat{\mathbf{u}}-\mathbf{b}\|_{\mathrm{N}^{-1}}$ are computed. The latter is plotted against the storage requirements (computed as $8 *$ ndofs with the htucker toolbox). For comparison purposes, Table 4.1 gives the storage requirements for the full, resp. sparse grid for $K=L=8$. More precisely, these numbers are the dimensions of the subspaces

$$
\begin{equation*}
E_{K} \otimes V_{L} \quad \text { and } \quad \sum_{k+\ell_{1}+\ldots+\ell_{d} \leq L} E_{k} \otimes V_{\ell_{1}}^{(1)} \otimes \cdots \otimes V_{\ell_{d}}^{(d)} \tag{4.1}
\end{equation*}
$$

respectively, times 8 , which is the byte length of the double precision floating point number. The matrix $\mathbf{N}^{-1}$ is approximated using $n=4$ and $s_{t} \equiv 5$, rel_eps $=10^{-8}$ and max_rank $=40$ (see Section 3.3). For each $d=1, \ldots, 5$, the result of the iteration with rel_eps $=10^{-10}$ is truncated to different hierarchical ranks $r_{t}=$ $1, \ldots, 6$ for all $t \in T$. The least squares residual of the prolongation onto level $K=L=11$ is again plotted against the storage requirements. Comparison of the two plots suggests a quasi-optimality property: the solution produced by the CG method is close to the exact solution truncated to the same hierarchical ranks.
4.2. Anisotropic diffusion. Let $\gamma_{1}:=1, \gamma_{2}:=10^{-1}, \gamma_{3}:=10^{-2}$, etc. Let $q: J \times D \rightarrow \mathbb{R}^{d \times d}$ be the constant function equal to the diagonal matrix $q \equiv$ $\operatorname{diag}\left(\gamma_{1}, \ldots, \gamma_{d}\right)$. Consider now (2.1)-(2.3) with $q$ so defined. In order to take the
anisotropy into account, we modify $\mathbf{g}_{\ell_{0}, \ldots, \ell_{d}}$ in the definition (3.21) of the parabolic BPX preconditioner as follows,

$$
\begin{equation*}
\mathbf{g}_{\ell_{0}, \ldots, \ell_{d}}=\left(\sum_{\mu=1}^{d} \gamma_{\mu} 2^{2 \ell_{\mu}}\right)+2^{2 \ell_{0}}\left(\sum_{\mu=1}^{d} \gamma_{\mu} 2^{2 \ell_{\mu}}\right)^{-1} \tag{4.2}
\end{equation*}
$$

and similarly for N . Note that this implies a modification of the norms that are used to measure the residuals. The following experiments are identical to the ones in the last section, and the conclusions for Figure 4.6-4.7 correspond to those of Figure 4.3-4.4. The convergence rate, however, now deteriorates only moderately in dependence on the dimension $d$, as opposed to the isotropic case. In fact, we observe close to linear scaling in the dimension $d=1, \ldots, 5$ for the computation time (total, as well as per iteration).

## 5. Conclusions

This paper merges a priori stable minimal residual Petrov-Galerkin space-time discretizations and adaptive low rank approximation of high-dimensional tensors in the hierarchical Tucker format for the solution of parabolic evolution equations.
The minimal residual Petrov-Galerkin discretization of the model diffusion equation yields a linear system of the form $\mathrm{B}^{\top} \mathrm{N}^{-1} \mathrm{Bu}=\mathrm{B}^{\top} \mathrm{N}^{-1} \mathbf{b}$ along with a preconditioner $M$. The matrices $N$ and $M$ are based on an extension of the BPX preconditioner into the space-time setting and render the linear system well-conditioned uniformly in the discretization level. This is crucial for the performance of the low rank variant of the conjugate gradient method: a well-conditioned system is required to maintain moderate ranks of the iterates, and consequently, low requirements on both, storage and computational time. The tensor structure of the system matrix (in the case of separable input data) and of the "parabolic BPX preconditioners" proposed here is very convenient: they can be approximated by means of the hierarchical Tucker format and applied to tensors in this format, in particular exploiting parallel architectures. Our numerical experiments demonstrate the potential of this combined approach. In particular, for the anisotropic diffusion, where the conductivity coefficient along the dimension $\mu$ decays exponentially in $\mu$, linear scaling of the computational time in the number of spatial dimensions $d=1, \ldots, 5$ was achieved.
An interesting pending question is the approximation of the tensor $\mathbf{h}$, which is the entrywise inverse of the tensor $\mathbf{g}$ defined in (3.22), in the hierarchical Tucker format. This is essential for the BPX-based preconditioner in higher dimensions: since the size of the full tensor $\mathbf{h}$ scales like $(L+1)^{d}$, where $L$ is the level of the spatial discretization, it quickly becomes prohibitive to assemble it before obtaining a low rank approximation. This currently sets the limit in our computations to essentially $d=5$. A possible solution may be the adaptive cross approximation proposed in [4, 30].
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Figure 4.2. Convergence curves of the CG method for different dimensions $d$. Isotropic case of Section 4.1.


Figure 4.3. Timings for the CG method until a least squares residual of $2 \times 10^{-3}$ is reached. Isotropic case of Section 4.1.



Figure 4.4. Error against storage cost: result of the CG method for different rel_eps (left) and of the best approximant truncated to ranks $1, \ldots, 6$ (right). Isotropic case of Section 4.1.


Figure 4.5. Convergence curves of the CG method for different dimensions $d$. Anisotropic case of Section 4.2.


Figure 4.6. Timings for the CG method until a least squares residual of $2 \times 10^{-3}$ is reached. Anisotropic case of Section 4.2.



Figure 4.7. Error against storage cost: result of the CG method for different rel_eps (left) and of the best approximant truncated to ranks $1, \ldots, 6$ (right). Anisotropic case of Section 4.2.

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