



Fast evaluation of boundary integral operators arising from an eddy current problem

S. Börm^{a,*}, J. Ostrowski^{b,1}

^a *Max-Planck-Institut für Mathematik in den Naturwissenschaften, Inselstraße 22–26, 04103 Leipzig, Germany*

^b *Sonderforschungsbereich 382, Transferbereich 26, Universität Tübingen*

Received 10 April 2003; received in revised form 8 August 2003; accepted 10 August 2003

Abstract

This paper deals with the \mathcal{H}^2 -matrix approximation of matrices that arise from a Galerkin boundary element (BEM) discretization in the context of the **E**-based eddy current model. The BEM operators are dense, thus need to be compressed. They are of complicated structure, i.e., some kernels and basis functions are vector valued, and test and basis functions are not always identical. The \mathcal{H}^2 -matrix approximation technique is applied to the kernels of the four different relevant boundary integral operators. Numerical experiments demonstrate the significant acceleration of an iterative solution procedure by means of matrix compression.

© 2003 Elsevier B.V. All rights reserved.

Keywords: Computational electromagnetism; \mathcal{H}^2 -matrix approximation by interpolation; Boundary integral operators; Eddy current problems

1. Introduction

We consider the simulation of the induction heating process. In this process, a slowly rotating conduction workpiece is exposed to an oscillating electro-magnetic field generated by applying an alternating current to an *inductor*, usually a coil (cf. Fig. 1). The field penetrates the workpiece and creates *eddy currents* due to Faraday's law. Ohmic losses of these currents heat the workpiece.

At the relevant frequency range of 10–40 kHz and in the presence of high conductivities, we can simplify Maxwell's equations governing the electro-magnetic processes to get the *eddy current model*, which neglects the displacement current [2,6].

In the frequency domain, i.e., for time-harmonic excitations with a constant angular frequency $\omega \in \mathbb{R}_{>0}$, it takes the form

* Corresponding author.

E-mail address: sbo@mis.mpg.de (S. Börm).

¹ This work was supported by Deutsche Forschungsgemeinschaft.

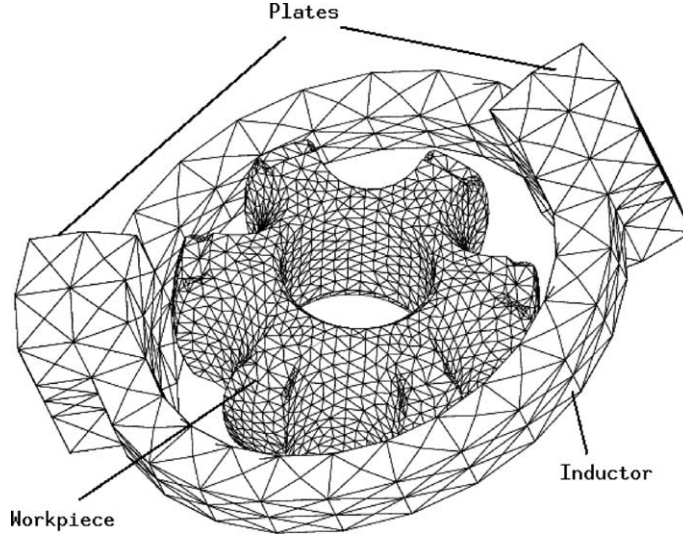


Fig. 1. Typical setting for induction heating: Inductor, workpiece and two plates.

$$\operatorname{div} \mathbf{E} = 0, \quad \text{in } \Omega^+, \quad (1)$$

$$\operatorname{curl} \frac{1}{\mu} \operatorname{curl} \mathbf{E} = -i\omega(\sigma \mathbf{E} + \mathbf{j}_0) \quad \text{in } \mathbb{R}^3, \quad (2)$$

$$[\mathbf{n} \times \mathbf{E}] = \left[\mathbf{n} \times \frac{1}{\mu} \operatorname{curl} \mathbf{E} \right] = 0 \quad \text{in } \partial\Omega^-, \quad (3)$$

$$\mathbf{E}(\mathbf{x}) = \mathcal{O}(|\mathbf{x}|^{-2}), \quad \operatorname{curl} \mathbf{E}(\mathbf{x}) = \mathcal{O}(|\mathbf{x}|^{-2}) \quad \text{for } |\mathbf{x}| \rightarrow \infty, \quad (4)$$

if a formulation based on the electric field \mathbf{E} is used. Here μ is the magnetic permeability, σ is the conductivity and \mathbf{n} is the outer normal. The interior of the items, i.e., workpiece, inductor, and plates, is denoted by Ω^- , the exterior vacuum is denoted by $\Omega^+ = \mathbb{R}^3 \setminus \Omega^-$, and \mathbf{j}_0 is the exciting current density in the inductor. The jump conditions (3), which are chosen at the interface of the items and the vacuum, are the transmission conditions of normal and tangential component of the electric field.

The domain Ω^- is approximated by a tetrahedral triangulation created from CAD data files. This triangulation induces a surface mesh of the boundary consisting of flat open triangles.

Finding a viable numerical scheme for solving the eddy current model is not an easy task. One important reason is that one has to cope with the unbounded exterior of a rotating workpiece with general, genuinely three-dimensional geometry.

An approach is introduced in [12], where the author presents a FEM/BEM-coupled scheme based on edge elements. The FEM part is used in Ω^- and the BEM part, which is needed for the exterior vacuum Ω^+ , consists of boundary integrals over $\partial\Omega^-$. So this scheme uses only elements on $\partial\Omega^-$ and inside Ω^- and can easily be applied to a moving Lagrangian mesh. The BEM part results from applying the Neumann trace and the Dirichlet trace to a Stratton–Chu kind of representation formula for the electric field in the exterior vacuum Ω^+ . The FEM part for the interior domain Ω^- is then coupled to this BEM part by using the jump conditions (3) on the boundary.

Here, we are only interested in the BEM operators, therefore we introduce so-called impedance boundary conditions (cf. [17])

$$\mathbf{n} \times (\mathbf{E} \times \mathbf{n}) = (1 - i) \sqrt{\frac{1}{2\sigma\mu\omega}} \operatorname{curl} \mathbf{E} \times \mathbf{n},$$

that allow us to separate the FEM part from the BEM part. Eq. (5) represents the discretized boundary integrals of the BEM part, see also [13]

$$\begin{pmatrix} M_{\mathfrak{R}} & -M_{\mathfrak{I}} & -\mathbf{B}^T & 0 \\ -M_{\mathfrak{I}} & -M_{\mathfrak{R}} & 0 & \mathbf{B}^T \\ -\mathbf{B} & 0 & -\mathbf{Q} & 0 \\ 0 & \mathbf{B} & 0 & \mathbf{Q} \end{pmatrix} \begin{pmatrix} \mathbf{E}_{\mathfrak{R}} \\ \mathbf{E}_{\mathfrak{I}} \\ \varphi_{\mathfrak{R}} \\ \varphi_{\mathfrak{I}} \end{pmatrix} = \text{right-hand side.} \quad (5)$$

The material parameters and the frequency then appear in the operators $M_{\mathfrak{R}}$ and $M_{\mathfrak{I}}$. If the workpiece is not simply connected, i.e. if there are holes in it, one needs additional matrices that are small enough to remain uncompressed, thus are of minor interest here.

The right-hand side of the equation arises from the exciting current in the inductor. The unknowns $\mathbf{E}_{\mathfrak{R}} \in \mathbb{R}^{\mathcal{E}}$ and $\mathbf{E}_{\mathfrak{I}} \in \mathbb{R}^{\mathcal{E}}$ are the real (\mathfrak{R}) and imaginary (\mathfrak{I}) parts of the electric field, discretized by surface edge elements (cf. Section 2.2). The unknowns $\varphi_{\mathfrak{R}} \in \mathbb{R}^{\mathcal{N}}$ and $\varphi_{\mathfrak{I}} \in \mathbb{R}^{\mathcal{N}}$ are the real and imaginary parts of a scaled scalar magnetic potential $\mathbf{grad} \phi = (\mu_0/\mu) \operatorname{curl} \mathbf{E}$, discretized by standard nodal basis functions on the surface.

Here, \mathcal{N} denotes the set of surface nodes, \mathcal{E} denotes that of surface edges, and the set of surface triangles is denoted by \mathcal{T} .

All these functions and the BEM matrices $M_{\mathfrak{R}} \in \mathbb{R}^{\mathcal{E} \times \mathcal{E}}$, $M_{\mathfrak{I}} \in \mathbb{R}^{\mathcal{E} \times \mathcal{E}}$, $\mathbf{Q} \in \mathbb{R}^{\mathcal{N} \times \mathcal{N}}$ and $\mathbf{B} \in \mathbb{R}^{\mathcal{N} \times \mathcal{E}}$ will be defined more precisely in the next section.

A direct solver cannot be used for Eq. (5) because it needs too much storage and is too slow. Instead, a fast iterative solver is applied, where matrix–vector multiplications dominate the total complexity.

The triangulation of the boundary Γ must be fine enough to meet two different demands. First, the geometry of the items must be described in a satisfactory way, and second, the desired precision of the solution must be achieved. For our application, this means that a number $\#\mathcal{T} \geq 10,000$ of surface triangles must be used for typical workpieces. The occurring BEM operators of Eq. (5) are dense. A matrix–vector multiplication for n unknowns needs $\mathcal{O}(n^2)$ operations, and the amount of storage is of the same order. Parts with $\#\mathcal{T} = 10,000$ surface triangles have approximately $\#\mathcal{E} = 15,000$ edges and $\#\mathcal{N} = 5000$ nodes. For the storage requirements of the matrices in (5) this means:

- Storing $M_{\mathfrak{R}}$ requires $15,000^2 \times \text{sizeof}(\text{double}) = 1.67$ Gbytes,
- storing $M_{\mathfrak{I}}$ requires $15,000^2 \times \text{sizeof}(\text{double}) = 1.67$ Gbytes,
- storing \mathbf{B} requires $15,000 \times 5000 \times \text{sizeof}(\text{double}) = 0.56$ Gbytes and
- storing \mathbf{Q} requires $5000^2 \times \text{sizeof}(\text{double}) = 0.18$ Gbytes.

We see that more than 4 Gbytes are needed, an amount of memory beyond the capacity of current desktop computers. Therefore, a compression technique must be applied to the four different boundary integral operators. This can be done by using the \mathcal{H}^2 -matrix approximation [4].

We remark that there is a close relationship of \mathcal{H}^2 -matrices to the panel clustering technique [11] and the fast multipole method for integral operators [8,14]. While the multipole technique applies an expansion specially designed for the kernel function under investigation in order to reach the, in some sense optimal, complexity of $\mathcal{O}(n \log^2 n)$, the \mathcal{H}^2 -matrix technique has a complexity of $\mathcal{O}(n \log^3 n)$ but can be applied to any asymptotically smooth (cf. (14)) kernel function.

The following section sketches several aspects of the implementation of our method and demonstrates its properties by numerical experiments.

2. Boundary element formulation

2.1. Bilinear forms

For our discrete method, we replace the boundary Γ by a polygonal approximation Γ_h given by the triangulation \mathcal{T}

$$\Gamma_h = \bigcup \{\bar{t} : t \in \mathcal{T}\}.$$

The question of how to choose a suitable triangulation, and therefore approximation Γ_h of Γ , is not a subject of this paper.

The matrices $M_{\mathfrak{R}}$, $M_{\mathfrak{S}}$, Q and B occurring in (5) are Galerkin discretizations of boundary integral operators corresponding to the bilinear forms

$$m_{\mathfrak{R}}(\mathbf{U}, \mathbf{E}) := \int_{\Gamma_h} \int_{\Gamma_h} C_m \langle \gamma_D \mathbf{U}(\mathbf{y}), \gamma_D \mathbf{E}(\mathbf{x}) \rangle \, d\mathbf{y} \, d\mathbf{x} + \int_{\Gamma_h} \int_{\Gamma_h} \langle \mathbf{curl}_\Gamma \mathbf{U}(\mathbf{y}), \mathbf{curl}_\Gamma \mathbf{E}(\mathbf{x}) \rangle \Phi(\mathbf{x}, \mathbf{y}) \, d\mathbf{y} \, d\mathbf{x}, \quad (6)$$

$$m_{\mathfrak{S}}(\mathbf{U}, \mathbf{E}) := \int_{\Gamma_h} \int_{\Gamma_h} C_m \langle \gamma_D \mathbf{U}(\mathbf{y}), \gamma_D \mathbf{E}(\mathbf{x}) \rangle \, d\mathbf{y} \, d\mathbf{x}, \quad (7)$$

$$q(\eta, \phi) := \int_{\Gamma_h} \int_{\Gamma_h} \langle \mathbf{curl}_\Gamma \eta(\mathbf{y}), \mathbf{curl}_\Gamma \phi(\mathbf{x}) \rangle \Phi(\mathbf{x}, \mathbf{y}) \, d\mathbf{y} \, d\mathbf{x}, \quad (8)$$

$$\begin{aligned} b(\eta, \mathbf{E}) := & -\frac{1}{2} \int_{\Gamma_h} \langle \mathbf{curl}_\Gamma \eta(\mathbf{x}), \gamma_D \mathbf{E}(\mathbf{x}) \rangle \, d\mathbf{x} + \int_{\Gamma_h} \int_{\Gamma_h} \langle \mathbf{curl}_\Gamma \eta(\mathbf{y}), \gamma_D \mathbf{E}(\mathbf{x}) \rangle \langle \mathbf{grad}_x \Phi(\mathbf{x}, \mathbf{y}), \mathbf{n}(\mathbf{x}) \rangle \, d\mathbf{y} \, d\mathbf{x}, \\ & - \int_{\Gamma_h} \int_{\Gamma_h} \langle \mathbf{curl}_\Gamma \eta(\mathbf{y}), \mathbf{n}(\mathbf{x}) \rangle \langle \mathbf{grad}_x \Phi(\mathbf{x}, \mathbf{y}), \gamma_D \mathbf{E}(\mathbf{x}) \rangle \, d\mathbf{y} \, d\mathbf{x}, \end{aligned} \quad (9)$$

where

$$C_m = \mu_0 \sqrt{\frac{\sigma \omega}{2\mu}},$$

\mathbf{n} is the outer normal of the surface Γ , \mathbf{E} is an electric field and \mathbf{U} a corresponding test function, while ϕ is the scalar magnetic potential mentioned above and η is the corresponding test function.

Φ is the fundamental solution of the Laplace operator in three space dimensions given by

$$\Phi(\mathbf{x}, \mathbf{y}) := \frac{1}{4\pi} \frac{1}{|\mathbf{x} - \mathbf{y}|}, \quad \mathbf{x}, \mathbf{y} \in \mathbb{R}^3, \mathbf{x} \neq \mathbf{y}. \quad (10)$$

γ_D is the tangential trace operator on Γ_h , \mathbf{curl}_Γ and \mathbf{curl}_Γ are the vector-valued and scalar-valued surface curls. These three operators are given by

$$\gamma_D \mathbf{E} := \mathbf{n} \times \left(\lim_{\epsilon \rightarrow +0} \mathbf{E}(\mathbf{x} + \epsilon \mathbf{n}) \times \mathbf{n} \right),$$

$$\mathbf{curl}_\Gamma \phi := \gamma_D(\mathbf{grad} \phi) \times \mathbf{n},$$

$$\mathbf{curl}_\Gamma \mathbf{E} := \langle \mathbf{n}, \mathbf{curl} \mathbf{E} \rangle.$$

2.2. Discretization

We want to use a piecewise polynomial conforming Galerkin approach to discretize the problem, so we need appropriate discrete spaces. This implies that the discretized potentials η and ϕ have to be continuous at the edges of the triangulation, while the discretized vector fields \mathbf{U} and \mathbf{E} have to have a continuous tangential component.

Further details on the discretization and the properties of the boundary integral operators are given in [12].

Let us consider a triangle $T = \overline{\mathbf{ABC}} \in \mathcal{T}$. The surface measure of the triangle $\overline{\mathbf{ABC}}$ is given by $S := \|(\mathbf{B} - \mathbf{A}) \times (\mathbf{C} - \mathbf{A})\|/2$. We suppose that the vertices \mathbf{A} , \mathbf{B} and \mathbf{C} are ordered counter-clockwise, so that the outer normal vector of T is given by

$$\mathbf{n} := \frac{(\mathbf{B} - \mathbf{A}) \times (\mathbf{C} - \mathbf{A})}{2S}.$$

We introduce the vectors

$$\mathbf{f} := \mathbf{B} - \mathbf{C} \quad \text{and} \quad \mathbf{t} := \frac{\mathbf{n} \times \mathbf{f}}{\|\mathbf{n} \times \mathbf{f}\|},$$

(cf. Fig. 2).

The local edge element basis function corresponding to the edge $e := \overline{\mathbf{BC}}$ is given by

$$\mathbf{b}_{T,e}(\mathbf{x}) := \frac{(\mathbf{x} - \mathbf{A}) \times \mathbf{n}}{2S}.$$

For $\mathbf{x} = \mathbf{C} + \alpha(\mathbf{B} - \mathbf{C})$, we have

$$\begin{aligned} \langle \mathbf{b}_{T,e}(\mathbf{x}), \mathbf{f} \rangle &= \frac{\langle (\mathbf{x} - \mathbf{A}) \times \mathbf{n}, \mathbf{f} \rangle}{2S} = \frac{\det(\mathbf{x} - \mathbf{A}, \mathbf{n}, \mathbf{B} - \mathbf{C})}{2S} = \frac{\det(\mathbf{C} - \mathbf{A}, \mathbf{n}, \mathbf{B} - \mathbf{C})}{2S} \\ &= \frac{\det(\mathbf{B} - \mathbf{A}, \mathbf{C} - \mathbf{A}, \mathbf{n})}{2S} = \frac{\langle (\mathbf{B} - \mathbf{A}) \times (\mathbf{C} - \mathbf{A}), \mathbf{n} \rangle}{2S} = \langle \mathbf{n}, \mathbf{n} \rangle = 1, \end{aligned}$$

i.e., the tangential component of $\mathbf{b}_{T,e}$ on the edge e is constant and equal to 1. Similar computations reveal that the tangential component on $\mathbf{b}_{T,e}$ on the other edges $\overline{\mathbf{AB}}$ and $\overline{\mathbf{CA}}$ are constant and equal to zero. This implies that we can build a global edge element basis function \mathbf{b}_e for each edge $e \in \mathcal{E}$ by combining the local edge element basis functions $\mathbf{b}_{T,e}$ corresponding to the triangles touching e .

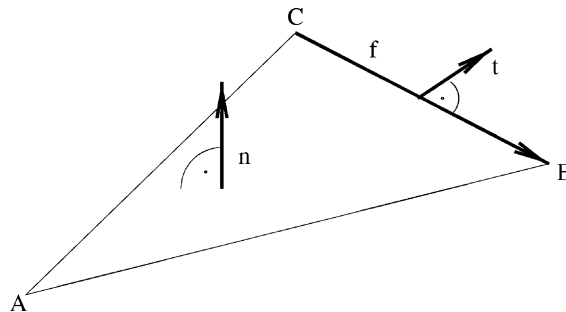


Fig. 2. Setting for the definition of basis functions.

The local nodal basis function corresponding to the vertex $v := \mathbf{A}$ is given by

$$\phi_{T,v}(\mathbf{x}) := -\frac{\|\mathbf{f}\|\langle \mathbf{x} - \mathbf{C}, \mathbf{t} \rangle}{2S}.$$

We have

$$\begin{aligned} \phi_{T,v}(\mathbf{A}) &= -\frac{\|\mathbf{f}\|\langle \mathbf{A} - \mathbf{C}, \mathbf{t} \rangle}{2S} = -\frac{\|\mathbf{f}\|\langle \mathbf{n} \times \mathbf{f}, \mathbf{A} - \mathbf{C} \rangle}{2S\|\mathbf{n} \times \mathbf{f}\|} = -\frac{\det(\mathbf{n}, \mathbf{B} - \mathbf{C}, \mathbf{A} - \mathbf{C})}{2S} \\ &= \frac{\det(\mathbf{n}, \mathbf{B} - \mathbf{A}, \mathbf{C} - \mathbf{A})}{2S} = \frac{\det(\mathbf{B} - \mathbf{A}, \mathbf{C} - \mathbf{A}, \mathbf{n})}{2S} = \langle \mathbf{n}, \mathbf{n} \rangle = 1. \end{aligned}$$

It is obvious that $\phi_{T,v}(\mathbf{B}) = \phi_{T,v}(\mathbf{C}) = 0$ holds. Since $\phi_{T,v}$ is affine, we can build a global nodal basis function ϕ_v for each vertex $v \in \mathcal{N}$ by combining the local nodal basis functions $\phi_{T,v}$ corresponding to the triangles touching v .

Discretizing the bilinear forms (6)–(9) leads to the matrices $M_{\mathfrak{R}} \in \mathbb{R}^{\mathcal{E} \times \mathcal{E}}$, $M_{\mathfrak{T}} \in \mathbb{R}^{\mathcal{E} \times \mathcal{E}}$, $Q \in \mathbb{R}^{\mathcal{N} \times \mathcal{N}}$ and $B \in \mathbb{R}^{\mathcal{N} \times \mathcal{E}}$ defined by

$$\begin{aligned} M_{\mathfrak{R},ij} &:= m_{\mathfrak{R}}(\mathbf{b}_i, \mathbf{b}_j), \quad M_{\mathfrak{T},ij} := m_{\mathfrak{T}}(\mathbf{b}_i, \mathbf{b}_j), \\ Q_{\iota\kappa} &:= q(\psi_{\iota}, \psi_{\kappa}) \quad \text{and} \quad B_{ij} := b(\psi_{\iota}, \mathbf{b}_j), \end{aligned}$$

for $i, j \in \mathcal{E}$ and $\iota, \kappa \in \mathcal{N}$.

2.3. Properties of the discretized matrices

The amount of work involved in assembling the matrices can be significantly reduced by making use of the fact that

$$m_{\mathfrak{R}}(\mathbf{U}, \mathbf{E}) = m_{\mathfrak{T}}(\mathbf{U}, \mathbf{E}) + g(\mathbf{curl}_T \mathbf{U}, \mathbf{curl}_T \mathbf{E}) \quad \text{and} \quad (11)$$

$$q(\eta, \phi) = \sum_{l=1}^3 g((\mathbf{curl}_T \eta)_l, (\mathbf{curl}_T \phi)_l), \quad (12)$$

hold for

$$g(\zeta, \theta) := \int_{\Gamma_h} \int_{\Gamma_h} \zeta(\mathbf{y}) \theta(\mathbf{x}) \Phi(\mathbf{x}, \mathbf{y}) \, d\mathbf{y} \, d\mathbf{x}. \quad (13)$$

An important property of the basis functions ϕ_v for $v \in \mathcal{N}$ and \mathbf{b}_e for $e \in \mathcal{E}$ is that $(\mathbf{curl}_T \phi_v)|_t$ and $(\mathbf{curl}_T \mathbf{b}_e)|_t$ are constant for each triangle $t \in \mathcal{T}$. This implies that we can represent the non-sparse part of $m_{\mathfrak{R}}(\cdot, \cdot)$ and $q(\cdot, \cdot)$ by discretizing $g(\cdot, \cdot)$ by piecewise constant functions $(\chi_t)_{t \in \mathcal{T}}$. This does not hold for the discretization of $b(\cdot, \cdot)$, where we have to use the piecewise linear basis functions \mathbf{b}_e .

All bilinear forms are double integrals with differential operators (namely \mathbf{curl}_T and \mathbf{curl}_T), the trace operator γ_D and the fundamental solution $\Phi(\cdot, \cdot)$ and its derivatives as integrands. The differential operators and the trace operator do not increase the support of the basis function and can therefore be considered “harmless”, which leaves us with the problem of discretizing the integral operators involving the non-local function $\Phi(\cdot, \cdot)$, namely those that correspond to the bilinear forms $g(\cdot, \cdot)$ and $b(\cdot, \cdot)$.

3. Matrix approximation

The integral kernels in (6)–(9) describe long range interactions between the boundary regions at \mathbf{x} and \mathbf{y} . Their strength depends on the inverse of the distance $|\mathbf{x} - \mathbf{y}|$. A common strategy for compression of the

matrices is to approximate the kernels in the so-called *far-field*, i.e., in regions that are far away from each other, whereas one sticks to exact kernels in the *near-field*.

Panel clustering methods are widely used [11,15,16]. They are based on degenerate approximations of the kernel function in the far-field. We construct the approximation by using an interpolation instead of the more traditional Taylor expansion. This leads to fast algorithms, see [7], that can be stated in the context of \mathcal{H}^2 -matrix techniques [4,10].

We replace the singularity function Φ by its Chebyshev interpolation, so we need only pointwise evaluations of Φ instead of the derivatives required by Taylor-based approaches. Since the interpolation of Φ on the entire domain $\Gamma_h \times \Gamma_h$ would not lead to good results due to the singularity at $\mathbf{x} = \mathbf{y}$, we consider sub-domains of the form $\tau \times \sigma$ and apply the interpolation locally. If a constant order m of the interpolation is used, then the resulting approximation is an \mathcal{H}^2 -matrix (cf. [4]). Using this structure, we can perform the matrix–vector multiplication and the discretization of the far-field in $O(nm^3)$ operations, where n is the number of the degrees of freedom. For some kernel functions, the complexity can be improved to $O(n)$ by using a non-constant order [5,15,16]. In our case, a constant order of $m = 2$ or 3 is sufficient.

The \mathcal{H}^2 -matrix approximation method consists of two main parts: the *preparation phase* of the compressed matrix representation, which needs to be performed only once even for several matrix–vector multiplications, and the *matrix–vector multiplication* itself. The preparation phase consists of three parts: we have to find a suitable splitting of $\Gamma_h \times \Gamma_h$ into sub-domains, we have to compute the matrices corresponding to the \mathcal{H}^2 -representation of the far-field blocks and we have to compute the coefficients of the near-field blocks.

3.1. Motivation

3.1.1. Approximation

Due to (11) and (12), we can use sparse matrices in order to reduce the problem of treating $m_{\mathbb{R}}$ and q to that of treating the auxiliary bilinear form g introduced in (13). This can be done by the methods introduced in [4] in the context of hierarchical matrices [3,9].

The kernel function Φ is *asymptotically smooth*, i.e., there are constants $C_{as}, c_0, d \in \mathbb{R}_{>0}$ such that

$$\left| \partial_x^\alpha \partial_y^\beta \Phi(\mathbf{x}, \mathbf{y}) \right| \leq C_{as} c_0^{|\alpha|+|\beta|} (\alpha + \beta)! \|\mathbf{x} - \mathbf{y}\|^{-d-|\alpha|-|\beta|}, \tag{14}$$

holds for all $\mathbf{x}, \mathbf{y} \in \mathbb{R}^3$ with $\mathbf{x} \neq \mathbf{y}$ and all multi-indices $\alpha, \beta \in \mathbb{N}_0^3$. In the case of the Laplace kernel, we have $c_0 = 1$ and $d = 1$.

Let $\tau, \sigma \subseteq \mathbb{R}^d$ be sub-domains of Γ_h such that $\text{dist}(\tau, \sigma) > 0$. We introduce the local bilinear form $g^{\tau, \sigma}$ given by

$$g^{\tau, \sigma}(\zeta, \theta) := \int_\tau \int_\sigma \zeta(\mathbf{y}) \theta(\mathbf{x}) \Phi(\mathbf{x}, \mathbf{y}) \, d\mathbf{y} \, d\mathbf{x}.$$

Eq. (14) implies that the function Φ is smooth on $\tau \times \sigma$, so we can approximate it on this sub-domain by polynomials and use the approximation of the kernel function to define an approximation of the bilinear form $g^{\tau, \sigma}$.

In order to keep our algorithm simple, we will not work with τ and σ directly, but use axis-parallel boxes: Let B^τ and B^σ be minimal d -dimensional axis-parallel boxes satisfying $\tau \subseteq B^\tau$ and $\sigma \subseteq B^\sigma$.

We apply m th order tensor product interpolation operators \mathcal{I}_m^τ and \mathcal{I}_m^σ given by

$$\mathcal{I}_m^\tau[u](\mathbf{x}) = \sum_{\mathbf{v} \in M} u(\mathbf{x}_\mathbf{v}^\tau) \mathcal{L}_\mathbf{v}^\tau(\mathbf{x}) \quad \text{and} \quad \mathcal{I}_m^\sigma[v](\mathbf{y}) = \sum_{\boldsymbol{\mu} \in M} v(\mathbf{x}_\boldsymbol{\mu}^\sigma) \mathcal{L}_\boldsymbol{\mu}^\sigma(\mathbf{y}), \tag{15}$$

where $(\mathbf{x}_v^\tau)_{v \in M}$ and $(\mathbf{x}_\mu^\sigma)_{\mu \in M}$ are interpolation points in B^τ and B^σ and $(\mathcal{L}_v^\tau)_{v \in M}$ and $(\mathcal{L}_\mu^\sigma)_{\mu \in M}$ are the corresponding Lagrange polynomials. M is a suitable index set, usually $\{0, \dots, m\}^3$ (cf. (25)).

We approximate Φ by

$$\tilde{\Phi}^{\tau,\sigma}(\mathbf{x}, \mathbf{y}) := (\mathcal{I}_m^\tau \otimes \mathcal{I}_m^\sigma)[\Phi](\mathbf{x}, \mathbf{y}) = \sum_{v \in M} \sum_{\mu \in M} \Phi(\mathbf{x}_v^\tau, \mathbf{x}_\mu^\sigma) \mathcal{L}_v^\tau(\mathbf{x}) \mathcal{L}_\mu^\sigma(\mathbf{y}). \quad (16)$$

3.1.2. Low-rank representation

Replacing Γ_h by τ and σ and Φ by $\tilde{\Phi}^{\tau,\sigma}$ in (13), we get

$$\begin{aligned} \tilde{\mathbf{g}}^{\tau,\sigma}(\zeta, \theta) &:= \int_\tau \int_\sigma \zeta(\mathbf{y}) \theta(\mathbf{x}) \tilde{\Phi}^{\tau,\sigma}(\mathbf{x}, \mathbf{y}) \, d\mathbf{y} \, d\mathbf{x} \\ &= \sum_{v \in M} \sum_{\mu \in M} \underbrace{\Phi(\mathbf{x}_v^\tau, \mathbf{x}_\mu^\sigma)}_{=: S_{v,\mu}^{\tau,\sigma}} \underbrace{\int_\tau \theta(\mathbf{x}) \mathcal{L}_v^\tau(\mathbf{x}) \, d\mathbf{x}}_{=: V_v^\tau(\theta)} \underbrace{\int_\sigma \zeta(\mathbf{y}) \mathcal{L}_\mu^\sigma(\mathbf{y}) \, d\mathbf{y}}_{=: V_\mu^\sigma(\zeta)} = \sum_{v \in M} \sum_{\mu \in M} S_{v,\mu}^{\tau,\sigma} V_v^\tau(\theta) V_\mu^\sigma(\zeta), \end{aligned} \quad (17)$$

i.e. the bilinear form can be expressed in terms of a, typically small, matrix $S^{\tau,\sigma}$ and a small number of functionals $(V_v^\tau)_{v \in M}$ and $(V_\mu^\sigma)_{\mu \in M}$.

The advantage of the new representation becomes obvious if we discretize the new bilinear form $\tilde{\mathbf{g}}^{\tau,\sigma}$ by introducing

$$\tilde{\mathbf{G}}_{ts}^{\tau,\sigma} := \tilde{\mathbf{g}}^{\tau,\sigma}(\chi_t, \chi_s), \quad V_{tv}^\tau := V_v^\tau(\chi_t) \quad \text{and} \quad V_{s\mu}^\sigma := V_\mu^\sigma(\chi_s),$$

for $t, s \in \mathcal{T}$. The Eq. (17) now takes the form

$$\tilde{\mathbf{G}}^{\tau,\sigma} = \mathbf{V}^\tau \mathbf{S}^{\tau,\sigma} (\mathbf{V}^\sigma)^\top. \quad (18)$$

We set

$$n^\tau := \#\{t \in \mathcal{T} : \tau \cap t \neq \emptyset\}, \quad n^\sigma := \#\{s \in \mathcal{T} : \sigma \cap s \neq \emptyset\} \quad \text{and} \quad k := \#\mathcal{M},$$

and find that storing $\tilde{\mathbf{G}}^{\tau,\sigma}$ as a dense matrix requires $n^\tau n^\sigma$ units of memory, while storing \mathbf{V}^τ , \mathbf{V}^σ and $\mathbf{S}^{\tau,\sigma}$ requires $n^\tau k + n^\sigma k + k^2$ units of memory. Typically k is much smaller than n^τ and n^σ , so the factorized representation is much more efficient.

3.1.3. Precision

We have seen that replacing the kernel function Φ by its interpolant $\tilde{\Phi}^{\tau,\sigma}$ leads to an efficient representation of the discretized matrix $\tilde{\mathbf{G}}^{\tau,\sigma}$. In order to be able to use this representation, we have to ensure that the error introduced by the interpolation can be controlled, i.e., that an estimate of the form

$$|\tilde{\Phi}^{\tau,\sigma}(\mathbf{x}, \mathbf{y}) - \Phi(\mathbf{x}, \mathbf{y})| \leq \epsilon_m,$$

holds for $\mathbf{x} \in \tau$, $\mathbf{y} \in \sigma$, where $\epsilon_m \in \mathbb{R}_{>0}$ depends favorably on the interpolation order m .

If we use Chebyshev interpolation for \mathcal{I}_m^τ and \mathcal{I}_m^σ , we have a constant $C_{\text{in}} \in \mathbb{R}_{>0}$ such that

$$\|\Phi - \mathcal{I}_m^\tau \otimes \mathcal{I}_m^\sigma[\Phi]\|_{\infty, B^\tau \times B^\sigma} \leq C_{\text{in}} \frac{3^{-m}}{(m+1)!} \text{diam}(B^\tau \times B^\sigma)^{m+1} \|D_{xy}^{m+1} \Phi\|_{\infty, B^\tau \times B^\sigma}, \quad (19)$$

holds for all m and all functions $u \in C^{m+1}(B^\tau \times B^\sigma)$ (cf. appendix of [4]), where $\text{diam}(B^\tau \times B^\sigma)$ denotes the Euclidean diameter of the axis-parallel box $B^\tau \times B^\sigma$ and $D^{m+1}\Phi$ is the total derivative of order $m+1$ of Φ .

Combining this estimate with (14), we find

$$\|\Phi - \tilde{\Phi}^{\tau,\sigma}\|_{\infty, B^\tau \times B^\sigma} \leq \frac{C_{\text{in}} C_{\text{as}}}{\text{dist}(B^\tau, B^\sigma)^d} 3^{-m} \left(\frac{c_0 \text{diam}(B^\tau \times B^\sigma)}{\text{dist}(B^\tau, B^\sigma)} \right)^{m+1}.$$

This implies that we have to be able to bound the diameter by the distance in order to reach a uniform bound for the approximation error, i.e., we require the *admissibility condition*

$$\text{diam}(B^\tau \times B^\sigma) \leq \eta \text{dist}(B^\tau, B^\sigma), \quad (20)$$

to hold for a parameter $\eta \in]0, 3/c_0[$ and find

$$\|\Phi - \tilde{\Phi}^{\tau,\sigma}\|_{\infty, B^\tau \times B^\sigma} \leq \epsilon_m := \frac{C_{\text{in}} C_{\text{as}} \eta^{d+1}}{\text{diam}(B^\tau \times B^\sigma)^d} (c_0 \eta / 3)^m, \quad (21)$$

so the interpolant converges exponentially on $B^\tau \times B^\sigma \supseteq \tau \times \sigma$ if the order m is increased.

3.2. Decomposition of $\Gamma_h \times \Gamma_h$

In the previous section, we have seen that we can efficiently approximate the local interpolants $\tilde{\Phi}^{\tau,\sigma}$ if the sub-domains τ, σ of Γ_h satisfy the admissibility condition (20).

3.2.1. Block partition

Obviously, the pair (Γ_h, Γ_h) does not satisfy this condition, so we have to split $\Gamma_h \times \Gamma_h$, the domain of integration of the bilinear form $g(\cdot, \cdot)$, into a collection P of sub-domains that either satisfy this condition or are so small that we can treat them directly without compromising the efficiency.

The family $P \subseteq \{(\tau, \sigma) : \tau, \sigma \subseteq \Gamma_h\}$ has to satisfy the following conditions:

$$\begin{aligned} \bigcup \{\overline{\tau \times \sigma} : (\tau, \sigma) \in P\} &= \Gamma_h \times \Gamma_h, \\ (\tau_1 \times \sigma_1) \cap (\tau_2 \times \sigma_2) \neq \emptyset &\Rightarrow (\tau_1, \sigma_1) = (\tau_2, \sigma_2) \quad \text{for all } (\tau_1, \sigma_1), (\tau_2, \sigma_2) \in P. \end{aligned}$$

According to the admissibility condition, we split P into the set P_{far} of *far-field* blocks and the set P_{near} of *near-field* blocks

$$P_{\text{far}} := \{(\tau, \sigma) \in P : \text{diam}(B^\tau \times B^\sigma) \leq \eta \text{dist}(B^\tau, B^\sigma)\}, \quad P_{\text{near}} := P \setminus P_{\text{far}}.$$

For each $(\tau, \sigma) \in P_{\text{far}}$, we can construct a local approximation $\tilde{g}^{\tau,\sigma}(\cdot, \cdot)$ of the form (13). For the remaining blocks $(\tau, \sigma) \in P_{\text{near}}$, we use the original local bilinear form $g^{\tau,\sigma}(\cdot, \cdot)$. The approximation of the *global* bilinear form $g(\cdot, \cdot)$ is then given as the sum of the local bilinear forms:

$$\tilde{g}(\zeta, \theta) = \sum_{(\tau,\sigma) \in P_{\text{far}}} \tilde{g}^{\tau,\sigma}(\zeta, \theta) + \sum_{(\tau,\sigma) \in P_{\text{near}}} g^{\tau,\sigma}(\zeta, \theta). \quad (22)$$

This corresponds to replacing the kernel function Φ by its piecewise m th order interpolant given by

$$\tilde{\Phi}(\mathbf{x}, \mathbf{y}) := \begin{cases} \tilde{\Phi}^{\tau,\sigma}(\mathbf{x}, \mathbf{y}) & \text{if } (\mathbf{x}, \mathbf{y}) \in \tau \times \sigma \text{ for } (\tau, \sigma) \in P_{\text{far}}, \\ \Phi(\mathbf{x}, \mathbf{y}) & \text{otherwise.} \end{cases}$$

The estimate (21) implies

$$|\Phi(\mathbf{x}, \mathbf{y}) - \tilde{\Phi}(\mathbf{x}, \mathbf{y})| \leq \epsilon_m,$$

for all $\mathbf{x}, \mathbf{y} \in \Gamma_h$, so we find

$$\begin{aligned}
|g(\zeta, \theta) - \tilde{g}(\zeta, \theta)| &= \left| \int_{\Gamma_h} \int_{\Gamma_h} \zeta(\mathbf{y}) \theta(\mathbf{x}) (\Phi(\mathbf{x}, \mathbf{y}) - \tilde{\Phi}(\mathbf{x}, \mathbf{y})) \, d\mathbf{y} \, d\mathbf{x} \right| \\
&\leq \|\Phi - \tilde{\Phi}\|_{\infty, \Gamma_h \times \Gamma_h} \int_{\Gamma_h} |\zeta(\mathbf{y})| \, d\mathbf{y} \int_{\Gamma_h} |\theta(\mathbf{x})| \, d\mathbf{x} \\
&\leq \epsilon_m \|\zeta\|_{L^1} \|\theta\|_{L^1} \leq \epsilon_m |\Gamma_h| \|\zeta\|_{L^2(\Gamma_h)} \|\theta\|_{L^2(\Gamma_h)}.
\end{aligned}$$

The construction of a good partition, i.e., a partition where the number of non-admissible blocks is small, is not trivial. Therefore, we will start by constructing a hierarchy of partitions of Γ_h and then use this hierarchy to create a partition of $\Gamma_h \times \Gamma_h$.

3.2.2. Cluster tree

We will construct the hierarchy of partitions of Γ_h by successively splitting domains. Since a pair (τ, σ) of sub-domains of Γ_h is admissible if the diameter of $B^\tau \times B^\sigma$ is smaller than the distance of these boxes (recall that B^τ and B^σ are the minimal axis-parallel boxes containing τ and σ), a good strategy is to split domains in such a way that the diameters of the newly created sub-domains are decreased as much as possible.

In order to keep the implementation simple, we consider only sub-domains τ that are the union of a set $\hat{\tau} \subseteq \mathcal{T}$ of triangles, and we identify the sub-domain τ with the set $\hat{\tau}$. Then successive splitting of domains, starting with the set \mathcal{T} of all triangles, leads to a tree structure, the *cluster tree*:

Definition 1. A tree \mathcal{C} is called a *cluster tree* for a set Γ if

- the set Γ is the root of \mathcal{C} , i.e., $\text{root}(\mathcal{C}) = \Gamma$, and
- if a node $\tau \in \mathcal{C}$ is not a leaf, then it is the (up to sets of measure zero) disjoint union of its sons, i.e.,

$$\tau = \bigcup \{\tau' : \tau' \in \text{sons}(\tau)\}.$$

Each node $\tau \in \mathcal{C}$ is called a *cluster*.

A cluster tree can be constructed from an arbitrary set of triangles by *binary space partitioning*: We start with the root cluster containing all the triangles, corresponding to the entire domain, split it into two son clusters and repeat the procedure recursively until the clusters contain less than a fixed number $C_{\text{lf}} \geq 1$ of triangles.

The splitting strategy is based on the geometry: We denote the center of each triangle $t \in \mathcal{T}$ by $\mu_t \in \mathbb{R}^3$, choose a suitable coordinate axis and split the set along this axis. This leads to the algorithm in Fig. 3.

```

procedure GeometricBisection( $\tau$ );
begin
  if  $\#\tau \geq C_{\text{lf}}$  then begin
    for  $j := 1$  to 3 do begin
      { find splitting coordinate }
       $s_j := \max\{\mu_{t,j} : t \in \mathcal{T}\}$ ;  $i_j := \min\{\mu_{t,j} : t \in \mathcal{T}\}$ ;  $\delta_j := s_j - i_j$ 
    end;
    Choose  $l \in \{1, 2, 3\}$  such that  $\delta_l = \max\{\delta_j : j \in \{1, \dots, 3\}\}$ ;
    Split  $\tau = \tau_1 \cup \tau_2$  with  $|\#\tau_1 - \#\tau_2| \leq 1$  and  $\mu_{t_1,l} \leq \mu_{t_2,l}$  for  $t_1 \in \tau_1, t_2 \in \tau_2$ ;
    GeometricBisection( $\tau_1$ ); GeometricBisection( $\tau_2$ );
    sons( $\tau$ ) :=  $\{\tau_1, \tau_2\}$ 
  end else sons( $\tau$ ) :=  $\emptyset$ 
end

```

Fig. 3. Cardinality (#) balanced geometric bisection.

Remark 2 (Geometric balancing). If the surface triangulation is not quasi-uniform, e.g., if it is the result of an adaptive refinement strategy, then a different splitting technique than that given in Fig. 3 is to be applied: instead of splitting the cluster τ into two clusters τ_1, τ_2 that are of similar cardinality, we use the middle of the coordinate interval $[i_l, s_l]$ in order to determine which son of τ has to contain which triangles

$$\tau_1 := \{t \in \tau : \mu_{t,l} \leq (s_l + i_l)/2\}, \quad \tau_2 := \{t \in \tau : \mu_{t,l} > (s_l + i_l)/2\},$$

with l still denoting the longest edge.

3.2.3. Construction of a block partition

The definition of the axis-parallel boxes B^τ carries over to clusters;

Definition 3 (Bounding boxes). Let $\tau \in \mathcal{C}$. The minimal axis-parallel box $B^\tau \subseteq \mathbb{R}^3$ satisfying $t \subseteq B$ for all $t \in \tau$ is called the *bounding box* of the cluster τ .

We will use the following simplified admissibility condition:

Remark 4 (Simplified admissibility condition). In some applications, the admissibility condition (20) is replaced by the condition

$$\max\{\text{diam}(B^\tau), \text{diam}(B^\sigma)\} \leq 2\eta' \text{dist}(B^\tau, B^\sigma). \tag{23}$$

By setting $\eta := 2\sqrt{2}\eta'$, we find

$$\text{diam}(B^\tau \times B^\sigma) \leq \sqrt{2} \max\{\text{diam}(B^\tau), \text{diam}(B^\sigma)\} \leq 2\sqrt{2}\eta' \text{dist}(B^\tau, B^\sigma),$$

so the simplified admissibility condition implies the original condition. If $\eta' < 3/(2\sqrt{2}c_0)$, then we get $\eta < 3/c_0$ and therefore exponential convergence.

Based on the criteria (20) or (23) and a cluster tree, we can find a partition of $\Gamma_h \times \Gamma_h$ by calling the algorithm `BlockPartition`($\Gamma_h, \Gamma_h, \emptyset$) of Fig. 4.

Remark 5 (Helmholtz kernel). If we consider the Helmholtz kernel

$$G(x, y) = \frac{\exp(i\kappa\|x - y\|)}{\|x - y\|},$$

the constant c_0 appearing in the asymptotic smoothness estimate (14) will scale linearly in κ , so the speed of the convergence will deteriorate for high wave numbers.

When dealing with a curve in two-dimensional space, a special multipole expansion can be used to achieve fast convergence even for high wave numbers [1], but it is not clear to us whether this approach can be extended to the three-dimensional problem we are considering.

```

procedure BlockPartition( $\tau, \sigma, \text{var } P$ );
begin
  if ( $\tau, \sigma$ ) is admissible then  $P := P \cup \{(\tau, \sigma)\}$ 
  else if sons( $\tau$ ) =  $\emptyset$  or sons( $\sigma$ ) =  $\emptyset$  then  $P := P \cup \{(\tau, \sigma)\}$ 
  else for  $\tau' \in \text{sons}(\tau)$  and  $\sigma' \in \text{sons}(\sigma)$  do BlockPartition( $\tau', \sigma', P$ )
end
    
```

Fig. 4. Construction of a partition of $\Gamma_h \times \Gamma_h$.

3.3. Matrix–vector multiplication

The matrix–vector multiplication $\mathbf{v} = \tilde{\mathbf{G}}\mathbf{u}$ with a vector $\mathbf{u} \in \mathbb{R}^{\mathcal{J}}$ corresponds to the evaluation of Eq. (22). We split this evaluation into four parts:

- **Forward transformation:** We transform \mathbf{u} into coefficients $\mathbf{u}^\sigma := \mathbf{V}^{\sigma^\top} \mathbf{u} \in \mathbb{R}^M$ corresponding to each cluster.
- **Transformed multiplication:** We evaluate the sum $\mathbf{v}^\tau := \sum_{\sigma \in \text{row}(\tau)} \mathbf{S}^{\tau, \sigma} \mathbf{u}^\sigma \in \mathbb{R}^M$ for $\text{row}(\tau) := \{\sigma \in \mathcal{C} : (\tau, \sigma) \in P_{\text{far}}\}$.
- **Backward transformation:** We transform the coefficients \mathbf{v}^τ back into the standard base in order to find $\mathbf{v}_{\text{far}} := \sum_{\tau \in \mathcal{C}} \mathbf{V}^\tau \mathbf{v}^\tau$.
- **Near-field computation:** We conclude the computation by adding the near-field part $\mathbf{v} := \mathbf{v}_{\text{far}} + \sum_{(\tau, \sigma) \in P_{\text{near}}} \mathbf{G}^{\tau, \sigma} \mathbf{u}$.

Due to

$$\tilde{\mathbf{G}}\mathbf{u} = \sum_{(\tau, \sigma) \in P_{\text{near}}} \mathbf{G}^{\tau, \sigma} \mathbf{u} + \sum_{(\tau, \sigma) \in P_{\text{far}}} \mathbf{V}^\tau \mathbf{S}^{\tau, \sigma} \mathbf{V}^{\sigma^\top} \mathbf{u} = \sum_{(\tau, \sigma) \in P_{\text{near}}} \mathbf{G}^{\tau, \sigma} \mathbf{u} + \sum_{(\tau, \sigma) \in P_{\text{far}}} \mathbf{V}^\tau \mathbf{S}^{\tau, \sigma} \mathbf{u}^\sigma = \sum_{(\tau, \sigma) \in P_{\text{near}}} \mathbf{G}^{\tau, \sigma} \mathbf{u} + \sum_{\tau \in \mathcal{C}} \mathbf{V}^\tau \mathbf{v}^\tau = \mathbf{v},$$

this four-step procedure indeed computes the matrix–vector product.

In order to find a fast algorithm for the matrix–vector multiplication, we will now introduce an alternative representation of the matrices \mathbf{V}^τ : Let $\tau \in \mathcal{C}$ be a cluster with $\text{sons}(\tau) \neq \emptyset$. Since we use the same order of interpolation for all clusters, we have

$$\mathcal{L}_v^\tau = \mathcal{I}^{\tau'} [\mathcal{L}_v^\tau] = \sum_{v' \in M} \mathcal{L}_v^\tau(\mathbf{x}_{v'}^{\tau'}) \mathcal{L}_{v'}^{\tau'} = \sum_{v' \in M} \mathbf{T}_{v'v}^{\tau', \tau} \mathcal{L}_{v'}^{\tau'},$$

with *transfer matrices* $\mathbf{T}^{\tau', \tau} \in \mathbb{R}^{M \times M}$ defined by

$$\mathbf{T}_{v'v}^{\tau', \tau} := \mathcal{L}_v^\tau(\mathbf{x}_{v'}^{\tau'}). \quad (24)$$

This alternative representation implies

$$\mathbf{v}_{lv}^\tau = \left(\sum_{\tau' \in \text{sons}(\tau)} \mathbf{V}^{\tau'} \mathbf{T}^{\tau', \tau} \right)_{lv}.$$

Using these equations, we find the recursive procedures for the computation of \mathbf{u}^τ and \mathbf{v}_{far} given in Fig. 5. Combining these procedures, we can derive the fast matrix–vector multiplication algorithm that is given in Fig. 6.

Remark 6 (Storage). The introduction of the transfer matrices $\mathbf{T}^{\tau', \tau}$ leads to a significant reduction in the amount of memory needed to store the \mathcal{H}^2 -approximation of the matrix \mathbf{G} : since we are able to construct \mathbf{V}^τ for all non-leaf clusters τ by using the transfer matrices, we need to store \mathbf{V}^τ only for leaf clusters.

This reduces the amount of storage required to store the \mathcal{H}^2 -matrix approximation to $\mathcal{O}(nm^3)$ (cf. [4]).

3.4. Treatment of $b(\cdot, \cdot)$

Since the bilinear forms $m_{\mathfrak{R}}(\cdot, \cdot)$ and $q(\cdot, \cdot)$ can be expressed by $g(\cdot, \cdot)$, we now have to treat $b(\cdot, \cdot)$ (cf. (9)) in order to be able to compress all matrices occurring in our boundary element formulation.

Since this bilinear form is based on $\mathbf{grad}\Phi$ instead of Φ , we have to find a degenerate approximation of the derivatives of the kernel function. In order to keep our algorithm simple, we use the derivatives of the approximation $\tilde{\Phi}$ of Φ , i.e., replace $\mathbf{grad}\Phi$ by $\mathbf{grad}\tilde{\Phi}$. Please note that $\mathbf{grad}\tilde{\Phi}$ exists almost everywhere, since the local interpolants $\tilde{\Phi}^{\tau, \sigma}$ are polynomials and therefore differentiable.

```

procedure FastForward( $\sigma$ ,  $\mathbf{u}$ , var ( $\mathbf{u}^\sigma$ ) $_{\sigma \in \mathcal{C}}$ );
begin
  if sons( $\sigma$ ) =  $\emptyset$  then  $\mathbf{u}^\sigma := \mathbf{V}^{\sigma \top} \mathbf{u}$ 
  else begin
    for  $\sigma' \in \text{sons}(\sigma)$  do FastForward( $\sigma'$ ,  $\mathbf{u}$ , ( $\mathbf{u}^{\sigma'}$ ));
     $\mathbf{u}^\sigma := \sum_{\sigma' \in \text{sons}(\sigma)} \mathbf{T}^{\sigma', \sigma \top} \mathbf{u}^{\sigma'}$ 
  end
end;

procedure FastBackward( $\tau$ , ( $\mathbf{v}^\tau$ ) $_{\tau \in \mathcal{C}}$ , var  $\mathbf{v}$ );
begin
  if sons( $\tau$ ) =  $\emptyset$  then  $\mathbf{v} := \mathbf{v} + \mathbf{V}^\tau \mathbf{v}^\tau$ 
  else
    for  $\tau' \in \text{sons}(\tau)$  do begin
       $\mathbf{v}^{\tau'} := \mathbf{v}^{\tau'} + \mathbf{T}^{\tau', \tau} \mathbf{v}^\tau$ ; FastBackward( $\tau'$ , ( $\mathbf{v}^\tau$ ),  $\mathbf{v}$ )
    end
  end
end

```

Fig. 5. Fast forward and backward transformations.

```

procedure MatrixVectorMultiplication( $\mathbf{u}$ , var  $\mathbf{v}$ );
begin
  FastForward( $\mathcal{T}$ ,  $\mathbf{u}$ , ( $\mathbf{u}^\sigma$ ));
  for  $\tau \in \mathcal{C}$  do  $\mathbf{v}^\tau := \sum_{\sigma \in \text{row}(\tau)} \mathbf{S}^{\tau, \sigma} \mathbf{u}^\sigma$ ;
  FastBackward( $\mathcal{T}$ , ( $\mathbf{v}^\tau$ ),  $\mathbf{v}_{\text{far}}$ );
   $\mathbf{v} := \mathbf{v}_{\text{far}} + \sum_{(\tau, \sigma) \in P_{\text{near}}} \mathbf{G}^{\tau, \sigma} \mathbf{u}$ 
end

```

Fig. 6. Matrix–vector multiplication.

We ignore the sparse parts of the bilinear form and use the same approach as before on the remainder, i.e., we replace the singularity function Φ by its local approximations $\tilde{\Phi}^{\tau, \sigma}$ for $(\tau, \sigma) \in P_{\text{far}}$. This leads to the following local bilinear forms:

$$\begin{aligned}
\tilde{b}^{\tau, \sigma}(\eta, \mathbf{E}) &= \int_{\tau} \int_{\sigma} \langle \mathbf{curl}_{\Gamma} \eta(\mathbf{y}), \gamma_D \mathbf{E}(\mathbf{x}) \rangle \langle \mathbf{grad}_x \tilde{\Phi}^{\tau, \sigma}(\mathbf{x}, \mathbf{y}), \mathbf{n}(\mathbf{x}) \rangle \, d\mathbf{y} \, d\mathbf{x} \\
&\quad - \int_{\tau} \int_{\sigma} \langle \mathbf{curl}_{\Gamma} \eta(\mathbf{y}), \mathbf{n}(\mathbf{x}) \rangle \langle \mathbf{grad}_x \tilde{\Phi}^{\tau, \sigma}(\mathbf{x}, \mathbf{y}), \gamma_D \mathbf{E}(\mathbf{x}) \rangle \, d\mathbf{y} \, d\mathbf{x} \\
&= \sum_{\nu \in M} \sum_{\mu \in M} \mathbf{S}_{\nu\mu}^{\tau, \sigma} \int_{\Gamma_{\tau}} \int_{\Gamma_{\sigma}} \left(\sum_{\ell=1}^3 (\mathbf{curl}_{\Gamma} \eta)_{\ell}(\mathbf{y}) (\gamma_D \mathbf{E})_{\ell}(\mathbf{x}) \langle \mathbf{grad} \mathcal{L}_{\nu}^{\tau}(\mathbf{x}), \mathbf{n}(\mathbf{x}) \rangle \mathcal{L}_{\mu}^{\sigma}(\mathbf{y}) \right. \\
&\quad \left. - \sum_{\ell=1}^3 (\mathbf{curl}_{\Gamma} \eta)_{\ell}(\mathbf{y}) \mathbf{n}_{\ell}(\mathbf{x}) \langle \mathbf{grad} \mathcal{L}_{\nu}^{\tau}(\mathbf{x}), \gamma_D \mathbf{E}(\mathbf{x}) \rangle \mathcal{L}_{\mu}^{\sigma}(\mathbf{y}) \right) \, d\mathbf{y} \, d\mathbf{x} \\
&= \sum_{\ell=1}^3 \sum_{\nu \in M} \sum_{\mu \in M} \mathbf{S}_{\nu\mu}^{\tau, \sigma} \int_{\sigma} (\mathbf{curl}_{\Gamma} \eta)_{\ell}(\mathbf{y}) \mathcal{L}_{\mu}^{\sigma}(\mathbf{y}) \, d\mathbf{y} \cdot \int_{\tau} (\gamma_D \mathbf{E})_{\ell}(\mathbf{x}) \langle \mathbf{grad} \mathcal{L}_{\nu}^{\tau}(\mathbf{x}), \mathbf{n}(\mathbf{x}) \rangle \\
&\quad - \mathbf{n}_{\ell}(\mathbf{x}) \langle \mathbf{grad} \mathcal{L}_{\nu}^{\tau}(\mathbf{x}), \gamma_D \mathbf{E}(\mathbf{x}) \rangle \, d\mathbf{x}.
\end{aligned}$$

Apart from the summation over ℓ , this representation is similar to that in (18), so we introduce matrices $V^{\tau,\ell} \in \mathbb{R}^{\mathcal{I} \times M}$ and $W^{\sigma,\ell} \in \mathbb{R}^{\mathcal{A} \times M}$ by setting

$$V_{ev}^{\tau,\ell} := \int_{\tau} (\gamma_D \mathbf{b}_e)_\ell(\mathbf{x}) \langle \mathbf{grad} \mathcal{L}_v^\tau(\mathbf{x}), \mathbf{n}(\mathbf{x}) \rangle - \mathbf{n}_\ell(\mathbf{x}) \langle \mathbf{grad} \mathcal{L}_v^\tau(\mathbf{x}), \gamma_D \mathbf{b}_e \rangle \, d\mathbf{x},$$

$$W_{v\mu}^{\sigma,\ell} := \int_{\sigma} (\mathbf{curl}_\Gamma \phi_v)_\ell(\mathbf{y}) \mathcal{L}_\mu^\sigma(\mathbf{y}) \, d\mathbf{y},$$

and find

$$\tilde{\mathbf{B}}^{\tau,\sigma} = \sum_{\ell=1}^3 V^{\tau,\ell} S^{\tau,\sigma} W^{\sigma,\ell\top},$$

where $\tilde{\mathbf{B}}^{\tau,\sigma}$ is the Galerkin discretization of the local bilinear form $\tilde{b}^{\tau,\sigma}(\cdot, \cdot)$ and $S^{\tau,\sigma}$ is the matrix from (17). Using this representation, we can treat the approximation of $b(\cdot, \cdot)$ by exactly the same techniques as that of $g(\cdot, \cdot)$.

4. Implementation

4.1. Interpolation

We use tensor product Chebyshev interpolation, i.e., the interpolation points $(\mathbf{x}_v^\tau)_{v \in M}$ and $(\mathbf{y}_\mu^\sigma)_{\mu \in M}$ in Eq. (16) are the Chebyshev points for the axis-parallel boxes B^τ and B^σ .

The computation of these points is straightforward: The m th order Chebyshev points $(x_i)_{i=0}^m$ for the interval $[-1, 1]$ are given by

$$x_i := \cos\left(\pi \frac{2i+1}{2m+2}\right).$$

For a given interval $[a, b] \subseteq \mathbb{R}$, the transformed Chebyshev points $(x_i^{[a,b]})_{i=0}^m$ are

$$x_i^{[a,b]} := \frac{b+a}{2} + \frac{b-a}{2} x_i,$$

and the corresponding one-dimensional Lagrange polynomials have the form

$$\mathcal{L}_i^{[a,b]}(x) := \prod_{j \neq i} \frac{x - x_j^{[a,b]}}{x_i^{[a,b]} - x_j^{[a,b]}}.$$

The axis-parallel box B^τ can be written as $B^\tau = [a_1, b_1] \times [a_2, b_2] \times [a_3, b_3]$, so the tensor product Chebyshev points are given by

$$\mathbf{x}_v^\tau := \left(x_{v_1}^{[a_1, b_1]}, x_{v_2}^{[a_2, b_2]}, x_{v_3}^{[a_3, b_3]}\right), \quad (25)$$

for $\mathbf{v} = (v_1, v_2, v_3) \in M := \{\mathbf{v} \in \mathbb{N}_0^3 : \|\mathbf{v}\|_\infty \leq m\}$. The corresponding Lagrange polynomials are given by

$$\mathcal{L}_v^\tau(\mathbf{x}) = \mathcal{L}_{v_1}^{[a_1, b_1]}(x_1) \mathcal{L}_{v_2}^{[a_2, b_2]}(x_2) \mathcal{L}_{v_3}^{[a_3, b_3]}(x_3),$$

for $\mathbf{x} = (x_1, x_2, x_3) \in \mathbb{R}^3$.

4.2. Setup of the \mathcal{H}^2 -matrices

In order to set up the \mathcal{H}^2 -matrix approximation, we have to create the coefficient matrices $S^{\tau,\sigma}$ for $(\tau, \sigma) \in P_{\text{far}}$, the basis matrices V^τ for $\tau \in \mathcal{C}$ with $\text{sons}(\tau) = \emptyset$, the transfer matrices $T^{\tau',\tau}$ for $\tau \in \mathcal{C}$ with $\text{sons}(\tau) \neq \emptyset$, and the near-field matrices $G^{\tau,\sigma}$ for $(\tau, \sigma) \in P_{\text{near}}$.

The computation of $S^{\tau,\sigma}$ is straightforward: using the interpolation points $(\mathbf{x}_\nu^\tau)_{\nu \in M}$ and $(\mathbf{x}_\mu^\sigma)_{\mu \in M}$ defined in Section 4.1, we have to evaluate the kernel function Φ

$$S_{\nu\mu}^{\tau,\sigma} = \Phi(\mathbf{x}_\nu^\tau, \mathbf{x}_\mu^\sigma).$$

The matrices V^τ satisfy

$$V_{iv}^\tau = \int_\tau \chi_t(\mathbf{x}) \mathcal{L}_v^\tau(\mathbf{x}) \, d\mathbf{x} = \int_t \mathcal{L}_v^\tau(\mathbf{x}) \, d\mathbf{x},$$

since χ_t is the characteristic function of the triangle t . This integral can be computed by standard quadrature techniques, since the integrand \mathcal{L}_v^τ is a polynomial.

Let $B^\tau = [a_1, b_1] \times [a_2, b_2] \times [a_3, b_3]$ and $B^{\tau'} = [c_1, d_1] \times [c_2, d_2] \times [c_3, d_3]$. Then the transfer matrix $T^{\tau',\tau}$ is given by

$$T_{\nu'\nu}^{\tau',\tau} = \mathcal{L}_{\nu'}^{\tau'}(\mathbf{x}_{\nu'}^{\tau'}) = \mathcal{L}_{\nu_1}^{[a_1,b_1]}(x_{\nu_1}^{[c_1,d_1]}) \mathcal{L}_{\nu_2}^{[a_2,b_2]}(x_{\nu_2}^{[c_2,d_2]}) \mathcal{L}_{\nu_3}^{[a_3,b_3]}(x_{\nu_3}^{[c_3,d_3]}),$$

i.e., is the Kronecker product of three matrices computed by evaluating one-dimensional Lagrange polynomials. This allows us to compute the matrices $T^{\tau',\tau}$ efficiently. Alternatively, we can reduce the memory requirements by storing the Kronecker factors instead of the full matrix.

The near-field matrices $G^{\tau,\sigma}$ are computed by a semi-analytical approach ² for the double integrals with singular kernels, where the interior integral is calculated analytically and the exterior integral is evaluated by a Gaussian quadrature scheme. Some of the analytical integrations can be looked up in [13].

Remark 7 (Compact storage). The matrices V^τ , $V^{\tau,\ell}$ and $W^{\tau,\ell}$ have to be stored for the leaf clusters only, since the matrix–vector multiplication algorithm uses the transfer matrices for all other clusters.

The transfer matrices $T^{\tau',\tau}$ and the coefficient matrices $S^{\tau,\sigma}$ coincide for the bilinear forms $g(\cdot, \cdot)$ and $b(\cdot, \cdot)$, so they have to be stored only once.

Most of the entries in the matrices V^τ , $V^{\tau,\ell}$, $W^{\tau,\ell}$ and $G^{\tau,\sigma}$ are zero, since only the basis functions whose supports intersect the domains τ or σ lead to non-zero entries. Since τ and σ are leaves of the cluster tree, they can be assumed to be quite small and therefore intersect only a small number of triangles. By storing only the non-zero entries, the matrices can be treated efficiently.

The same holds for the near-field matrices: they are sparse due to the locality of leaves of the cluster tree, so we can apply the usual techniques for sparse matrices.

Differently from standard \mathcal{H}^2 -techniques, degrees of freedom located in the edges and nodes can appear multiple times during the course of the matrix–vector multiplication due to the fact that our clustering technique is based on the triangles of the grid, not on the degrees of freedom. The compression rates in Fig. 1 show that the performance of our method is still good.

² Private communication with Dr. Olaf Steinbach, University of Stuttgart.

The diagonal part of $M_{\mathfrak{R}}$ and the matrix $M_{\mathfrak{I}}$ consist of scalar products of linear edge functions, therefore they are sparse and do not have to be compressed. The non-diagonal part of $M_{\mathfrak{R}}$ is the scalar single layer potential of Eq. (13) with constant basis functions.

All in all, the elaborate operators can be compressed without major difficulties by the \mathcal{H}^2 -matrix approximation technique. A large number of matrix entries must be calculated and stored only once and can be reused for different operators. This automatically saves a lot of storage and enhances the efficiency of the algorithm.

4.3. Preconditioning

A preconditioned conjugate residual method is used for solving Eq. (5).

We use the block diagonal preconditioner P given by

$$P := \begin{pmatrix} M_{\mathfrak{R}} + M_{\mathfrak{I}} & 0 & 0 & 0 \\ 0 & M_{\mathfrak{R}} + M_{\mathfrak{I}} & 0 & 0 \\ 0 & 0 & Q & 0 \\ 0 & 0 & 0 & Q \end{pmatrix}, \quad (26)$$

which is an approximation of the diagonal part of the operator itself, see (5). For the strongly coupled FEM/BEM formulation, it is possible to show that their exist bounds for the spectrum of the preconditioned system that do not depend on the meshwidth.

Therefore we can expect the preconditioner to work well for the impedance model, too. The condition numbers were computed for a small test example, and decreased from 38,000 for the original system to 4 for the preconditioned system.

We use the conjugate gradients method with a simple Jacobi preconditioner to solve the diagonal blocks of P approximately.

All the operators of P , namely $M_{\mathfrak{R}}$, $M_{\mathfrak{I}}$ and Q , are also part of Eq. (5), thus are already compressed if the \mathcal{H}^2 -matrix approximation is applied. So the \mathcal{H}^2 -matrix approximation also pays for the preconditioner.

The near-field part Q_{near} of the operator Q can be inverted by a *sparse Gaussian elimination scheme* as long as the dimensions are moderate, giving us a preconditioner for the third and fourth row of P .

5. Numerical experiments

The predicted behavior of the \mathcal{H}^2 -matrix approximation method was tested for the three operators $M_{\mathfrak{R}}$, Q , and B of Eq. (5) by using the typical geometry of the induction heating setting, as shown in Fig. 7. The storage and cpu-time requirements for solving Eq. (5) with the uncompressed standard method (i.e., semi-analytic quadrature of constant order used to compute a dense matrix) and the interpolation-based \mathcal{H}^2 -method were compared. The order of interpolation was set to be 2 and the simplified admissibility condition (23) with $\eta' = 0.99$ was used. This turned out to lead to sufficiently small approximation errors.

Table 1 shows the performance of the \mathcal{H}^2 -matrix approximation compared with the use of dense matrices.

The solver was stopped when the residual had decreased to a value below 0.0001 times the residual in the first step. The documented time is the time needed for filling the matrices and solving the system. The number of unknowns is defined as

$$n := 2 \times (\#\mathcal{E} + \#\mathcal{N} + \text{number of holes in the workpiece}),$$

and for the storage requirements one finds

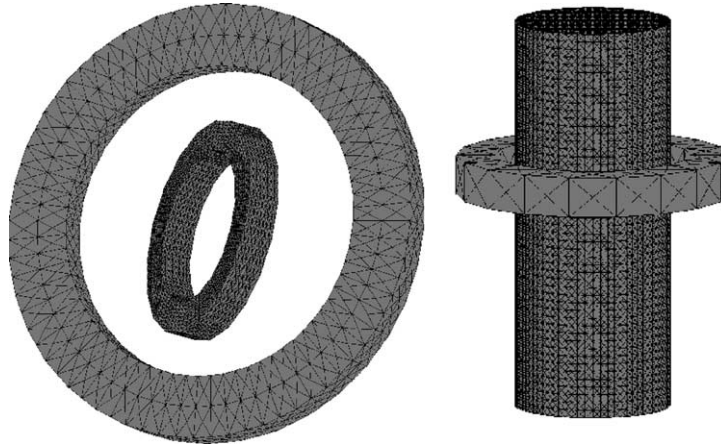


Fig. 7. Test geometries A and B.

Table 1
Time, storage, and errors for the impedance model

| n | Std Strg (MB) | \mathcal{H}^2 Strg (MB) | Std time (min) | \mathcal{H}^2 Time (min) | Rel. error |
|--------|---------------|---------------------------|----------------|----------------------------|-------------------|
| 2948 | 22.7 | 15.5 | 20 | 15 | 1.4 ₋₃ |
| 6916 | 125.3 | 35.0 | 114 | 43 | 2.5 ₋₃ |
| 11,420 | 342.0 | 71.0 | 402 | 96 | 1.5 ₋₃ |
| 23,840 | 954.1 | 93.4 | 876 | 144 | 2.2 ₋₃ |
| 46,724 | 5725.0 | 333.0 | – | 540 | – |

Standard storage = $(2 \times (\#\mathcal{E})^2 + (\#\mathcal{E} \times \#\mathcal{N}) + (\#\mathcal{N})^2) \times \text{sizeof}(\text{double})$ bytes,

\mathcal{H}^2 storage = Storage needed for the \mathcal{H}^2 approximation.

The relative error is defined as the difference between the surface current $\gamma_{D\mathbf{j}} := \sigma \cdot \gamma_D \mathbf{E}$ of the \mathcal{H}^2 -solution and that of the solution obtained without compression

$$\text{Rel. error} := \int_{\Gamma} \frac{\|\gamma_{D\mathbf{j}_{\mathcal{H}^2}}(\mathbf{x}) - \gamma_{D\mathbf{j}_{st.}}(\mathbf{x})\|}{\|\gamma_{D\mathbf{j}_{st.}}(\mathbf{x})\|} dS_{\mathbf{x}},$$

because the current is the most important entity for the calculation of the inductive heating.

The first four rows were produced by interpolation on the geometry A of Fig. 7 and the last on geometry B. Storage requirements and calculation times are strongly reduced, and geometries consisting of 20,000 surface faces can be calculated. In each case the time for filling the matrices amounts to 90% of the total time. The experiments were carried out on a Sun Ultra 450 computer with a 300 MHz Ultrasparc II processor.

In order to investigate the influence of the admissibility parameter η' (cf. (23)) and the order m of the polynomial expansion (cf. (15)), we use a fixed approximation of test geometry B from Fig. 7 with 3312 edges and 1106 vertices, leading to 8836 degrees of freedom, and examine the results for different combinations of parameters in Table 2. The table contains the times for creating an approximation (build), performing a matrix–vector multiplication (MVM), the amount of required storage (storage) and the error (Op. error) of the approximate matrix in the relative operator norm, i.e., $\|K - K_{\mathcal{H}^2}\|_2 / \|K\|_2$. We can see

Table 2
Influence of the admissibility parameter η' and the approximation order m

| η' | m | Build (min) | Storage (MB) | MVM (s) | Op. error |
|----------|-----|-------------|--------------|---------|-------------------|
| 1.4 | 2 | 38.1 | 52 | 3.3 | 1.1 ₋₄ |
| 1.1 | | 44.5 | 59 | 3.8 | 9.2 ₋₅ |
| 0.8 | | 61.5 | 79 | 4.5 | 4.5 ₋₅ |
| 0.5 | | 70.9 | 90 | 5.0 | 1.7 ₋₅ |
| 1.4 | 1 | 13.7 | 28 | 3.5 | 6.7 ₋₄ |
| | 2 | 38.1 | 52 | 3.3 | 1.1 ₋₄ |
| | 3 | 50.4 | 61 | 5.5 | 2.1 ₋₅ |
| | 4 | 70.3 | 78 | 8.2 | 1.0 ₋₅ |
| | 5 | 101.8 | 108 | 11.0 | 4.3 ₋₆ |
| Standard | | 127.4 | 205 | 8.6 | 0 |

that, in order to achieve high accuracy, increasing the approximation order is more advisable than decreasing the admissibility parameter η' . The slight increase when switching from $m = 1$ to 2 is due to the fact that the maximal leaf size C_{lf} (cf. Fig. (3)) is chosen to be proportional to m^3 , so the cluster tree for $m = 2$ is much smaller than that for $m = 1$, which results in better performance.

A special feature is implemented due to the fact that the workpiece rotates. This means that the current has to be calculated for several positions of the workpiece per rotation. Therefore, it is desirable to find a clever way for reusing entities, which have already been calculated:

The BEM operators change only for items that are moving relative to each other. Parts of the operators which describe workpiece/workpiece interactions (w/w) or inductor/inductor interactions (i/i) need to be calculated only once. But how to reach this goal in the \mathcal{H}^2 context where everything is linked together in the tree? The option chosen here is to apply the geometric bisection algorithm (cf. Fig. 3) separately to workpiece and inductor. Then each cluster consists exclusively of workpiece triangles or inductor triangles. In this case, the matrices V , W and T do not change during the rotation and have to be calculated only once (except for the useless T of the root). After merging the two resulting trees under one big root, one finds the situation of Fig. 8. If the block partition algorithm (cf. Fig. 4) is now applied to this new tree, each pair of

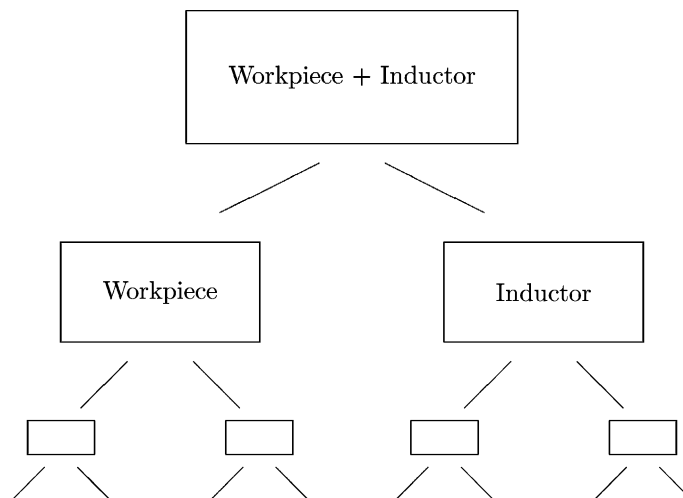


Fig. 8. Cluster tree for the inductor/workpiece domain.

admissible clusters belongs either to w/w or i/i or to the interaction w/i between workpiece and inductor. The near-field and the matrices S have to be refreshed only for the w/i pairs. This is a big advantage because the interesting part with the biggest number of triangles is the workpiece, and w/w does not have to be refreshed.

References

- [1] S. Amini, A. Profit, Multi-level fast multipole solution of the scattering problem, *Eng. Anal. Bound. Elem.* 27 (2003) 547–564.
- [2] H. Ammari, A. Buffa, J.-C. Nédélec, A justification of eddy currents model for the Maxwell equations, *SIAM J. Appl. Math.* 60 (2000) 1805–1823.
- [3] S. Börm, L. Grasedyck, W. Hackbusch, Introduction to hierarchical matrices with applications, *Eng. Anal. Bound. Elem.* 27 (2003) 405–422.
- [4] S. Börm, W. Hackbusch, \mathcal{H}^2 -matrix approximation of integral operators by interpolation, *Appl. Numer. Math.* 43 (2002) 129–143.
- [5] S. Börm, M. Löhndorf, and J.M. Melenk, Approximation of integral operators by variable-order interpolation, Technical Report 82, Max Planck Institute for Mathematics in the Sciences, 2002, *Numerische Mathematik* (to appear).
- [6] H. Dirks, Quasi-stationary fields for microelectronic applications, *Electr. Eng.* 79 (1996) 145–155.
- [7] K. Giebermann, Multilevel approximation of boundary integral operators, *Computing* 67 (2001) 183–207.
- [8] L. Greengard, V. Rokhlin, A fast algorithm for particle simulations, *J. Comput. Phys.* 73 (1987) 325–348.
- [9] W. Hackbusch, A sparse matrix arithmetic based on \mathcal{H} -matrices. Part I: Introduction to \mathcal{H} -matrices, *Computing* 62 (1999) 89–108.
- [10] W. Hackbusch, B. Khoromskij, S. Sauter, H. Bungartz, R. Hoppe, C. Zenger, On \mathcal{H}^2 -matrices, in: *Lectures on Applied Mathematics*, Springer, Berlin, 2000, pp. 9–29.
- [11] W. Hackbusch, Z.P. Nowak, On the fast matrix multiplication in the boundary element method by panel clustering, *Numer. Math.* 54 (1989) 463–491.
- [12] R. Hiptmair, Symmetric coupling for eddy current problems, *SIAM J. Numer. Anal.* 40 (2002) 41–65.
- [13] J. Ostrowski, *Boundary Element Methods for Inductive Hardening*, PhD thesis, University of Tübingen, Germany, 2003. Available from <http://w210.ub.uni-tuebingen.de/dbt/volltexte/2003/672>.
- [14] V. Rokhlin, Rapid solution of integral equations of classical potential theory, *J. Comput. Phys.* 60 (1985) 187–207.
- [15] S. Sauter, Variable order panel clustering (extended version), Tech. Rep. 52, Max-Planck-Institut für Mathematik, Leipzig, Germany, 1999.
- [16] S. Sauter, Variable order panel clustering, *Computing* 64 (2000) 223–261.
- [17] C. Schwab, O. Sterz, A scalar BEM for time harmonic eddy current problems with impedance boundary conditions, in: M. Günther, D. Hecht, U. van Rienen (Eds.), *Scientific Computing in Electrical Engineering, Lecture Notes in Computational Science and Engineering*, vol. 18, Springer, Berlin, 2001, pp. 129–136.