# SCHRIFTENREIHE DER FAKULTÄT FÜR MATHEMATIK 

Time-Harmonic Electro-Magnetic Scattering in Exterior Weak Lipschitz Domains with Mixed Boundary Conditions

by
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# Time-Harmonic Electro-Magnetic Scattering in Exterior Weak Lipschitz Domains with Mixed Boundary Conditions 

FRANK OSTERBRINK AND DIRK PAULY

Abstract. This paper treats the time-harmonic electro-magnetic scattering or radiation problem governed by Maxwell's equations, i.e.,

$$
\begin{array}{rlll}
\operatorname{rot} E+i \omega \mu H=F & \text { in } \Omega, & E \times \nu=0 & \text { on } \Gamma_{1}, \\
\operatorname{rot} H-i \omega \varepsilon E=G & \text { in } \Omega, & H \times \nu=0 & \text { on } \Gamma_{2},
\end{array}
$$

where $\omega \in \mathbb{C} \backslash\{0\}$ and $\Omega \subset \mathbb{R}^{3}$ is an exterior weak Lipschitz domain with boundary $\Gamma:=\partial \Omega$ divided into two disjoint parts $\Gamma_{1}$ and $\Gamma_{2}$. We will present a solution theory using the framework of polynomially weighted Sobolev spaces for the rotation and divergence. For the physically interesting case $\omega \in \mathbb{R} \backslash\{0\}$ we will show a Fredholm alternative type result to hold using the principle of limiting absorption introduced by Eidus in the 1960's. The necessary a-priori-estimate and polynomial decay of eigenfunctions for the Maxwell equations will be obtained by transferring well known results for the Helmholtz equation using a suitable decomposition of the fields $E$ and $H$. The crucial point for existence is a local version of Weck's selection theorem, also called Maxwell compactness property.

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## 1. Introduction

The equations that describe the behavior of electro-magnetic fields in a region $\Omega \subset \mathbb{R}^{3}$, first completely formulated by J. C. Maxwell in 1864, are

$$
\begin{array}{rlrl}
-\operatorname{rot} H+\partial_{t} D & =-I, & \operatorname{rot} E+\partial_{t} B=0, & \\
\text { in } \Omega, \\
\operatorname{div} D & =\rho, & \operatorname{div} B=0, & \\
\text { in } \Omega,
\end{array}
$$

where $E, H$ are the electric resp. magnetic field, $D, B$ represent the displacement current and magnetic induction and $I, \rho$ describe the current density resp. the charge density. Excluding, e.g., ferromagnetic resp. ferroelectric materials, the parameters linking $E$ and $H$ with $D$ and $B$ are often assumed to be of the linear form $D=\varepsilon E$ and $B=\mu H$, where $\varepsilon$ and $\mu$ are matrix-valued functions describing the permittivity and permeability of the medium filling $\Omega$. Here we are especially interested in the case of

[^0]an exterior domain $\Omega \subset \mathbb{R}^{3}$, i.e., a connected open subset of with compact complement. Applying the divergence to the first two equations we see that the latter two equations are implicitly included in the first two and may be omitted. Hence, neglecting the static case, Maxwell's equations reduce to
$$
-\operatorname{rot} H+\partial_{t}(\varepsilon E)=F, \quad \operatorname{rot} E+\partial_{t}(\mu H)=G, \quad \text { in } \Omega,
$$
with arbitrary right hand sides $F, G$. Among the wide range of phenomena described by these equations one important case is the discussion of "time-harmonic" electro-magnetic fields where all fields vary sinusoidally in time with a single frequency $\omega \in \mathbb{C} \backslash\{0\}$, i.e.,
$$
E(x, t)=E(x) e^{i \omega t}, \quad H(x, t)=H(x) e^{i \omega t}, \quad G(x, t)=G(x) e^{i \omega t}, \quad F(x, t)=F(x) e^{i \omega t}
$$

Substituting this ansatz into the equations (or using Fourier-transformation in time) and assuming that $\varepsilon$ and $\mu$ are time-independent we are lead to what is called "time-harmonic Maxwell's equations":

$$
\begin{equation*}
\operatorname{rot} E+i \omega \mu H=G, \quad-\operatorname{rot} H+i \omega \varepsilon E=F, \quad \text { in } \Omega . \tag{1.1}
\end{equation*}
$$

This system equipped with suitable boundary conditions describes, e.g., the scattering of time-harmonic electro-magnetic waves which is of high interest in many applications like geophysics, medicine, electrical engineering, biology and many others.

First existence results concerning boundary value problems for the time-harmonic Maxwell system in bounded and exterior domains have been given by Müller [13], [12]. He studied isotropic and homogeneous media and used integral equation methods. Using alternating differential forms, Weyl [29] investigated theseequations on Riemannian manifolds of arbitrary dimension, while Werner [28] was able to transfer Müller's results to the case of inhomogeneous but isotropic media. However, for general inhomogeneous anisotropic media and arbitrary exterior domains, boundary integral methods are less useful since they heavily depend on the explicit knowledge of the fundamental solution and strong assumptions on boundary regularity. That is why Hilbert space methods are a promising alternative. Unfortunately, Maxwell's equations are non elliptic, hence it is in general not possible to estimate all first derivatives of a solution. In [9] Leis could overcome this problem by transforming the boundary value problem for Maxwell's system into a boundary value problem for the Helmholtz equation, assuming that the medium filling $\Omega$, is inhomogeneous and anisotropic within a bounded subset of $\Omega$. Nevertheless, he still needed boundary regularity to gain equivalence of both problems. But also for nonsmooth boundaries Hilbert space methods are expedient. In fact, as shown by Leis [10], it is sufficient that $\Omega$ satisfies a certain selection theorem, later called Weck's selection theorem or Maxwell compactness property, which holds for a class of boundaries much larger than those accessible by the detour over $\mathrm{H}^{1}$. (cf. Weck [24], Costabel [2] and Picard, Weck, Witsch [20] ). See [11] for a detailed monograph and [1] for the most recent result and an overview. The most recent result regarding a solution theory is due to Pauly [16] (see also [14]) and in its structure comparable to the results of Picard [18] and Picard, Weck \& Witsch [20]. While all these results above have been obtained for full boundary conditions, in the present paper we study the case of mixed boundary conditions. More precisely, we are interested in solving the system (1.1) for $\omega \in \mathbb{C} \backslash\{0\}$ in an exterior domain $\Omega \subset \mathbb{R}^{3}$, where we assume that $\Gamma:=\partial \Omega$ is decomposed into two relatively open subsets $\Gamma_{1}$ and its complement $\Gamma_{2}:=\Gamma \backslash \bar{\Gamma}_{1}$ and impose homogeneous boundary conditions, which in classical terms can be written as

$$
\begin{equation*}
\nu \times E=0 \text { on } \Gamma_{1}, \quad \nu \times H=0 \text { on } \Gamma_{2}, \quad(\nu: \text { outward unit normal }) \tag{1.2}
\end{equation*}
$$

Conveniently, we can apply the same methods as in [15] ( see also Picard, Weck \& Witsch [20], Weck \& Witsch $[27],[25])$ to construct a solution. Indeed, most of the proofs carry over practically verbatim. For $\omega \in \mathbb{C} \backslash \mathbb{R}$ the solution theory is obtained by standard Hilbert space methods as $\omega$ belongs to the resolvent set of the Maxwell operator. In the case of $\omega \in \mathbb{R} \backslash\{0\}$, i.e., $\omega$ is in the continuous spectrum of the Maxwell operator, we use the limiting absorption principle introduced by Eidus [4] and approximate solutions to $\omega \in \mathbb{R} \backslash\{0\}$ by solutions corresponding to $\omega \in \mathbb{C} \backslash \mathbb{R}$. This will be sufficient to show a generalized Fredholm alternative (cf. our main result, Theorem 3.10) to hold. The essential ingredients needed for the limit process are

- the polynomial decay of eigensolutions,
- an a-priori-estimate for solutions corresponding to non-real frequencies
- a Helmholtz-type decomposition,
- and Weck's local selection theorem (WLST), that is,

$$
\mathbf{R}_{\Gamma_{1}}(\Omega) \cap \varepsilon^{-1} \mathbf{D}_{\Gamma_{2}}(\Omega) \longleftrightarrow \mathrm{L}_{\mathrm{loc}}^{2}(\bar{\Omega}) \text { is compact. }
$$

While the first two are obtained by transferring well known results for the scalar Helmholtz equation to the time-harmonic Maxwell equations using a suitable decomposition of the fields $E$ and $H$, Lemma 4.1, the last one is an assumption on the quality of the boundary. As we will see, WLST is an immediate consequence of Weck's selection theorem (WST), i.e.,

$$
\mathbf{R}_{\Gamma_{1}}(\Theta) \cap \varepsilon^{-1} \mathbf{D}_{\Gamma_{2}}(\Theta) \longleftrightarrow \mathrm{L}^{2}(\Theta) \text { is compact, }
$$

which holds in bounded weak Lipschitz domains $\Theta \subset \mathbb{R}^{3}$, but fails in unbounded like exterior domains (cf. Bauer, Pauly, Schomburg [1] and the references therein ). For strong Lipschitz- domains see Jochmann [7] and Fernandes, Gilardis [5].

## 2. Preliminaries and Notations

Let $\mathbb{Z}, \mathbb{N}, \mathbb{R}$ and $\mathbb{C}$ be the usual sets of integers, natural, real and complex numbers, respectively. Furthermore, let $i$ be the imaginary unit, $\operatorname{Re} z, \operatorname{Im} z$ and $\bar{z}$ real part, imaginary part and complex conjugate of $z \in \mathbb{C}$, as well as

$$
\mathbb{R}_{+}:=\{s \in \mathbb{R} \mid s>0\}, \quad \mathbb{C}_{+}:=\{z \in \mathbb{C} \mid \operatorname{Im} z \geq 0\} \quad \text { and } \quad \mathbb{I}:=\{(2 m+1) / 2 \mid m \in \mathbb{Z} \backslash\{0\}\}
$$

For $x \in \mathbb{R}^{n}$ with $x \neq 0$ we set $r(x):=|x|$ and $\xi(x):=x /|x|\left(|\cdot|:\right.$ Euclidean norm in $\left.\mathbb{R}^{n}\right)$. Moreover, $\mathrm{U}(\tilde{r})$ resp. $\mathrm{B}(\tilde{r})$ indicate the open resp. closed ball of radius $\tilde{r}$ in $\mathbb{R}^{n}$ centered in the origin and we define

$$
\mathrm{S}(\tilde{r}):=\mathrm{B}(\tilde{r}) \backslash \mathrm{U}(\tilde{r}), \quad \check{\mathrm{U}}(\tilde{r}):=\mathbb{R}^{3} \backslash \mathrm{~B}(\tilde{r}), \quad \mathrm{G}(\tilde{r}, \hat{r}):=\check{\mathrm{U}}(\tilde{r}) \cap \mathrm{U}(\hat{r}),
$$

with $\hat{r}>\tilde{r}$. If $f: \mathrm{X} \longrightarrow \mathrm{Y}$ is a function mapping X to Y the restriction of $f$ to a subset $\mathrm{U} \subset \mathrm{X}$ will be marked with $\left.f\right|_{U}$ and $\mathcal{D}(f), \mathcal{N}(f)$, supp $f$ resp. $\mathcal{R}(f)$ denote domain of definition, kernel, support and range of $f$. Given two functions $f, g: \mathbb{R}^{n} \longrightarrow \mathbb{C}^{k}$ we write

$$
f=\mathcal{O}(g) \text { for } r \longrightarrow \infty \quad \text { if and only if } \quad \exists c>0 \exists \tilde{r}>0 \forall x \in \check{\mathrm{U}}(\tilde{r}): \quad|f(x)| \leq c \cdot g(x)
$$

For $\mathrm{X}, \mathrm{Y}$ subspaces of a normed vector space $\mathrm{V}, \mathrm{X}+\mathrm{Y}, \mathrm{X} \dot{+} \mathrm{Y}$ and $\mathrm{X} \oplus \mathrm{Y}$ indicate the sum, the direct sum and the orthogonal sum of X and Y , where in the last case we presume the existence of a scalar product $\langle\cdot, \cdot\rangle_{\mathrm{V}}$ on V. Moreover, $\langle\cdot, \cdot\rangle_{\mathrm{X} \times \mathrm{Y}}$ resp. $\|\cdot\|_{\mathrm{X} \times \mathrm{Y}}$ denote the natural scalar product resp. induced norm on $\mathrm{X} \times \mathrm{Y}$ and if $\mathrm{X}=\mathrm{Y}$, we stay "with X instead of $\mathrm{X} \times \mathrm{X}$.
2.1. General Assumptions and Weighted Spaces. Unless stated otherwise, from now on and throughout this paper it is assumed that $\Omega \subset \mathbb{R}^{3}$ is an exterior weak Lipschitz domain with weak Lipschitz interface in the sense of [1, Definition 2.3, Definition 2.5], which in principle means that $\Gamma:=\partial \Omega$ is a Lipschitz-manifold and $\Gamma_{1}$ resp. $\Gamma_{2}$ are Lipschitz-submanifolds of $\Gamma$. For later purpose we fix $r_{0}>0$ such that $\mathbb{R}^{3} \backslash \Omega \Subset \mathrm{U}\left(r_{0}\right)$ and define for arbitrary $\tilde{r} \geq r_{0}$

$$
\Omega(\tilde{r}):=\Omega \cap \mathrm{U}(\tilde{r})
$$

With $r_{k}:=2^{k} r_{0}, k \in \mathbb{N}$ and $\tilde{\eta} \in \mathrm{C}^{\infty}(\mathbb{R})$ such that

$$
\begin{equation*}
0 \leq \tilde{\eta} \leq 1, \quad \operatorname{supp} \tilde{\eta} \subset(-\infty, 2-\delta),\left.\quad \tilde{\eta}\right|_{(-\infty, 1+\delta)}=1, \tag{2.1}
\end{equation*}
$$

for some $0<\delta<1$, we define functions $\eta, \check{\eta}, \eta_{k}, \check{\eta}_{k} \in \mathrm{C}^{\infty}\left(\mathbb{R}^{3}\right)$ by

$$
\eta(x):=\tilde{\eta}\left(r(x) / r_{0}\right), \quad \check{\eta}(x):=1-\eta(x), \quad \eta_{k}(x):=\tilde{\eta}\left(r(x) / r_{k}\right), \quad \text { resp. } \quad \check{\eta}_{k}(x):=1-\eta_{k}(x)
$$

meaning

$$
\begin{array}{lllll}
\operatorname{supp} \eta \subset \mathrm{B}\left(r_{1}\right) & \text { with } \quad \eta=1 \text { on } \mathrm{U}\left(r_{0}\right), & \text { resp. } & \operatorname{supp} \eta_{k} \subset \mathrm{U}\left(r_{k+1}\right) & \text { with } \\
\operatorname{supp} \check{\eta} \subset \check{\mathrm{U}}\left(r_{0}\right) & \text { with } \check{\eta}=1 \text { on } \check{\mathrm{U}}\left(r_{1}\right), & & \operatorname{supp} \check{\eta}_{k} \subset \check{\mathrm{U}}\left(r_{k}\right) & \text { with } \mathrm{U}\left(r_{k}\right) \\
\check{\eta}_{k}=1 \text { on } \check{\mathrm{U}}\left(r_{k+1}\right) .
\end{array}
$$

These functions will later be utilized for particular cut-off procedures.
Next we introduce our notations for Lebesgue and Sobolev spaces needed in the following discussion. Note that we will not indicate whether the elements of these spaces are scalar functions or vector fields. This will be always clear from the context. The example ${ }^{\mathrm{i}}$

$$
E:=\nabla \ln (r) \in \mathrm{H}_{\mathrm{loc}}^{1}(\check{\mathrm{U}}(1)), \quad \nu \times\left. E\right|_{\mathrm{S}(1)}=0, \quad \operatorname{rot} E=0 \in \mathrm{~L}^{2}(\check{\mathrm{U}}(1)), \quad \operatorname{div} E=r^{-2} \in \mathrm{~L}^{2}(\check{\mathrm{U}}(1)),
$$

shows that a standard $\mathrm{L}^{2}$-setting is not appropriate for exterior domains. Even for square-integrable right hand sides we cannot expect to find square-integrable solutions. Indeed it turns out that we have to work in weighted Lebesgue and Sobolev spaces to develop a solution theory. With $\rho:=\left(1+r^{2}\right)^{1 / 2}$ we introduce for an arbitrary domain $\Omega \subset \mathbb{R}^{3}, t \in \mathbb{R}$ and $m \in \mathbb{N}$ :

$$
\begin{gathered}
\mathrm{L}_{\mathrm{t}}^{2}(\Omega):=\left\{w \in \mathrm{~L}_{\mathrm{loc}}^{2}(\Omega) \mid \rho^{t} w \in \mathrm{~L}^{2}(\Omega)\right\} . \\
\mathbf{H}_{\mathrm{t}}^{\mathrm{m}}(\Omega):=\left\{w \in \mathrm{~L}_{\mathrm{t}}^{2}(\Omega)|\forall| \alpha \mid \leq m: \partial^{\alpha} w \in \mathrm{~L}_{\mathrm{t}}^{2}(\Omega)\right\}, \\
\mathrm{H}_{\mathrm{t}}^{\mathrm{m}}(\Omega):=\left\{w \in \mathrm{~L}_{\mathrm{t}}^{2}(\Omega)|\forall| \alpha \mid \leq m: \partial^{\alpha} w \in \mathrm{~L}_{\mathrm{t}+|\alpha|}^{2}(\Omega)\right\}, \\
\mathbf{R}_{\mathrm{t}}(\Omega):=\left\{E \in \mathrm{~L}_{\mathrm{t}}^{2}(\Omega) \mid \operatorname{rot} E \in \mathrm{~L}_{\mathrm{t}}^{2}(\Omega)\right\}, \quad \mathrm{R}_{\mathrm{t}}(\Omega):=\left\{E \in \mathrm{~L}_{\mathrm{t}}^{2}(\Omega) \mid \operatorname{rot} E \in \mathrm{~L}_{\mathrm{t}+1}^{2}(\Omega)\right\}, \\
\mathbf{D}_{\mathrm{t}}(\Omega):=\left\{H \in \mathrm{~L}_{\mathrm{t}}^{2}(\Omega) \mid \operatorname{div} H \in \mathrm{~L}_{\mathrm{t}}^{2}(\Omega)\right\}, \quad \mathrm{D}_{\mathrm{t}}(\Omega):=\left\{H \in \mathrm{~L}_{\mathrm{t}}^{2}(\Omega) \mid \operatorname{div} H \in \mathrm{~L}_{\mathrm{t}+1}^{2}(\Omega)\right\} .
\end{gathered}
$$

where $\alpha=\left(\alpha_{1}, \alpha_{2}, \alpha_{3}\right) \in \mathbb{N}^{3}$ is a multi-index and $\partial^{\alpha} w:=\partial_{1}^{\alpha_{1}} \partial_{2}^{\alpha_{2}} \partial_{3}^{\alpha_{3}} w$, rot $E$ resp. div $H$ are the usual distributional or weak derivatives. Equipped with the induced norms

$$
\begin{aligned}
& \|w\|_{\mathrm{L}_{\mathbf{t}}^{2}(\Omega)}^{2}:=\left\|\rho^{t} w\right\|_{\mathrm{L}^{2}(\Omega)}^{2}, \\
& \|w\|_{\mathbf{H}_{\mathrm{t}}^{\mathrm{m}}(\Omega)}^{2}:=\sum_{|\alpha| \leq m}\left\|\partial^{\alpha} w\right\|_{\mathrm{L}_{\mathrm{t}}^{2}(\Omega)}^{2}, \\
& \|w\|_{\mathrm{H}_{\mathrm{t}}^{\mathrm{m}}(\Omega)}^{2}:=\sum_{|\alpha| \leq m}\left\|\partial^{\alpha} w\right\|_{\mathrm{L}_{\mathrm{t}+|\alpha|}^{2}(\Omega)}^{2}, \\
& \|E\|_{\mathbf{R}_{\mathrm{t}}(\Omega)}^{2}:=\|E\|_{\mathrm{L}_{\mathrm{t}}^{2}(\Omega)}^{2}+\|\operatorname{rot} E\|_{\mathrm{L}_{\mathrm{t}}^{2}(\Omega)}^{2}, \quad\|E\|_{\mathrm{R}_{\mathrm{t}}(\Omega)}^{2}:=\|E\|_{\mathrm{L}_{\mathrm{t}}^{2}(\Omega)}^{2}+\|\operatorname{rot} E\|_{\mathrm{L}_{\mathrm{t}+1}^{2}(\Omega)}^{2}, \\
& \|H\|_{\mathbf{D}_{\mathrm{t}}(\Omega)}^{2}:=\|H\|_{\mathrm{L}_{\mathrm{t}}^{2}(\Omega)}^{2}+\|\operatorname{div} H\|_{\mathrm{L}_{\mathrm{t}}^{2}(\Omega)}^{2}, \quad\|H\|_{\mathrm{D}_{\mathrm{t}}(\Omega)}^{2}:=\|H\|_{\mathrm{L}_{\mathrm{t}}^{2}(\Omega)}^{2}+\|\operatorname{div} H\|_{\mathrm{L}_{t+1}^{2}(\Omega)}^{2},
\end{aligned}
$$

they become Hilbert spaces. As usual, the subscript "loc" resp. "vox" indicates local square-integrability resp. bounded support. Please note, that the bold spaces with weight $t=0$ correspond to the classical Lebesque and Sobolev spaces and for bounded domains "non-weighted" and weighted spaces even coincide:

$$
\Omega \subset \mathbb{R}^{3} \text { bounded } \Longrightarrow \forall t \in \mathbb{R}:\left\{\begin{array}{l}
\mathbf{H}_{\mathrm{t}}^{1}(\Omega)=\mathrm{H}_{\mathrm{t}}^{1}(\Omega)=\mathrm{H}_{0}^{1}(\Omega) \\
\mathbf{R}_{\mathrm{t}}(\Omega)=\mathrm{R}_{\mathrm{t}}(\Omega)=\mathrm{R}_{0}(\Omega) \\
\mathbf{D}_{\mathrm{t}}(\Omega)=\mathrm{D}_{\mathrm{t}}(\Omega)=\mathrm{D}_{0}(\Omega)
\end{array}\right.
$$

Besides the usual set $\complement^{\infty}(\Omega)$ of test fields (resp. test functions) we introduce

$$
\mathrm{C}_{\Gamma_{\mathrm{i}}}^{\infty}(\Omega):=\left\{\left.\varphi\right|_{\Omega} \mid \varphi \in \dot{\mathrm{C}}^{\infty}\left(\mathbb{R}^{3}\right) \text { and } \operatorname{dist}\left(\operatorname{supp} \varphi, \Gamma_{i}\right)>0\right\}, \quad \mathrm{i}=1,2
$$

to formulate boundary conditions in the weak sense:

$$
\begin{align*}
& \mathrm{H}_{\mathrm{t}, \mathrm{\Gamma}_{\mathrm{i}}}^{\mathrm{m}}(\Omega):=\overline{\mathrm{C}_{\mathrm{r}_{\mathrm{i}}}^{\infty}(\Omega)}\|\cdot\|_{\mathrm{H}_{\mathrm{t}}^{\mathrm{m}}(\Omega)}, \quad \mathrm{R}_{\mathrm{t}, \Gamma_{\mathrm{i}}}(\Omega):=\overline{\mathrm{C}_{\Gamma_{\mathrm{i}}}^{\infty}(\Omega)}\|\cdot\|_{\mathrm{R}_{\mathrm{t}}(\Omega)}, \quad \mathrm{D}_{\mathrm{t}, \mathrm{\Gamma}_{\mathrm{i}}}(\Omega):=\overline{\mathrm{C}_{\Gamma_{\mathrm{i}}}^{\infty}(\Omega)}\|\cdot\|_{\mathrm{D}_{\mathrm{t}}(\Omega)} . \tag{2.2}
\end{align*}
$$

[^1]These spaces indeed generalize vanishing scalar, tangential and normal Dirichlet boundary conditions even and in particular to boundaries for which the notion of a normal vector may not make any sense. Moreover, 0 at the lower left corner denotes vanishing rotation resp. divergence, e.g.,

$$
{ }_{0} \mathrm{R}_{\mathrm{t}}(\Omega):=\left\{E \in \mathrm{R}_{\mathrm{t}}(\Omega) \mid \operatorname{rot} E=0\right\}, \quad{ }_{0} \mathrm{D}_{\mathrm{t}, \Gamma_{1}}(\Omega):=\left\{H \in \mathrm{D}_{\mathrm{t}, \Gamma_{1}}(\Omega) \mid \operatorname{div} H=0\right\},
$$

and if $t=0$ in any of the definitions given above, we will skip the weight, such that, e.g.,

$$
\mathbf{H}^{\mathrm{m}}(\Omega)=\mathbf{H}_{0}^{\mathrm{m}}(\Omega), \quad \mathrm{R}_{\Gamma_{1}}(\Omega)=\mathrm{R}_{0, \Gamma_{1}}(\Omega), \quad \mathbf{D}_{\Gamma_{1}}(\Omega)=\mathbf{D}_{0, \Gamma_{1}}(\Omega),
$$

Finally, in order to shorten notation, we declare

$$
\mathrm{X}_{<\mathrm{s}}:=\bigcap_{\mathrm{t}<\mathrm{s}} \mathrm{X}_{\mathrm{t}} \quad \text { and } \quad \mathrm{X}_{>\mathrm{s}}:=\bigcup_{\mathrm{t}>\mathrm{s}} \mathrm{X}_{\mathrm{t}} \quad(s \in \mathbb{R}),
$$

for $\mathrm{X}_{\mathrm{t}}$ being any of the spaces above and write, e.g.,

$$
\mathbf{H}_{\mathrm{t}}^{\mathrm{m}}=\mathbf{H}_{\mathrm{t}}^{\mathrm{m}}(\Omega), \quad \mathrm{R}_{\mathrm{t}, \Gamma_{1}}=\mathrm{R}_{\mathrm{t}, \Gamma_{1}}(\Omega), \quad \mathbf{D}_{\mathrm{t}}=\mathbf{D}_{\mathrm{t}}(\Omega), \quad \mathbf{H}_{\mathrm{t}, \Gamma_{2}}^{\mathrm{m}}=\mathrm{H}_{\mathrm{t}, \Gamma_{2}}^{\mathrm{m}}(\Omega)
$$

if $\Omega=\mathbb{R}^{3}$.

For the material parameters $\varepsilon$ and $\mu$ we suppose that they are $\kappa$ - admissible with $\kappa \geq 0$.
Definition 2.1. Let $\kappa \geq 0$. We call a transformation $\gamma \kappa$-admissible, if

- $\gamma: \Omega \longrightarrow \mathbb{R}^{3 \times 3}$ is an $\mathrm{L}^{\infty}$-matrix field,
- $\gamma$ is symmetric, i.e.,

$$
\forall E, H \in \mathrm{~L}^{2}(\Omega): \quad\langle E, \gamma H\rangle_{\mathrm{L}^{2}(\Omega)}=\langle\gamma E, H\rangle_{\mathrm{L}^{2}(\Omega)}
$$

- $\gamma$ is uniformly positive definite, i.e.,

$$
\exists c>0 \forall E \in \mathrm{~L}^{2}(\Omega): \quad\langle E, \gamma E\rangle_{\mathrm{L}^{2}(\Omega)} \geq c \cdot\|E\|_{\mathrm{L}^{2}(\Omega)}^{2},
$$

- $\gamma$ is asymptotically a multiple of the identity, i.e.,

$$
\gamma=\gamma_{0} \cdot \mathbb{1}+\hat{\gamma} \text { with } \gamma_{0} \in \mathbb{R}_{+} \text {and } \hat{\gamma}=\mathcal{O}\left(r^{-\kappa}\right) \text { as } r \longrightarrow \infty .
$$

Then $\varepsilon, \mu$ are pointwise invertible and $\varepsilon^{-1}, \mu^{-1}$ defined by

$$
\varepsilon^{-1}(x):=(\varepsilon(x))^{-1} \quad \text { and } \quad \mu^{-1}(x):=(\mu(x))^{-1}, \quad x \in \Omega
$$

are also $\kappa$-admissible. Moreover,

$$
\langle\cdot, \cdot\rangle_{\varepsilon}:=\langle\varepsilon \cdot, \cdot\rangle_{L^{2}(\Omega)} \quad \text { and } \quad\langle\cdot, \cdot\rangle_{\mu}:=\langle\mu \cdot, \cdot\rangle_{L^{2}(\Omega)}
$$

define scalar products on $L^{2}(\Omega)$ inducing norms equivalent to the standard ones. Consequently,

$$
\mathrm{L}_{\varepsilon}^{2}(\Omega):=\left(\mathrm{L}^{2}(\Omega),\langle\cdot, \cdot\rangle_{\varepsilon}\right), \quad \mathrm{L}_{\mu}^{2}(\Omega):=\left(\mathrm{L}^{2}(\Omega),\langle\cdot, \cdot\rangle_{\mu}\right) \quad \text { and } \quad \mathrm{L}_{\Lambda}^{2}(\Omega):=\mathrm{L}_{\varepsilon}^{2}(\Omega) \times \mathrm{L}_{\mu}^{2}(\Omega)
$$

are Hilbert spaces and we will write

$$
\|\cdot\|_{\varepsilon},\|\cdot\|_{\mu}\|\cdot\|_{\Lambda}, \quad \oplus_{\varepsilon}, \oplus_{\mu}, \oplus_{\Lambda} \quad \text { and } \quad \perp_{\varepsilon}, \perp_{\mu}, \perp_{\Lambda}
$$

to indicate the norm, the orthogonal sum and the orthogonal complement in these spaces. For further simplification and to shorten notation we also introduce for $\varepsilon=\varepsilon_{0} \cdot \mathbb{1}+\hat{\varepsilon}$ and $\mu=\mu_{0} \cdot \mathbb{1}+\hat{\mu}$ the formal matrix operators

$$
\begin{aligned}
& \Lambda:=\left(\begin{array}{cc}
\varepsilon & 0 \\
0 & \mu
\end{array}\right) \quad, \quad \Lambda^{-1}:=\left(\begin{array}{cc}
\varepsilon^{-1} & 0 \\
0 & \mu^{-1}
\end{array}\right) \quad, \quad \hat{\Lambda}:=\left(\begin{array}{cc}
\hat{\varepsilon} & 0 \\
0 & \hat{\mu}
\end{array}\right) \quad, \\
& \Lambda(E, H)=(\varepsilon E, \mu H) \quad \Lambda^{-1}(E, H)=\left(\varepsilon^{-1} E, \mu^{-1} H\right) \quad \hat{\Lambda}(E, H)=(\hat{\varepsilon} E, \hat{\mu} H) \\
& \operatorname{Rot}:=\left(\begin{array}{cc}
0 & -\operatorname{rot} \\
\operatorname{rot} & 0
\end{array}\right) \quad, \quad \mathrm{M}:=i \Lambda^{-1} \operatorname{Rot}=\left(\begin{array}{cc}
0 & -i \varepsilon^{-1} \operatorname{rot} \\
i \mu^{-1} \operatorname{rot} & 0
\end{array}\right), \\
& \operatorname{Rot}(E, H)=(-\operatorname{rot} H, \operatorname{rot} E) \quad \mathrm{M}(E, H)=\left(-i \varepsilon^{-1} \operatorname{rot} H, i \mu^{-1} \operatorname{rot} E\right)
\end{aligned}
$$

$$
\begin{array}{ccc}
\Xi:=\left(\begin{array}{cc}
0 & -\xi \times \\
\xi \times & 0
\end{array}\right) & \Lambda_{0}:=\left(\begin{array}{cc}
\varepsilon_{0} & 0 \\
0 & \mu_{0}
\end{array}\right) & \widetilde{\Lambda}_{0}:=\left(\begin{array}{cc}
\mu_{0} & 0 \\
0 & \varepsilon_{0}
\end{array}\right) \\
\Xi(E, H)=(-\xi \times H, \xi \times E) & & \Lambda_{0}(E, H)=\left(\varepsilon_{0} E, \mu_{0} H\right)
\end{array} \quad \widetilde{\Lambda}_{0}(E, H)=\left(\mu_{0} E, \varepsilon_{0} H\right) ~ l
$$

Recall $\xi(x)=x / r(x)$.
We end this section with a Lemma, showing that the spaces defined in (2.2) indeed generalize vanishing scalar, tangential and normal boundary conditions.
Lemma 2.2. For $t \in \mathbb{R}$ and $i \in\{1,2\}$ the following inclusions hold:
(a) $H_{\mathrm{t}, \Gamma_{\mathrm{i}}}^{\mathrm{m}}(\Omega) \subset \mathbf{H}_{\mathrm{t}, \Gamma_{\mathrm{i}}}^{\mathrm{m}}(\Omega), \quad \mathrm{R}_{\mathrm{t}, \Gamma_{\mathrm{i}}}(\Omega) \subset \mathbf{R}_{\mathrm{t}, \Gamma_{\mathrm{i}}}(\Omega), \quad \mathrm{D}_{\mathrm{t}, \Gamma_{\mathrm{i}}}(\Omega) \subset \mathbf{D}_{\mathrm{t}, \Gamma_{\mathrm{i}}}(\Omega)$
(b) $\nabla \mathbf{H}_{\mathrm{t}, \Gamma_{\mathrm{i}}}^{1}(\Omega) \subset{ }_{0} \mathrm{R}_{\mathrm{t}, \Gamma_{\mathrm{i}}}(\Omega), \quad \nabla \mathrm{H}_{\mathrm{t}, \Gamma_{\mathrm{i}}}^{1}(\Omega) \subset{ }_{0} \mathrm{R}_{\mathrm{t}+1, \Gamma_{\mathrm{i}}}(\Omega)$
(c) $\operatorname{rot} \mathbf{R}_{\mathrm{t}, \Gamma_{\mathrm{i}}}(\Omega) \subset{ }_{0} \mathrm{D}_{\mathrm{t}, \Gamma_{\mathrm{i}}}(\Omega), \quad \operatorname{rot} \mathrm{R}_{\mathrm{t}, \Gamma_{\mathrm{i}}}(\Omega) \subset{ }_{0} \mathrm{D}_{\mathrm{t}+1, \Gamma_{\mathrm{i}}}(\Omega)$

Additionally we have for $i, j \in\{1,2\}, i \neq j$ :

$$
\begin{aligned}
& \mathbf{H}_{\mathrm{t}, \Gamma_{\mathrm{i}}}^{1}(\Omega)=\mathcal{H}_{\mathrm{t}, \Gamma_{\mathrm{i}}}^{1}(\Omega):=\left\{w \in \mathbf{H}_{\mathrm{t}}^{1}(\Omega) \mid \forall \Phi \in \mathrm{C}_{\Gamma_{\mathrm{j}}}^{\infty}(\Omega):\langle w, \operatorname{div} \Phi\rangle_{\mathrm{L}^{2}(\Omega)}=-\langle\nabla w, \Phi\rangle_{\mathrm{L}^{2}(\Omega)}\right\}, \\
& \mathbf{R}_{\mathrm{t}, \Gamma_{\mathrm{i}}}(\Omega)=\mathcal{R}_{\mathrm{t}, \Gamma_{\mathrm{i}}}(\Omega):=\left\{E \in \mathbf{R}_{\mathrm{t}}(\Omega) \mid \forall \Phi \in \mathrm{C}_{\Gamma_{\mathrm{j}}}^{\infty}(\Omega):\langle E, \operatorname{rot} \Phi\rangle_{\mathrm{L}^{2}(\Omega)}=\langle\operatorname{rot} E, \Phi\rangle_{\mathrm{L}^{2}(\Omega)}\right\}, \\
& \mathbf{D}_{\mathrm{t}, \Gamma_{\mathrm{i}}}(\Omega)=\mathcal{D}_{\mathrm{t}, \Gamma_{\mathrm{i}}}(\Omega):=\left\{H \in \mathbf{D}_{\mathrm{t}}(\Omega) \mid \forall \phi \in \mathrm{C}_{\Gamma_{\mathrm{j}}}^{\infty}(\Omega):\langle H, \nabla \phi\rangle_{\mathrm{L}^{2}(\Omega)}=-\langle\operatorname{div} H, \phi\rangle_{\mathrm{L}^{2}(\Omega)}\right\},
\end{aligned}
$$

and

$$
\begin{aligned}
& \mathrm{H}_{\mathrm{t}, \Gamma_{\mathrm{i}}}^{1}(\Omega)=\mathcal{H}_{\mathrm{t}, \Gamma_{\mathrm{i}}}^{1}(\Omega):=\left\{w \in \mathrm{H}_{\mathrm{t}}^{1}(\Omega) \mid \forall \Phi \in \mathrm{C}_{\Gamma_{\mathrm{j}}}^{\infty}(\Omega):\langle w, \operatorname{div} \Phi\rangle_{\mathrm{L}^{2}(\Omega)}=-\langle\nabla w, \Phi\rangle_{\mathrm{L}^{2}(\Omega)}\right\}, \\
& \mathrm{R}_{\mathrm{t}, \Gamma_{\mathrm{i}}}(\Omega)=\mathcal{R}_{\mathrm{t}, \Gamma_{\mathrm{i}}}(\Omega):=\left\{E \in \mathrm{R}_{\mathrm{t}}(\Omega) \mid \forall \Phi \in \mathrm{C}_{\Gamma_{\mathrm{j}}}^{\infty}(\Omega):\langle E, \operatorname{rot} \Phi\rangle_{\mathrm{L}^{2}(\Omega)}=\langle\operatorname{rot} E, \Phi\rangle_{\mathrm{L}^{2}(\Omega)}\right\}, \\
& \mathrm{D}_{\mathrm{t}, \Gamma_{\mathrm{i}}}(\Omega)=\mathcal{D}_{\mathrm{t}, \Gamma_{\mathrm{i}}}(\Omega):=\left\{H \in \mathrm{D}_{\mathrm{t}}(\Omega) \mid \forall \phi \in \mathrm{C}_{\Gamma_{\mathrm{j}}}^{\infty}(\Omega):\langle H, \nabla \phi\rangle_{\mathrm{L}^{2}(\Omega)}=-\langle\operatorname{div} H, \phi\rangle_{\mathrm{L}^{2}(\Omega)}\right\},
\end{aligned}
$$

where ( using continuity of the $\mathrm{L}^{2}$-scalar product) we may also replace $\mathrm{C}_{\Gamma_{j}}^{\infty}(\Omega)$ by

$$
\mathbf{H}_{\mathrm{s}, \Gamma_{\mathrm{j}}}^{1}(\Omega), \mathbf{R}_{\mathrm{s}, \Gamma_{j}}(\Omega), \mathbf{D}_{\mathrm{s}, \Gamma_{\mathrm{j}}}(\Omega) \quad \text { resp. } \quad \mathbf{H}_{\mathrm{s}, \Gamma_{\mathrm{j}}}^{1}(\Omega), \mathrm{R}_{\mathrm{s}, \Gamma_{\mathrm{j}}}(\Omega), \mathrm{D}_{\mathrm{s}, \Gamma_{\mathrm{j}}}(\Omega),
$$

with $s+t \geq 0$ resp. $s+t \geq-1$.
Proof. As representatives of the arguments we show

$$
\text { (i) } \operatorname{rot} \mathbf{R}_{\mathrm{t}, \Gamma_{2}}(\Omega) \subset{ }_{0} \mathrm{D}_{\mathrm{t}, \Gamma_{2}}(\Omega) \quad \text { and } \quad \text { (ii) } \mathbf{R}_{\mathrm{t}, \Gamma_{1}}(\Omega)=\boldsymbol{\mathcal { R }}_{\mathrm{t}, \Gamma_{1}}(\Omega) .
$$

For $E \in \operatorname{rot} \mathbf{R}_{\mathrm{t}, \Gamma_{2}}(\Omega)$ there exists a sequence $\left(\mathcal{E}_{n}\right)_{n \in \mathbb{N}} \subset \mathrm{C}_{\Gamma_{2}}^{\infty}(\Omega)$ such that $\operatorname{rot} \mathcal{E}_{n} \longrightarrow E$ in $\mathrm{L}_{\mathrm{t}}^{2}(\Omega)$. Then

$$
\forall \phi \in \dot{C}^{\infty}(\Omega): \quad\langle E, \nabla \phi\rangle_{\mathrm{L}^{2}(\Omega)}=\lim _{n \rightarrow \infty}\left\langle\operatorname{rot} \mathcal{E}_{n}, \nabla \phi\right\rangle_{\mathrm{L}^{2}(\Omega)}=-\lim _{n \rightarrow \infty}\left\langle\operatorname{div}\left(\operatorname{rot} \mathcal{E}_{n}\right), \phi\right\rangle_{\mathrm{L}^{2}(\Omega)}=0,
$$

hence $E$ has vanishing divergence and $\left(E_{n}\right)_{n \in \mathbb{N}}$ defined by $E_{n}:=\operatorname{rot} \mathcal{E}_{n}$ satisfies

$$
\left(E_{n}\right)_{n \in \mathbb{N}} \subset C_{\Gamma_{2}}^{\infty}(\Omega), \quad E_{n} \xrightarrow{\mathrm{~L}_{\mathrm{t}}^{2}(\Omega)} E \quad \text { and } \quad \operatorname{div} E_{n}=\operatorname{div}\left(\operatorname{rot} \mathcal{E}_{n}\right)=0 \xrightarrow{\mathrm{~L}_{\mathrm{t}+1}^{2}(\Omega)} 0 .
$$

Thus $E \in{ }_{0} \mathrm{D}_{\mathrm{t}, \Gamma_{2}}(\Omega)$ and we have shown (i). Let us show (ii). We clearly have $\mathbf{R}_{\mathrm{t}, \Gamma_{1}}(\Omega) \subset \boldsymbol{\mathcal { R }}_{\mathrm{t}, \Gamma_{1}}(\Omega)$. For the other direction, let $E \in \mathcal{R}_{\mathrm{t}, \Gamma_{1}}(\Omega)$ and $\delta>0$. Using the cut-off function from above we define $\left(E_{k}\right)_{k \in \mathbb{N}}$ by $E_{k}:=\eta_{k} E$. Then $E_{k} \in \mathcal{R}_{\tilde{\Gamma}_{1}}\left(\Omega\left(2 r_{k}\right)\right), \tilde{\Gamma}_{1}:=\Gamma_{1} \cup \mathrm{~S}\left(2 r_{k}\right)$, since for $\Phi \in \mathrm{C}_{\Gamma_{2}}^{\infty}\left(\Omega\left(2 r_{k}\right)\right)$ it holds by $\eta_{k} \Phi \in \mathrm{C}_{\Gamma_{2}}^{\infty}(\Omega)$

$$
\begin{aligned}
\left\langle E_{k}, \operatorname{rot} \Phi\right\rangle_{\mathrm{L}^{2}\left(\Omega\left(2 r_{k}\right)\right)} & =\left\langle\eta_{k} E, \operatorname{rot} \Phi\right\rangle_{\mathrm{L}^{2}(\Omega)} \\
& =\left\langle E, \operatorname{rot}\left(\eta_{k} \Phi\right)\right\rangle_{\mathrm{L}^{2}(\Omega)}-\left\langle E, \nabla \eta_{k} \times \Phi\right\rangle_{\mathrm{L}^{2}(\Omega)} \\
& =\left\langle\eta_{k} \operatorname{rot} E+\nabla \eta_{k} \times E, \Phi\right\rangle_{\mathrm{L}^{2}\left(\Omega\left(2 r_{k}\right)\right.}=\left\langle\operatorname{rot} E_{k}, \Phi\right\rangle_{\mathrm{L}^{2}\left(\Omega\left(2 r_{k}\right)\right)} .
\end{aligned}
$$

Now observe, that by means of monotone convergence we have ${ }^{\text {ii }}$

$$
\left\|E-E_{k}\right\|_{\mathbf{R}_{\mathrm{t}}(\Omega)}=\left\|\check{\eta}_{k} E\right\|_{\mathbf{R}_{\mathrm{t}}(\Omega)} \leq c \cdot\left(\|E\|_{\mathbf{R}_{\mathrm{t}}\left(\check{\mathrm{U}}\left(r_{k}\right)\right)}+\frac{1}{2^{k}} \cdot\|E\|_{\mathrm{L}_{\mathrm{t}}^{2}(\Omega)}\right) \longrightarrow 0
$$

hence we can choose $\hat{k}>0$ such that $\left\|E-E_{\hat{k}}\right\|_{\mathbf{R}_{\mathrm{t}}(\Omega)}<\delta / 2$. Moreover $\Omega\left(2 r_{\hat{k}}\right)=\Omega \cap \mathrm{U}\left(2 r_{\hat{k}}\right)$ is a bounded weak Lipschitz domain and therefore (cf. [1, Section 3.3]) $\boldsymbol{\mathcal { R }}_{\tilde{\Gamma}_{1}}\left(\Omega\left(2 r_{\hat{k}}\right)\right)=\mathbf{R}_{\tilde{\Gamma}_{1}}\left(\Omega\left(2 r_{\hat{k}}\right)\right)$, yielding the existence of some $\Psi \in \mathrm{C}_{\tilde{\Gamma}_{1}}^{\infty}\left(\Omega\left(2 r_{\hat{k}}\right)\right)$ such that

$$
\left\|E_{\hat{k}}-\Psi\right\|_{\mathbf{R}_{\mathbf{t}}\left(\Omega\left(2 r_{\hat{k}}\right)\right)} \leq c \cdot\left\|E_{\hat{k}}-\Psi\right\|_{\mathbf{R}\left(\Omega\left(2 r_{\hat{k}}\right)\right)}<\delta / 2
$$

Extending $\Psi$ by zero to $\Omega$ we obtain (by abuse of notation) $\Psi \in C_{\Gamma_{1}}^{\infty}(\Omega)$ with

$$
\|E-\Psi\|_{\mathbf{R}_{\mathrm{t}}(\Omega)} \leq\left\|E-E_{\hat{k}}\right\|_{\mathbf{R}_{\mathrm{t}}(\Omega)}+\left\|E_{\hat{k}}-\Psi\right\|_{\mathbf{R}_{\mathrm{t}}\left(\Omega\left(2 r_{\hat{k}}\right)\right)}<\frac{\delta}{2}+\frac{\delta}{2}=\delta
$$

which completes the proof.
2.2. Some Functional Analysis. Let $H_{1}$ and $H_{2}$ be Hilbert spaces. With $L\left(H_{1}, H_{2}\right)$ and $B\left(H_{1}, H_{2}\right)$ we introduce the sets of linear resp. bounded linear operators mapping $\mathrm{H}_{1}$ to $\mathrm{H}_{2}$. For $\mathrm{A}: \mathcal{D}(\mathrm{A}) \subset \mathrm{H}_{1} \longrightarrow \mathrm{H}_{2}$ linear, closed, and densely defined, the adjoint $A^{*}: \mathcal{D}\left(A^{*}\right) \subset H_{2} \longrightarrow H_{1}$ is characterized by

$$
\langle\mathrm{A} x, y\rangle_{\mathrm{H}_{2}}=\left\langle x, \mathrm{~A}^{*} y\right\rangle_{\mathrm{H}_{1}} \quad \forall x \in \mathcal{D}(\mathrm{~A}), y \in \mathcal{D}\left(\mathrm{~A}^{*}\right) .
$$

By the projection theorem we have the following Helmholtz type decompositions

$$
\mathrm{H}_{1}=\overline{\mathcal{R}\left(\mathrm{A}^{*}\right)} \oplus \mathcal{N}(\mathrm{A}), \quad \text { and } \quad \mathrm{H}_{2}=\overline{\mathcal{R}(\mathrm{A})} \oplus \mathcal{N}\left(\mathrm{A}^{*}\right)
$$

which we use to define the corresponding reduced operators $\mathcal{A}:=\left.\mathrm{A}\right|_{\mathcal{N}(\mathrm{A})^{\perp}}, \mathcal{A}^{*}:=\left.\mathrm{A}^{*}\right|_{\mathcal{N}\left(\mathrm{A}^{*}\right)^{\perp}}$, i.e.,

$$
\begin{array}{rlrl}
\mathcal{A}: \mathcal{D}(\mathcal{A}) & \subset \overline{\mathcal{R}\left(\mathrm{A}^{*}\right)} \longrightarrow \overline{\mathcal{R}(\mathrm{A})} & \text { resp. } & \mathcal{A}^{*}: \mathcal{D}\left(\mathcal{A}^{*}\right) \subset \overline{\mathcal{R}(\mathrm{A})} \longrightarrow \overline{\mathcal{R}\left(\mathrm{A}^{*}\right)} \\
\mathcal{D}(\mathcal{A})=\mathcal{D}(\mathrm{A}) \cap \overline{\mathcal{R}\left(\mathrm{A}^{*}\right)} & & \mathcal{D}\left(\mathcal{A}^{*}\right)=\mathcal{D}\left(\mathrm{A}^{*}\right) \cap \overline{\mathcal{R}(\mathrm{A})}
\end{array}
$$

These operators are also closed, densely defined and indeed adjoint to each other. Moreover, by definition $\mathcal{A}$ and $\mathcal{A}^{*}$ are injective and therefore the inverse operators

$$
\mathcal{A}^{-1}: \mathcal{R}(\mathrm{A}) \longrightarrow \mathcal{D}(\mathcal{A}) \quad \text { and } \quad\left(\mathcal{A}^{*}\right)^{-1}: \mathcal{R}\left(\mathrm{A}^{*}\right) \longrightarrow \mathcal{D}\left(\mathcal{A}^{*}\right)
$$

exist. Moreover, the pair $\left(\mathcal{A}, \mathcal{A}^{*}\right)$ satisfies the following result of the so called Functional Analysis Toolbox, see e.g. [17, Section 2], from which we will derive some Poincaré type estimates for the time-harmonic Maxwell operator $(\mathcal{M}-\omega)$ (cf. Remark 3.11 and Remark 3.7).
Lemma 2.3. The following assertions are equivalent:
(1) $\exists c_{\mathrm{A}} \in(0, \infty) \quad \forall x \in \mathcal{D}(\mathcal{A}): \quad\|x\|_{\mathrm{H}_{1}} \leq c_{\mathrm{A}}\|\mathrm{A} x\|_{\mathrm{H}_{2}}$.
$\left(1^{*}\right) \exists c_{\mathrm{A}^{*}} \in(0, \infty) \quad \forall y \in \mathcal{D}\left(\mathcal{A}^{*}\right): \quad\|y\|_{\mathrm{H}_{2}} \leq c_{\mathrm{A}^{*}}\left\|\mathrm{~A}^{*} y\right\|_{\mathrm{H}_{1}}$.
(2) $\mathcal{R}(\mathrm{A})=\mathcal{R}(\mathcal{A})$ is closed in $\mathrm{H}_{2}$.
(2*) $\mathcal{R}\left(\mathrm{A}^{*}\right)=\mathcal{R}\left(\mathcal{A}^{*}\right)$ ist closed in $\mathrm{H}_{1}$.
(3) $\mathcal{A}^{-1}: \mathcal{R}(\mathrm{A}) \longrightarrow \mathcal{D}(\mathcal{A})$ is continuous.
$\left(3^{*}\right)\left(\mathcal{A}^{*}\right)^{-1}: \mathcal{R}\left(\mathrm{A}^{*}\right) \longrightarrow \mathcal{D}\left(\mathcal{A}^{*}\right)$ is continuous.
Note that

$$
\left\|\mathcal{A}^{-1}\right\|_{\mathcal{R}(\mathrm{A}), \mathcal{R}\left(\mathrm{A}^{*}\right)}=c_{\mathrm{A}}=c_{\mathrm{A}^{*}}=\left\|\left(\mathcal{A}^{*}\right)^{-1}\right\|_{\mathcal{R}\left(\mathrm{A}^{*}\right), \mathcal{R}(\mathrm{A})} .
$$

[^2]
## 3. Solution Theory for Time-Harmonic Maxwell Equations

As mentioned above we treat the time-harmonic Maxwell equations with mixed boundary conditions

$$
\begin{array}{rll}
-\operatorname{rot} H+i \omega \varepsilon E=F & \text { in } \Omega, & E \times \nu=0 \quad \text { on } \Gamma_{1}, \\
\operatorname{rot} E+i \omega \mu H=G & \text { in } \Omega, & H \times \nu=0 \quad \text { on } \Gamma_{2}, \tag{3.1}
\end{array}
$$

in an exterior weak Lipschitz domain $\Omega \subset \mathbb{R}^{3}$ and for frequencies $\omega \in \mathbb{C} \backslash\{0\}$. Moreover, we suppose that the material parameters $\varepsilon$ and $\mu$ are $\kappa$-admissible for $\kappa \geq 0$. Using the abbreviations from above and rewriting

$$
u:=(E, H) \quad, \quad f:=i \Lambda^{-1}(-F, G),
$$

the weak formulation of these boundary value problem reads:

$$
\begin{equation*}
\text { For } f \in \mathrm{~L}_{\mathrm{loc}}^{2}(\bar{\Omega}) \times \mathrm{L}_{\mathrm{loc}}^{2}(\bar{\Omega}) \text { find } u \in \mathrm{R}_{\mathrm{loc}, \Gamma_{1}}(\bar{\Omega}) \times \mathrm{R}_{\mathrm{loc}}(\bar{\Omega}) \text { such that }(\mathrm{M}-\omega) u=f . \tag{3.2}
\end{equation*}
$$

We shall solve this problem using polynomially weighted Hilbert spaces. In doing so we avoid additional assumptions on boundary regularity for $\Omega$, since only a compactness result comparable to Rellich's selection theorem is needed. More precisely, we will show that $\Omega$ satisfies "Weck's (local) selection theorem", also called "(local) Maxwell compactness property", which in fact is also an assumption on the quality of the boundary and in some sense supersedes assumptions on boundary regularity.

Definition 3.1. Let $\gamma$ be $\kappa$-admissible with $\kappa \geq 0 . \Omega$ satisfies "Weck's local selection theorem" (WLST) (or has the "local Maxwell compactness property"), if the embedding

$$
\begin{equation*}
\mathbf{R}_{\Gamma_{1}}(\Omega) \cap \gamma^{-1} \mathbf{D}_{\Gamma_{2}}(\Omega) \longleftrightarrow \mathrm{L}_{\mathrm{loc}}^{2}(\bar{\Omega}) \tag{3.3}
\end{equation*}
$$

is compact. A bounded weak Lipschitz domain $\Theta$ with weak Lipschitz Interface satisfies "Weck's selection theorem" (WST) ( or posses the "Maxwell compactness property") if the embedding

$$
\begin{equation*}
\mathbf{R}_{\Gamma_{1}}(\Theta) \cap \gamma^{-1} \mathbf{D}_{\Gamma_{2}}(\Theta) \longleftrightarrow \mathrm{L}^{2}(\Theta) \tag{3.4}
\end{equation*}
$$

is compact.
Remark 3.2. Note that Weck's (local) selection theorem is essentially independent of $\gamma$ meaning that a domain $\Omega \subset \mathbb{R}^{3}$ satisfies WST resp. WLST, if and only if the imbedding

$$
\mathbf{R}_{\Gamma_{1}}(\Omega) \cap \mathbf{D}_{\Gamma_{2}}(\Omega) \longleftrightarrow \mathrm{L}^{2}(\Omega) \quad \text { resp. } \quad \mathbf{R}_{\Gamma_{1}}(\Omega) \cap \mathbf{D}_{\Gamma_{2}}(\Omega) \longleftrightarrow \mathrm{L}_{\mathrm{loc}}^{2}(\bar{\Omega})
$$

is compact. The proof is practically identical with the one of [19, Lemma 2] (see also [24], [22]) .
Lemma 3.3. Let $\gamma$ be $\kappa$-admissible with $\kappa \geq 0$. Then the following statements are equivalent:
(a) $\Omega$ satisfies WLST
(b) For all $\tilde{r}>r_{0}$ the imbedding

$$
\mathbf{R}_{\tilde{\Gamma}_{1}}(\Omega(\tilde{r})) \cap \gamma^{-1} \mathbf{D}_{\Gamma_{2}}(\Omega(\tilde{r})) \longleftrightarrow \mathrm{L}^{2}(\Omega(\tilde{r}))
$$

with $\tilde{\Gamma}_{1}:=\Gamma_{1} \cup \mathrm{~S}(\tilde{r})$ is compact.
(c) For all $\tilde{r}>r_{0}$ the imbedding

$$
\mathbf{R}_{\Gamma_{1}}(\Omega(\tilde{r})) \cap \gamma^{-1} \mathbf{D}_{\tilde{\Gamma}_{2}}(\Omega(\tilde{r})) \longleftrightarrow \mathrm{L}^{2}(\Omega(\tilde{r}))
$$

with $\tilde{\Gamma}_{2}:=\Gamma_{2} \cup \mathrm{~S}(\tilde{r})$ is compact.
(d) For all $s, t \in \mathbb{R}$ with $t<s$ the imbedding

$$
\mathbf{R}_{\mathrm{s}, \Gamma_{1}}(\Omega) \cap \gamma^{-1} \mathbf{D}_{\mathrm{s}, \Gamma_{2}}(\Omega) \longleftrightarrow \mathrm{L}_{\mathrm{t}}^{2}(\Omega)
$$

is compact.
Proof.
$(a) \Longrightarrow(b)$ :
Let $\tilde{r}>r_{0}$. By Remark 3.2 it is sufficient to show the compactness of

$$
\mathbf{R}_{\tilde{\Gamma}_{1}}(\Omega(\tilde{r})) \cap \mathbf{D}_{\Gamma_{2}}(\Omega(\tilde{r})) \longleftrightarrow \mathrm{L}^{2}(\Omega(\tilde{r})) .
$$

Therefore let $\left(E_{n}\right)_{n \in \mathbb{N}} \subset \mathbf{R}_{\tilde{\Gamma}_{1}}(\Omega(\tilde{r})) \cap \mathbf{D}_{\Gamma_{2}}(\Omega(\tilde{r}))$ be bounded, choose $r_{0}<\hat{r}<\tilde{r}$ and a cut-off function $\chi \in \mathrm{C}^{\infty}\left(\mathbb{R}^{3}\right)$ with $\operatorname{supp} \chi \subset \mathrm{U}(\tilde{r})$ and $\left.\chi\right|_{\mathrm{B}(\hat{r})}=1$. Then, for every $n \in \mathbb{N}$ we have

$$
E_{n}=\check{E}_{n}+\hat{E}_{n}:=\chi E_{n}+(1-\chi) E_{n}, \quad \operatorname{supp} \check{E}_{n} \subset \Omega(\tilde{r}), \quad \operatorname{supp} \hat{E}_{n} \subset \mathrm{G}\left(r_{0}, \tilde{r}\right)
$$

splitting $\left(E_{n}\right)_{n \in \mathbb{N}}$ into $\left(\check{E}_{n}\right)_{n \in \mathbb{N}}$ and $\left(\hat{E}_{n}\right)_{n \in \mathbb{N}}$. Extending $\check{E}_{n}$ resp. $\hat{E}_{n}$ by zero, we obtain (by abuse of notation) sequences

$$
\left(\check{E}_{n}\right)_{n \in \mathbb{N}} \subset \mathbf{R}_{\Gamma_{1}}(\Omega) \cap \mathbf{D}_{\Gamma_{2}}(\Omega) \quad \text { and } \quad\left(\hat{E}_{n}\right)_{n \in \mathbb{N}} \subset \mathbf{R}_{\mathrm{S}(\tilde{\tilde{r})}}(\mathrm{B}(\tilde{r})) \cap \mathbf{D}(\mathrm{B}(\tilde{r}))
$$

which are bounded in the respective spaces. Thus, using Weck's local selection theorem and Remark 3.2, we can choose a subsequence $\left(\check{E}_{\pi(n)}\right)_{n \in \mathbb{N}}$ of $\left(\check{E}_{n}\right)_{n \in \mathbb{N}}$ converging in $\mathrm{L}_{\text {loc }}^{2}(\bar{\Omega})$. The corresponding subsequence $\left(\hat{E}_{\pi(n)}\right)_{n \in \mathbb{N}}$ is of course also bounded in $\mathbf{R}_{\mathrm{S}(\tilde{\mathrm{r}})}(\mathrm{B}(\tilde{r})) \cap \mathbf{D}(\mathrm{B}(\tilde{r}))$ and by [23, Theorem 2.2], even in $\mathbf{H}^{1}(\mathrm{~B}(\tilde{r}))$, hence (Rellich's selection theorem) has a subsequence $\left(\hat{E}_{\tilde{\pi}(n)}\right)_{n \in \mathbb{N}}$ converging in $\mathrm{L}^{2}(\mathrm{~B}(\tilde{r}))$. Thus

$$
\begin{aligned}
\| E_{\tilde{\pi}(n)}- & E_{\tilde{\pi}(m)} \|_{\mathrm{L}^{2}(\Omega(\tilde{r}))} \\
\leq c \cdot & \left\{\left\|\chi\left(E_{\tilde{\pi}(n)}-E_{\tilde{\pi}(m)}\right)\right\|_{\mathrm{L}^{2}(\Omega(\tilde{r}))}+\left\|(1-\chi)\left(E_{\tilde{\pi}(n)}-E_{\tilde{\pi}(m)}\right)\right\|_{\mathrm{L}^{2}(\Omega(\tilde{r}))}\right\} \\
& \leq c \cdot\left\{\left\|\check{E}_{\tilde{\pi}(n)}-\check{E}_{\tilde{\pi}(m)}\right\|_{\mathrm{L}^{2}(\Omega(\tilde{r}))}+\left\|\hat{E}_{\tilde{\pi}(n)}-\hat{E}_{\tilde{\pi}(m)}\right\|_{\mathrm{L}^{2}(\mathrm{~B}(\tilde{r}))}\right\} \xrightarrow{m, n \rightarrow \infty} 0
\end{aligned}
$$

meaning $\left(E_{\tilde{\pi}(n)}\right)_{n \in \mathbb{N}} \subset\left(E_{n}\right)_{n \in \mathbb{N}}$ is a Cauchy sequence and therefore converging in $\mathrm{L}^{2}(\Omega(\tilde{r}))$.
$(b) \Longrightarrow(d)$ :
Let $s, t \in \mathbb{R}$ with $s>t$ and $\left(E_{n}\right)_{n \in \mathbb{N}} \subset \mathbf{R}_{\mathrm{s}, \Gamma_{1}}(\Omega) \cap \gamma^{-1} \mathbf{D}_{\mathrm{s}, \Gamma_{2}}(\Omega)$ be bounded. Then there exists a subsequence $\left(E_{\pi(n)}\right)_{n \in \mathbb{N}} \subset\left(E_{n}\right)_{n \in \mathbb{N}}$ converging weakly in $\mathbf{R}_{\mathrm{s}, \Gamma_{1}}(\Omega) \cap \gamma^{-1} \mathbf{D}_{\mathrm{s}, \Gamma_{2}}(\Omega)$ to some vector field $E \in \mathbf{R}_{\mathrm{s}, \Gamma_{1}}(\Omega) \cap \gamma^{-1} \mathbf{D}_{\mathrm{s}, \Gamma_{2}}(\Omega)$. We now construct a subsequence $\left(E_{\tilde{\pi}(n)}\right)_{n \in \mathbb{N}}$ of $\left(E_{\pi(n)}\right)_{n \in \mathbb{N}}$ converging in $\mathrm{L}_{\text {loc }}^{2}(\bar{\Omega})$ to the same limit $E$. Therefore, observe that

$$
\left(E_{\pi(n), 1}\right)_{n \in \mathbb{N}} \quad \text { with } \quad E_{\pi(n), 1}:=\eta_{1} E_{\pi(n)}
$$

is bounded in $\mathbf{R}_{\tilde{\Gamma}_{1}}\left(\Omega\left(r_{2}\right) \cap \gamma^{-1} \mathbf{D}_{\Gamma_{2}}\left(\Omega\left(r_{2}\right)\right), \tilde{\Gamma}_{1}:=\Gamma_{1} \cup \mathrm{~S}\left(r_{2}\right)\right.$ such that by assumption there exists a subsequence $\left(E_{\pi_{1}(n), 1}\right)_{n \in \mathbb{N}}$ converging in $\mathrm{L}^{2}\left(\Omega\left(r_{2}\right)\right)$. Then $\left(E_{\pi_{1}(n)}\right)_{n \in \mathbb{N}} \subset\left(E_{\pi(n)}\right)_{n \in \mathbb{N}}$ is converging in $\mathrm{L}^{2}\left(\Omega\left(r_{1}\right)\right)$ and as $\left(E_{\pi_{1}(n)}\right)_{n \in \mathbb{N}}$ is also weakly convergent in $\mathrm{L}^{2}\left(\Omega\left(r_{1}\right)\right)$, we have

$$
E_{\pi_{1}(n)} \longrightarrow E \quad \text { in } \quad \mathrm{L}^{2}\left(\Omega\left(r_{1}\right)\right)
$$

Multiplying $\left(E_{\pi_{1}(n)}\right)_{n \in \mathbb{N}}$ with $\eta_{2}$ we obtain a sequence $\left(E_{\pi_{1}(n), 2}\right)_{n \in \mathbb{N}}, E_{\pi_{1}(n), 2}=\eta_{2} E_{\pi_{1}(n)}$ bounded in $\mathbf{R}_{\tilde{\Gamma}_{1}}\left(\Omega\left(r_{3}\right)\right) \cap \gamma^{-1} \mathbf{D}_{\Gamma_{2}}\left(\Omega\left(r_{3}\right)\right), \tilde{\Gamma}_{1}:=\Gamma_{1} \cup \mathrm{~S}\left(r_{3}\right)$ and as before we construct a subsequence $\left(E_{\pi_{2}(n), 2}\right)_{n \in \mathbb{N}}$ converging in $\mathrm{L}^{2}\left(\Omega\left(r_{3}\right)\right)$, giving a subsequence $\left(E_{\pi_{2}(n)}\right)_{n \in \mathbb{N}} \subset\left(E_{\pi_{1}(n)}\right)_{n \in \mathbb{N}}$ with

$$
E_{\pi_{2}(n)} \longrightarrow E \quad \text { in } \quad \mathrm{L}^{2}\left(\Omega\left(r_{2}\right)\right)
$$

Continuing like this, we successively construct subsequences $\left(E_{\pi_{k}(n)}\right)_{n \in \mathbb{N}}$ with $E_{\pi_{k}(n)} \longrightarrow E$ in $\mathrm{L}^{2}\left(\Omega\left(r_{k}\right)\right)$ and switching to the diagonal sequence we indeed end up with a sequence $\left(E_{\tilde{\pi}(n)}\right)_{n \in \mathbb{N}}$, $\tilde{\pi}(n):=\pi_{n}(n)$, with $E_{\tilde{\pi}(n)} \longrightarrow E$ in $\mathrm{L}_{\text {loc }}^{2}(\bar{\Omega})$. Now Lemma A. 1 implies for arbitrary $\theta>0$

$$
\left\|E_{\tilde{\pi}(n)}-E\right\|_{\mathrm{L}_{t}^{2}(\Omega)} \leq c \cdot\left\|E_{\tilde{\pi}(n)}-E\right\|_{\mathrm{L}^{2}(\Omega(\delta))}+\theta
$$

with $c, \delta \in(0, \infty)$ independent of $E_{\tilde{\pi}(n)}$. Hence

$$
\limsup _{n \rightarrow \infty}\left\|E_{\tilde{\pi}(n)}-E\right\|_{L_{t}^{2}(\Omega)} \leq \theta
$$

and we obtain $E_{\tilde{\pi}(n)} \longrightarrow E$ in $\mathrm{L}_{\mathrm{t}}^{2}(\Omega)$.
$(d) \Longrightarrow(a):$
For $\left(E_{n}\right)_{n \in \mathbb{N}}$ bounded in $\mathbf{R}_{\Gamma_{1}}(\Omega) \cap \gamma^{-1} \mathbf{D}_{\Gamma_{2}}(\Omega)$ assertion (c) implies the existence of a subsequence $\left(E_{\pi(n)}\right)_{n \in \mathbb{N}}$ converging in $\mathrm{L}_{-1}^{2}(\Omega)$ to some $E \in \mathrm{~L}_{-1}^{2}(\Omega)$. Then $E \in \mathrm{~L}_{\text {loc }}^{2}(\bar{\Omega})$ and as

$$
\forall \tilde{r}>0: \quad\left\|E_{\pi(n)}-E\right\|_{\mathrm{L}^{2}(\Omega(\tilde{r}))} \leq(1+\tilde{r})^{1 / 2} \cdot\left\|E_{\pi(n)}-E\right\|_{\mathrm{L}_{-1}^{2}(\Omega)}
$$

we obtain $\left(E_{\pi(n)}\right)_{n \in \mathbb{N}} \longrightarrow E$ in $\mathrm{L}_{\mathrm{loc}}^{2}(\bar{\Omega})$.
Similar arguments to those corresponding to (b) show the assertion for $(c)$.

As shown by Bauer, Pauly \& Schomburg [1, Theorem 4.7], bounded weak Lipschitz domains satisfy Weck's selection theorem and by Lemma 3.3 (a) this directly implies:

Theorem 3.4. Exterior weak Lipschitz domains satisfy Weck's local selection theorem.
A first step to a solution theory for (3.2) with arbitrary $\omega \in \mathbb{C} \backslash\{0\}$ is the following result
Theorem 3.5. The Maxwell operator

$$
\mathcal{M}: \mathbf{R}_{\Gamma_{1}}(\Omega) \times \mathbf{R}_{\Gamma_{2}}(\Omega) \subset \mathrm{L}_{\Lambda}^{2}(\Omega) \longrightarrow \mathrm{L}_{\Lambda}^{2}(\Omega), u \longmapsto \mathrm{M} u,
$$

is self-adjoint and reduced by the closure of its range

$$
\overline{\mathcal{R}(\mathcal{M})}=\varepsilon^{-1} \overline{\operatorname{rot} \mathbf{R}_{\Gamma_{2}}(\Omega)} \times \mu^{-1} \overline{\operatorname{rot} \mathbf{R}_{\Gamma_{1}}(\Omega)}
$$

We note that here, in the case of an exterior domain $\Omega$, the respective ranges are not closed.
Proof. The proof is straightforward using Lemma 2.2, i.e., the equivalence of the definition of weak and strong boundary conditions.

Thus $\sigma(\mathcal{M}) \subset \mathbb{R}$, meaning every $\omega \in \mathbb{C} \backslash \mathbb{R}$ is contained in the resolvent set of $\mathcal{M}$ and for given $f \in \mathrm{~L}_{\Lambda}^{2}(\Omega)$ we obtain a unique solution of (3.2) by $u:=(\mathcal{M}-\omega)^{-1} f \in \mathbf{R}_{\Gamma_{1}}(\Omega) \times \mathbf{R}_{\Gamma_{2}}(\Omega)$. Moreover, using the resolvent estimate $\left\|(\mathcal{M}-\omega)^{-1}\right\| \leq|\operatorname{Im} \omega|^{-1}$ and the differential equation we obtain

$$
\|u\|_{\mathbf{R}(\Omega)} \leq c \cdot\left\{\|u\|_{L_{\Lambda}^{2}(\Omega)}+\|f\|_{L_{\Lambda}^{2}(\Omega)}+|\omega|\|u\|_{L_{\Lambda}^{2}(\Omega)}\right\} \leq c \cdot \frac{1+|\omega|}{|\operatorname{Im} \omega|} \cdot\|f\|_{L_{\Lambda}^{2}(\Omega)},
$$

and hence:
Theorem 3.6. For $\omega \in \mathbb{C} \backslash \mathbb{R}$ the solution operator

$$
\mathcal{L}_{\omega}:=(\mathcal{M}-\omega)^{-1}: \mathrm{L}_{\Lambda}^{2}(\Omega) \longrightarrow \mathbf{R}_{\Gamma_{1}}(\Omega) \times \mathbf{R}_{\Gamma_{2}}(\Omega)
$$

is continuous with $\left\|\mathcal{L}_{\omega}\right\|_{L_{\Lambda}^{2}(\Omega), \mathbf{R}(\Omega)} \leq c \cdot \frac{1+|\omega|}{|\operatorname{Im} \omega|}$, where $c$ is independent of $\omega$ and $f$.
Remark 3.7. Let $\omega \in \mathbb{C} \backslash \mathbf{R}$. By Lemma 2.3 the following statements are equivalent to the boundedness of $\mathcal{L}_{\omega}$ :

- (Frierich/Poincaré type estimate) There exists $c>0$ such that

$$
\|u\|_{\mathbf{R}(\Omega)} \leq c \cdot\|(\mathrm{M}-\omega) u\|_{\mathrm{L}_{\Lambda}^{2}(\Omega)} \quad \forall u \in \mathbf{R}_{\Gamma_{1}}(\Omega) \times \mathbf{R}_{\Gamma_{2}}(\Omega)
$$

- (Closed range) The range

$$
\mathcal{R}(\mathcal{M}-\omega)=(\mathrm{M}-\omega)\left(\mathbf{R}_{\Gamma_{1}}(\Omega) \times \mathbf{R}_{\Gamma_{2}}(\Omega)\right)
$$

is closed in $\mathrm{L}_{\Lambda}^{2}(\Omega)$.

The case $\omega \in \mathbb{R} \backslash\{0\}$ is much more challenging, since we want to solve in the continuous spectrum of the Maxwell operator. Clearly this cannot be done for every $f \in \mathrm{~L}_{\Lambda}^{2}(\Omega)$, since otherwise we would have $\mathcal{R}(\mathcal{M}-\omega)=\mathrm{L}_{\Lambda}^{2}(\Omega)$ and therefore $(\mathcal{M}-\omega)^{-1}$ would be continuous (cf. Lemma 2.3) or in other words $\omega \notin \sigma(\mathcal{M})$. Thus we have to restrict ourselves to certain subspaces of $\mathrm{L}_{\Lambda}^{2}$ or generalize our solution concept. Actually, we will do both and show existence as well as uniqueness of weaker, so called "radiating solutions", by switching to data $f \in \mathrm{~L}_{\mathrm{s}}^{2}(\Omega) \times \mathrm{L}_{\mathrm{s}}^{2}(\Omega)$ for some $s>1 / 2$.

Definition 3.8. Let $\omega \in \mathbb{R} \backslash\{0\}$ and $f \in \mathrm{~L}_{\mathrm{loc}}^{2}(\Omega)$. We call $u$ "(radiating) solution" of (3.2), if

$$
u \in \mathbf{R}_{<-\frac{1}{2}, \Gamma_{1}}(\Omega) \times \mathbf{R}_{<-\frac{1}{2}, \Gamma_{2}}(\Omega)
$$

and satisfies

$$
\begin{array}{r}
(\mathrm{M}-\omega) u=f \\
\left(\Lambda_{0}+\sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u \in \mathrm{~L}_{>-\frac{1}{2}}^{2}(\Omega) \tag{3.6}
\end{array}
$$

Remark 3.9. Since

$$
\left(\Lambda_{0}+\sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u=\Lambda_{0}\left(E+\sqrt{\frac{\mu_{0}}{\varepsilon_{0}}} \xi \times H, H-\sqrt{\frac{\varepsilon_{0}}{\mu_{0}}} \xi \times E\right)
$$

the last condition is just the classical Silver-Müller radiation condition which describes the behavior of the electro-magnetic field at infinity and is needed to distinguish outgoing from incoming waves (interchanging + and - in (3.6) would yield incoming waves ).
In order to construct such a solution $u$ we follow the "limiting absorption principle" introduced by Eidus and approximate $u$ by solutions $\left(u_{n}\right)_{n \in \mathbb{N}}$ associated with frequencies $\left(\omega_{n}\right)_{n \in \mathbb{N}} \subset \mathbb{C} \backslash \mathbb{R}$ converging to $\omega \in \mathbb{R} \backslash\{0\}$. This leads to statement (4) of our main result, where the following abbreviations are used:

$$
\mathcal{N}_{\text {gen }}(\mathcal{M}-\omega):=\{u \mid u \text { is a radiating solution of }(\mathrm{M}-\omega) u=0\} \quad(\text { generalized kernel of } \mathcal{M}-\omega)
$$

$$
\left.\sigma_{\text {gen }}(\mathcal{M}):=\left\{\omega \in \mathbb{C} \backslash\{0\} \mid \mathcal{N}_{\text {gen }}(\omega) \neq\{0\}\right\} \quad \text { (generalized spectrum of } \mathcal{M}\right)
$$

Theorem 3.10 ( Fredholm alternative). Let $\Omega \subset \mathbb{R}^{3}$ be an exterior weak Lipschitz domain with boundary $\Gamma$ and weak Lipschitz boundary parts $\Gamma_{1}$ and $\Gamma_{2}:=\Gamma \backslash \bar{\Gamma}_{1}$. Furthermore let $\omega \in \mathbb{R} \backslash\{0\}$ and $\varepsilon, \mu$ $\kappa$-admissible with $\kappa>1$. Then we have:
(1) $\mathcal{N}_{\text {gen }}(\mathcal{M}-\omega) \subset \bigcap_{t \in \mathbb{R}}\left(\mathbf{R}_{\mathrm{t}, \Gamma_{1}}(\Omega) \cap \varepsilon^{-1}{ }_{0} \mathrm{D}_{\mathrm{t}, \Gamma_{2}}(\Omega)\right) \times\left(\mathbf{R}_{\mathrm{t}, \Gamma_{2}}(\Omega) \cap \mu^{-1}{ }_{0} \mathrm{D}_{\mathrm{t}, \Gamma_{1}}(\Omega)\right)^{\mathrm{iii}}$.
(2) $\operatorname{dim} \mathcal{N}_{\text {gen }}(\mathcal{M}-\omega)<\infty$.
(3) $\sigma_{\operatorname{gen}}(\mathcal{M}) \subset \mathbb{R} \backslash\{0\}$ and $\sigma_{\text {gen }}(\mathcal{M})$ has no accumulation point in $\mathbb{R} \backslash\{0\}$.

$$
\begin{aligned}
& \text { iii We even have } \\
& \qquad \mathcal{N}_{\mathrm{gen}}(\mathcal{M}-\omega) \subset \bigcap_{t \in \mathbb{R}}\left(\mathbf{R}_{\mathrm{t}, \Gamma_{1}}(\Omega) \cap \varepsilon^{-1} \operatorname{rot} \mathbf{R}_{\mathrm{t}, \Gamma_{2}}(\Omega)\right) \times\left(\mathbf{R}_{\mathrm{t}, \Gamma_{2}}(\Omega) \cap \mu^{-1} \operatorname{rot} \mathbf{R}_{\mathrm{t}, \Gamma_{1}}(\Omega)\right) .
\end{aligned}
$$

(4) For all $f \in \mathrm{~L}_{>\frac{1}{2}}^{2}(\Omega)$ there exists a (radiating) solution $u$ of (3.2), if and only if

$$
\begin{equation*}
\forall v \in \mathcal{N}_{\operatorname{gen}}(\mathcal{M}-\omega): \quad\langle f, v\rangle_{\mathrm{L}_{\Lambda}^{2}(\Omega)}=0 \tag{3.7}
\end{equation*}
$$

Moreover, we can choose $u$ such that

$$
\begin{equation*}
\forall v \in \mathcal{N}_{\text {gen }}(\mathcal{M}-\omega): \quad\langle u, v\rangle_{L_{\Lambda}^{2}(\Omega)}=0 . \tag{3.8}
\end{equation*}
$$

Then $u$ is uniquely determined.
(5) For all $s,-t>1 / 2$ the solution operator

$$
\mathcal{L}_{\omega}: \mathrm{L}_{\mathrm{s}}^{2}(\Omega) \cap \mathcal{N}_{\mathrm{gen}}(\mathcal{M}-\omega)^{\perp_{\Lambda}} \longrightarrow\left(\mathbf{R}_{\mathrm{t}, \Gamma_{1}}(\Omega) \times \mathbf{R}_{\mathrm{t}, \Gamma_{2}}(\Omega)\right) \cap \mathcal{N}_{\mathrm{gen}}(\mathcal{M}-\omega)^{\perp_{\Lambda}}
$$

defined by (4) is continuous.

Remark 3.11. Under the conditions of Theorem 3.10 the following statements are equivalent to the boundedness of $\mathcal{L}_{\omega}$ (cf. Lemma 2.3):

- (Frierich/Poincaré type estimate) For all $s,-t>1 / 2$ there exists $c>0$ such that

$$
\|u\|_{\mathbf{R}_{\mathrm{t}}(\Omega)} \leq c \cdot\|(\mathrm{M}-\omega) u\|_{\mathrm{L}_{\mathrm{s}}^{2}(\Omega)} \quad \forall u \in\left(\mathbf{R}_{\mathrm{t}, \Gamma_{1}}(\Omega) \times \mathbf{R}_{\mathrm{t}, \Gamma_{2}}(\Omega)\right) \cap \mathcal{N}_{\mathrm{gen}}(\mathcal{M}-\omega)^{\perp_{\Lambda}}
$$

- (Closed range) For all $s,-t>1 / 2$ the range

$$
\mathcal{R}(\mathrm{M}-\omega)=(\mathrm{M}-\omega)\left(\mathbf{R}_{\mathrm{t}, \Gamma_{1}}(\Omega) \times \mathbf{R}_{\mathrm{t}, \Gamma_{2}}(\Omega)\right)
$$

is closed in $\mathrm{L}_{\mathrm{s}}^{2}(\Omega)$.

By the same indirect arguments as in [15, Corollary 3.9] (see also [14, Section 4.9]), we get even stronger estimates for the solution operator $\mathcal{L}_{\omega}$.

Corollary 3.12. Let $\Omega \subset \mathbb{R}^{3}$ be an exterior weak Lipschitz domain with boundary $\Gamma$ and weak Lipschitz boundary parts $\Gamma_{1}$ and $\Gamma_{2}:=\Gamma \backslash \bar{\Gamma}_{1}$. Furthermore let $s,-t>1 / 2$, $\varepsilon, \mu \kappa$-admissible with $\kappa>1$ and $K \Subset \mathbb{C}_{+} \backslash\{0\}$ with $\bar{K} \cap \sigma_{\operatorname{gen}}(\mathcal{M})=\emptyset$. Then
(1) there exist constants $c>0$ and $\hat{t}>-1 / 2$ such that for all $\omega \in \bar{K}$ and $f \in \mathrm{~L}_{\mathrm{s}}^{2}(\Omega)$

$$
\left\|\mathcal{L}_{\omega} f\right\|_{\mathbf{R}_{\mathrm{t}}(\Omega)}+\left\|\left(\Lambda_{0}+\sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) \mathcal{L}_{\omega} f\right\|_{\mathrm{L}_{\hat{t}}^{2}(\Omega)} \leq c \cdot\|f\|_{\mathrm{L}_{\mathrm{s}}^{2}(\Omega)}
$$

holds, implying that $\mathcal{L}_{\omega}: \mathrm{L}_{\mathrm{s}}^{2}(\Omega) \longrightarrow \mathbf{R}_{\mathrm{t}, \Gamma_{1}}(\Omega) \times \mathbf{R}_{\mathrm{t}, \Gamma_{2}}(\Omega)$ is equicontinuous w.r.t. $\omega \in \bar{K}$.
(2) the mapping

$$
\begin{array}{rlcc}
\mathcal{L}: \bar{K} & \longrightarrow \mathrm{~B}\left(\mathrm{~L}_{\mathrm{s}}^{2}(\Omega), \mathbf{R}_{\mathrm{t}, \Gamma_{1}}(\Omega) \times \mathbf{R}_{\mathrm{t}, \Gamma_{2}}(\Omega)\right) \\
\omega & \longmapsto & \mathcal{L}_{\omega}
\end{array}
$$

is uniformly continuous.

## 4. Polynomial Decay and A-Priori Estimate

As stated before, we will construct a solution $u$ in the case of $\omega \in \mathbb{R} \backslash\{0\}$ by solving (3.2) for $\omega_{n}=$ $\omega+i \sigma_{n} \in \mathbb{C}_{+} \backslash \mathbb{R}$ and sending $\sigma_{n} \longrightarrow 0$ ( using $\left(\omega_{n}\right)_{n \in \mathbb{N}} \in \mathbb{C}_{-} \backslash \mathbb{R}$ instead will lead to "incoming" solutions ). The essential ingredients to generate convergence are the polynomial decay of eigensolutions, an a-prioriestimate for solutions corresponding to non-real frequencies and Weck's local selection theorem. While the latter one is already satisfied (cf. Theorem 3.4), we obtain the first two in the spirit of [27] using the following decomposition Lemma introduced in [14] (see also [15], [16]).

Lemma 4.1. Let $\omega \in K \Subset \mathbb{C} \backslash\{0\}, \varepsilon, \mu \kappa$-admissible with $\kappa \geq 0$ and $s, t \in \mathbb{R}$ such that $0 \leq s \in \mathbb{R} \backslash \mathbb{I}$ and $t \leq s \leq t+\kappa$. Moreover, assume that $u \in \mathbf{R}_{\mathrm{t}}(\Omega)$ satisfies $(\mathrm{M}-\omega) u=f \in \mathrm{~L}_{\mathrm{s}}^{2}(\Omega)$. Then we have

$$
f_{1}:=\left(\mathrm{C}_{\mathrm{Rot}, \check{\eta}}-i \omega \check{\eta} \hat{\Lambda}\right) u-i \check{\eta} \Lambda f \in \mathrm{~L}_{\mathrm{s}}^{2}
$$

and by decomposing

$$
f_{1}=f_{\mathrm{R}}+f_{\mathrm{D}}+f_{\mathcal{S}} \in{ }_{0} \mathrm{R}_{\mathrm{s}} \dot{+}_{0} \mathrm{D}_{\mathrm{s}} \dot{+} \mathcal{S}_{s}
$$

according to [26, Theorem 4], we obtain

$$
f_{2}:=f_{\mathrm{D}}+\frac{i}{\omega} \widetilde{\Lambda}_{0}^{-1} \operatorname{Rot} f_{\mathcal{S}} \in{ }_{0} \mathrm{D}_{\mathrm{s}} .
$$

Additionally, u may be decomposed in

$$
u=\eta u+u_{1}+u_{2}+u_{3}
$$

where
(1) $\eta u \in \mathrm{R}_{\mathrm{vox}}(\Omega)$ and for all $\hat{t} \in \mathbb{R}$

$$
\|\eta u\|_{R_{\hat{t}}(\Omega)} \leq c \cdot\left\{\|f\|_{\mathbf{L}_{\mathbf{s}}^{2}(\Omega)}+\|u\|_{\mathrm{L}_{\mathrm{s}-\kappa}^{2}(\Omega)}\right\}
$$

(2) $u_{1}:=-\frac{i}{\omega} \Lambda_{0}^{-1}\left(f_{\mathrm{R}}+f_{\mathcal{S}}\right) \in \mathrm{R}_{\mathrm{s}}$ and

$$
\left\|u_{1}\right\|_{\mathrm{R}_{\mathrm{s}}} \leq c \cdot\left\|f_{1}\right\|_{\mathrm{L}_{\mathrm{s}}^{2}}
$$

(3) $u_{2}:=\mathcal{F}^{-1}\left(\rho^{-2}(1-i r \Xi) \mathcal{F}\left(f_{2}\right)\right) \in \mathbf{H}_{\mathrm{s}}^{1} \cap_{0} \mathrm{D}_{\mathrm{s}}$ and

$$
\left\|u_{2}\right\|_{\mathbf{H}_{\mathbf{s}}^{1}} \leq c \cdot\left\|f_{2}\right\|_{L_{\mathbf{s}}^{2}},
$$

(4) $u_{3}:=\tilde{u}-u_{2} \in \mathbf{H}_{\mathrm{t}}^{2} \cap_{0} \mathrm{D}_{\mathrm{t}}$ and for all $\hat{t} \leq t$

$$
\left\|u_{3}\right\|_{\mathbf{H}_{\hat{t}}^{2}} \leq c \cdot\left\{\left\|u_{3}\right\|_{\mathbf{L}_{\hat{t}}^{2}}+\left\|u_{2}\right\|_{\mathbf{H}_{\hat{t}}^{1}}\right\}
$$

where $\tilde{u}:=i \omega^{-1} \Lambda_{0}^{-1}\left(\operatorname{Rot} \check{\eta} u-f_{\mathrm{D}}\right) \in \mathbf{H}_{\mathrm{t}}^{1} \cap_{0} \mathrm{D}_{\mathrm{t}}$
with constants $c \in(0, \infty)$ independent of $u, f$ or $\omega$. These fields solve the following equations:

$$
\begin{gathered}
\left(\operatorname{Rot}+i \omega \Lambda_{0}\right) \check{\eta}_{1} u=f_{1}, \quad\left(\operatorname{Rot}+i \omega \Lambda_{0}\right) \tilde{u}=f_{2}, \quad\left(\operatorname{Rot}+i \omega \Lambda_{0}\right) u_{3}=\left(1-\omega \Lambda_{0}\right) u_{2} \\
\left(\Delta+\omega^{2} \varepsilon_{0} \mu_{0}\right) u_{3}=\left(1-i \omega \widetilde{\Lambda}_{0}\right) f_{2}-\left(1+\omega^{2} \varepsilon_{0} \mu_{0}\right) u_{2}
\end{gathered}
$$

Moreover the following estimates hold for all $\hat{t} \leq t$ and uniformly w.r.t. $\lambda \in K$, $u$ and $f$ :

$$
\begin{aligned}
& \circ\left\|f_{2}\right\|_{L_{\mathrm{s}}^{2}} \leq c \cdot\left\|f_{1}\right\|_{L_{\mathrm{s}}^{2}} \leq c \cdot\left\{\|f\|_{\mathrm{L}_{\mathrm{s}}^{2}(\Omega)}+\|u\|_{L_{\mathrm{s}-\kappa}^{2}(\Omega)}\right\} \\
& \circ\|u\|_{\mathbf{R}_{\hat{t}}(\Omega)} \leq c \cdot\left\{\|f\|_{\mathrm{L}_{\mathrm{s}}^{2}(\Omega)}+\|u\|_{L_{\mathrm{s}-\kappa}^{2}(\Omega)}+\left\|u_{3}\right\|_{\mathrm{L}_{\hat{t}}^{2}}\right\} \\
& \circ\left\|\left(\Delta+\omega^{2} \varepsilon_{0} \mu_{0}\right) u_{3}\right\|_{\mathrm{L}_{\mathrm{s}}^{2}} \leq c \cdot\left\{\|f\|_{\mathrm{L}_{\mathrm{s}}^{2}(\Omega)}+\|u\|_{\mathrm{L}_{\mathrm{s}-\kappa}^{2}(\Omega)}\right\} \\
& \circ\left\|\left(\operatorname{Rot}-i \lambda \sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u\right\|_{\mathrm{L}_{\hat{t}}^{2}} \leq c \cdot\left\{\|f\|_{\mathrm{L}_{\mathrm{s}}^{2}(\Omega)}+\|u\|_{L_{\mathrm{s}-\kappa}^{2}(\Omega)}+\left\|\left(\operatorname{Rot}-i \lambda \sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u_{3}\right\|_{L_{\hat{t}}^{2}}\right\}
\end{aligned}
$$

Here $\mathcal{S}_{s}$ is a finite dimensional subspace of $\dot{C}^{\infty}\left(\mathbb{R}^{3}\right), \mathcal{F}$ the Fourier-transformation and

$$
\mathrm{C}_{A, B}:=A B-B A
$$

the commutator of $A$ and $B$.
Basically, this lemma allows us to split $u$ into two parts. One part (consisting of $\eta u, u_{1}$ and $u_{2}$ ) has better integrability properties and the other part (consisting of $u_{3}$ ) is more regular and satisfies a Helmholtz equation in the whole of $\mathbb{R}^{3}$. Thus we can use well known results from the theory for Helmholtz equation (cf. Appendix, Section B ) to establish corresponding results for Maxwell's equations. We start with the polynomial decay of solutions, especially of eigensolutions, which will lead to assertions (1) - (3) of our main theorem. Moreover, this will also show, that the solution $u$ we are going to construct, can be
chosen to be perpendicular to the generalized kernel of the time-harmonic Maxwell operator. As in the proof of [16, Theorem 4.2] we obtain ( see also Appendix, Section C) the following.

Lemma 4.2 (Polynomial decay of solutions ). Let $J \subset \mathbb{R} \backslash\{0\}$ be some interval, $\omega \in J, \varepsilon, \mu \kappa$-admissible with $\kappa>1$ and $s \in \mathbb{R} \backslash \mathbb{I}$ with $s>1 / 2$. If

$$
u \in \mathbf{R}_{>-\frac{1}{2}}(\Omega) \text { satisfies }(\mathrm{M}-\omega) u=: f \in \mathrm{~L}_{\mathrm{s}}^{2}(\Omega)
$$

we have

$$
u \in \mathbf{R}_{\mathrm{s}-1}(\Omega) \quad \text { and } \quad\|u\|_{\mathbf{R}_{\mathrm{s}-1}(\Omega)} \leq c \cdot\left\{\|f\|_{\mathrm{L}_{\mathrm{s}}^{2}(\Omega)}+\|u\|_{\mathrm{L}^{2}(\Omega(\delta))}\right\}
$$

with $c, \delta \in(0, \infty)$ independent of $\omega, u$ and $f$.
In short: If a solution $u$ satisfies $u \in \mathbf{R}_{\mathrm{t}}(\Omega)$ for some $t>-1 / 2$ and the right hand side $f=(\mathrm{M}-\omega) u$ has a better integrability properties, meaning $f \in \mathrm{~L}_{\mathrm{s}}^{2}(\Omega)$ for some $s>1 / 2$, then $u$ is also better integrable, i.e., $u \in \mathbf{R}_{\mathrm{s}}(\Omega)$. Especially, if

$$
u \in \mathbf{R}_{>-\frac{1}{2}}(\Omega) \quad \text { and } \quad f \in \mathrm{~L}_{\mathrm{s}}^{2}(\Omega) \quad \forall s \in \mathbb{R}
$$

then $u \in \mathbf{R}_{\mathbf{s}}(\Omega)$ for all $s \in \mathbb{R}$, which is called "polynomial decay".
Corollary 4.3. Let $\omega \in \mathbb{R} \backslash\{0\}$ and assume $\varepsilon, \mu$ are $\kappa$-admissible with $\kappa>1$ and

$$
u \in \mathbf{R}_{<-\frac{1}{2}, \Gamma_{1}}(\Omega) \times \mathbf{R}_{<-\frac{1}{2}, \Gamma_{2}}(\Omega)
$$

is a radiating solution (cf. Definition 3.8) of $(\mathrm{M}-\omega) u=0$. Then:

$$
u \in \bigcap_{t \in \mathbb{R}}\left(\mathbf{R}_{\mathrm{t}, \Gamma_{1}}(\Omega) \times \mathbf{R}_{\mathrm{t}, \Gamma_{2}}(\Omega)\right) .
$$

Proof. According to Lemma 4.2 it suffices to show $u \in \mathbf{R}_{\mathrm{t}}(\Omega)$ for some $t>-1 / 2$. Therefore remember that $u$ is a radiating solution, hence the radiation condition (3.6) holds and there exists $\hat{t}>-1 / 2$ such that

$$
\begin{equation*}
\left(\Lambda_{0}+\sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u \in \mathrm{~L}_{\hat{\mathrm{t}}}^{2}(\Omega) . \tag{4.1}
\end{equation*}
$$

On the other hand we have

$$
\begin{aligned}
& \left\|\left(\Lambda_{0}+\sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u\right\|_{\mathrm{L}_{\hat{t}}^{2}\left(\mathrm{G}\left(r_{0}, \tilde{r}\right)\right)}^{2} \\
& \quad=\left\|\Lambda_{0} u\right\|_{\mathrm{L}_{\hat{t}}^{2}\left(\mathrm{G}\left(r_{0}, \tilde{r}\right)\right)}^{2}+2 \sqrt{\varepsilon_{0} \mu_{0}} \operatorname{Re}\left\langle\Xi u, \Lambda_{0} u\right\rangle_{\mathrm{L}_{\hat{t}}^{2}\left(\mathrm{G}\left(r_{0}, \tilde{r}\right)\right)}+\varepsilon_{0} \mu_{0}\|\Xi u\|_{\mathrm{L}_{\hat{t}}^{2}\left(\mathrm{G}\left(r_{0}, \tilde{r}\right)\right)}^{2}
\end{aligned}
$$

and using Lemma A. 3 (cf. Appendix, Section A) with

$$
\phi(s):=\left(1+s^{2}\right)^{\hat{t}}, \quad \Phi:=\phi \circ r, \quad \psi(s)=\int_{\max \left\{r_{0}, s\right\}}^{\tilde{r}} \phi(\sigma) d \sigma, \quad \Psi=\psi \circ r,
$$

as well as the differential equation, we conclude

$$
\begin{aligned}
& \operatorname{Re}\left\langle\Xi u, \Lambda_{0} u\right\rangle_{\mathrm{L}_{\hat{t}}^{2}\left(\mathrm{G}\left(r_{0}, \tilde{r}\right)\right)}=\operatorname{Re}\left\langle\Phi \Xi u, \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}\left(\mathrm{G}\left(r_{0}, \tilde{r}\right)\right)} \\
& =\operatorname{Re}\left(\left\langle\Psi \operatorname{Rot} u, \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\Omega(\tilde{r}))}+\left\langle\Psi u, \widetilde{\Lambda}_{0} \operatorname{Rot} u\right\rangle_{\mathrm{L}^{2}(\Omega(\tilde{r}))}\right) \\
& =\operatorname{Re}\left(\left\langle-i \omega \Psi \Lambda u, \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\Omega(\tilde{r}))}+\left\langle\Psi u,-i \omega \widetilde{\Lambda}_{0} \Lambda u\right\rangle_{\mathrm{L}^{2}(\Omega(\tilde{r}))}\right) \\
& =\operatorname{Re} \underbrace{i \omega\left\langle\Psi \Lambda u,\left(\widetilde{\Lambda}_{0}-\Lambda_{0}\right) u\right\rangle_{L^{2}(\Omega(\tilde{r}))}}_{\in i \mathbb{R}}=0,
\end{aligned}
$$

hence

$$
\|u\|_{\mathrm{L}_{\hat{t}}^{2}\left(\mathrm{G}\left(r_{0}, \tilde{r}\right)\right)} \leq c \cdot\left\|\left(\Lambda_{0}+\sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u\right\|_{\mathrm{L}_{\hat{t}}^{2}\left(\mathrm{G}\left(r_{0}, \tilde{r}\right)\right)}
$$

with $c \in(0, \infty)$ independent of $\tilde{r}$. Now the monotone convergence theorem and (4.1) show

$$
\|u\|_{\mathrm{L}_{\hat{t}}^{2}\left(\check{\mathrm{U}}\left(r_{0}\right)\right)} \leq c \cdot\left\|\left(\Lambda_{0}+\sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u\right\|_{\mathrm{L}_{\hat{t}}^{2}\left(\check{\mathrm{U}}\left(r_{0}\right)\right.}<\infty
$$

which already implies $u \in \mathrm{~L}_{\hat{\mathrm{t}}}^{2}(\Omega)$ and completes the proof.

The next step is an a-priori estimate for solutions corresponding to non-real frequencies, which will later guarantee that our solution satisfies the radiation condition (3.6) and has the proper integrability. The proof of it is practically identical with the proof of [16, Lemma 6.3] (cf. Appendix, Section C).

Lemma 4.4 (A-priori estimate for Maxwell's equations). Let $J \Subset \mathbb{R} \backslash\{0\}$ be some interval, $-t, s>1 / 2$ and $\varepsilon, \mu \kappa$-admissible with $\kappa>1$. Then there exists constants $c, \delta \in(0, \infty)$ and some $\hat{t}>-1 / 2$, such that for all $\omega \in \mathbb{C}_{+}$with $\omega^{2}=\lambda^{2}+i \lambda \sigma, \lambda \in J, \sigma \in\left(0,{\sqrt{\varepsilon_{0} \mu_{0}}}^{-1}\right]$ and $f \in \mathrm{~L}_{\mathrm{s}}^{2}(\Omega)$ we have:

$$
\left\|\mathcal{L}_{\omega} f\right\|_{\mathbf{R}_{\mathrm{t}}(\Omega)}+\left\|\left(\Lambda_{0}+\sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) \mathcal{L}_{\omega} f\right\|_{\mathrm{L}_{\mathrm{t}}^{2}(\Omega)} \leq c \cdot\left\{\|f\|_{L_{s}^{2}(\Omega)}+\left\|\mathcal{L}_{\omega} f\right\|_{\mathrm{L}^{2}(\Omega(\delta))}\right\}
$$

## 5. Proof of the Main Result

Before we start with the proof of Theorem 3.10 we provide some Helmholtz type decompositions, which will be useful in the following. These are immediate consequences of the projection theorem and Lemma 2.2.

Lemma 5.1. We have

$$
\begin{aligned}
\mathrm{L}_{\varepsilon}^{2}(\Omega) & =\overline{\nabla \mathbf{H}_{\Gamma_{1}}^{1}(\Omega)} \oplus_{\varepsilon} \varepsilon^{-1}{ }_{0} \mathrm{D}_{\Gamma_{2}}(\Omega), & \mathrm{L}_{\varepsilon}^{2}(\mu) & =\overline{\nabla \mathbf{H}_{\Gamma_{2}}^{1}(\Omega)} \oplus_{\mu} \mu^{-1}{ }_{0} \mathrm{D}_{\Gamma_{1}}(\Omega), \\
\mathbf{R}_{\Gamma_{1}}(\Omega) & =\overline{\nabla \mathbf{H}_{\Gamma_{1}}^{1}(\Omega)} \oplus_{\varepsilon}\left(\mathbf{R}_{\Gamma_{1}}(\Omega) \cap \varepsilon^{-1}{ }_{0} \mathrm{D}_{\Gamma_{2}}(\Omega)\right), & \mathbf{R}_{\Gamma_{2}}(\Omega) & =\overline{\nabla \mathbf{H}_{\Gamma_{2}}^{1}(\Omega)} \oplus_{\mu}\left(\mathbf{R}_{\Gamma_{2}}(\Omega) \cap \mu^{-1}{ }_{0} \mathrm{D}_{\Gamma_{1}}(\Omega)\right),
\end{aligned}
$$

where the closures are taken in $\mathrm{L}^{2}(\Omega)$.
Proof. Let $\gamma \in\{\varepsilon, \mu\}$ and $i, j \in\{1,2\}$ with $i \neq j$. The linear operator

$$
\nabla_{i}: \mathbf{H}_{\Gamma_{\mathrm{i}}}^{1}(\Omega) \subset \mathrm{L}^{2}(\Omega) \longrightarrow \mathrm{L}_{\gamma}^{2}(\Omega)
$$

is densely defined and closed with adjoint (cf. Lemma 2.2)

$$
-\operatorname{div}_{j} \gamma: \gamma^{-1} \mathbf{D}_{\Gamma_{j}}(\Omega) \subset \mathrm{L}_{\gamma}^{2}(\Omega) \longrightarrow \mathrm{L}^{2}(\Omega)
$$

The projection theorem yields

$$
\mathrm{L}_{\gamma}^{2}(\Omega)=\overline{\mathcal{R}\left(\nabla_{i}\right)} \oplus_{\gamma} \mathcal{N}\left(\operatorname{div}_{j} \gamma\right)
$$

The remaining assertion follows by $\nabla \mathbf{H}_{\Gamma_{\mathrm{i}}}^{1}(\Omega) \subset \mathrm{R}_{\Gamma_{\mathrm{i}}}(\Omega)$.

Proof of Theorem 3.10. Let $\omega \in \mathbb{R} \backslash\{0\}$ and $\varepsilon, \mu$ be $\kappa$-admissible for some $\kappa>1$.
(1): The assertion follows by Corollary 4.3 and the differential equation

$$
(\mathrm{M}-\omega) u=0 \Longleftrightarrow u=i \omega^{-1} \Lambda^{-1} \operatorname{Rot} u
$$

using the fact that (cf. Lemma 2.2)

$$
\operatorname{rot} \mathbf{R}_{\mathrm{t}, \Gamma_{1}}(\Omega) \subset{ }_{0} \mathrm{D}_{\mathrm{t}, \Gamma_{1}}(\Omega) \quad \text { resp. } \quad \operatorname{rot} \mathbf{R}_{\mathrm{t}, \Gamma_{2}}(\Omega) \subset{ }_{0} \mathrm{D}_{\mathrm{t}, \Gamma_{2}}(\Omega)
$$

(2): Assume $\operatorname{dim} \mathcal{N}_{\operatorname{gen}}(\mathcal{M}-\omega)=\infty$. Using (1) there exists a $\mathrm{L}_{\Lambda}^{2}$-orthonormal sequence $\left(u_{n}\right)_{n \in \mathbb{N}} \subset$ $\mathcal{N}_{\text {gen }}(\mathcal{M}-\omega)$ converging weakly in $L^{2}(\Omega)$ to 0 . By the differential equation this sequence is bounded in $\left(\mathbf{R}_{\Gamma_{1}}(\Omega) \cap \varepsilon^{-1}{ }_{0} \mathrm{D}_{\Gamma_{2}}(\Omega)\right) \times\left(\mathbf{R}_{\Gamma_{2}}(\Omega) \cap \mu^{-1}{ }_{0} \mathrm{D}_{\Gamma_{1}}(\Omega)\right)$. Hence, due to Weck's local selection theorem, we can choose a subsequence, $\left(u_{\pi(n)}\right)_{n \in \mathbb{N}}$ converging to 0 in $\mathrm{L}_{\text {loc }}^{2}(\bar{\Omega})\left(\left(u_{\pi(n)}\right)_{n \in \mathbb{N}}\right.$ also converges weakly on
every bounded subset ). Now let $1<s \in \mathbb{R} \backslash \mathbb{I}$. Then Lemma 4.2 guarantees the existence of $c, \delta \in(0, \infty)$ independent of $\left(u_{\pi(n)}\right)_{n \in \mathbb{N}}$ such that

$$
1=\left\|u_{\pi(n)}\right\|_{\mathrm{L}_{\Lambda}^{2}(\Omega)} \leq c \cdot\left\|u_{\pi(n)}\right\|_{\mathbf{R}_{\mathrm{s}-1}(\Omega)} \leq c \cdot\left\|u_{\pi(n)}\right\|_{\mathrm{L}^{2}(\Omega(\delta))} \xrightarrow{n \rightarrow \infty} 0
$$

holds; a contradiction.
(3): $\mathcal{M}$ is a selfadjoint operator, hence we clearly have $\sigma_{\operatorname{gen}}(\mathcal{M}) \subset \mathbb{R} \backslash\{0\}$. Now assume $\widetilde{\omega} \in \mathbb{R} \backslash\{0\}$ is an accumulation point of $\sigma_{\operatorname{gen}}(\mathcal{M})$. Then we can choose a sequence $\left(\omega_{n}\right)_{n \in \mathbb{N}} \subset \mathbb{R} \backslash\{0\}$ with $\omega_{n} \neq \omega_{m}$ for $n \neq m, \omega_{n} \longrightarrow \widetilde{\omega}$ and a corresponding sequence $\left(u_{n}\right)_{n \in \mathbb{N}}$ with $u_{n} \in \mathcal{N}_{\operatorname{gen}}\left(\mathcal{M}-\omega_{n}\right) \backslash\{0\}$. As $\mathcal{M}$ is selfadjoint, eigenvectors associated to different eigenvalues are orthogonal provided they are well enough integrable (which is given by (1)) and thus by normalizing $\left(u_{n}\right)_{n \in \mathbb{N}}$ we end up with an $\mathrm{L}_{\Lambda}^{2}$-orthonormal sequence. Continuing as in (2), we again obtain a contradiction.
(4): First of all, if a solution $u$ satisfies (3.8), it is uniquely determined as for the homogeneous problem $u \in \mathcal{N}_{\text {gen }}(\mathcal{M}-\omega)$ together with (1) and (3.8) implies $u=0$. Moreover, using Lemma 2.2 and (1), we obtain

$$
\langle f, v\rangle_{\mathrm{L}_{\Lambda}^{2}(\Omega)}=\langle(\mathrm{M}-\omega) u, v\rangle_{\mathrm{L}_{\Lambda}^{2}(\Omega)}=\langle u,(\mathrm{M}-\omega) v\rangle_{\mathrm{L}_{\Lambda}^{2}(\Omega)}=0 \quad \forall v \in \mathcal{N}_{\operatorname{gen}}(\mathcal{M}-\omega),
$$

meaning (3.7) is necessary. In order to show, that (3.7) is also sufficient, we use Eidus' principle of limiting absorption. Therefore let $s>1 / 2$ and $f \in \mathrm{~L}_{\mathrm{s}}^{2}(\Omega)$ satisfy (3.7). We take a sequence $\left(\sigma_{n}\right)_{n \in \mathbb{N}} \subset \mathbb{R}_{+}$with $\sigma_{n} \longrightarrow 0$ and construct a sequence of frequencies

$$
\left(\omega_{n}\right)_{n \in \mathbb{N}}, \quad \omega_{n}:=\sqrt{\omega^{2}+i \sigma_{n} \omega} \in \mathbb{C}_{+} \backslash \mathbb{R}
$$

converging to $\omega$. Since $\mathcal{M}$ is a selfadjoint operator we obtain (cf. Section 3 ) a corresponding sequence of solutions $\left(u_{n}\right)_{n \in \mathbb{N}}, u_{n}:=\mathcal{L}_{\omega_{n}} f \in \mathbf{R}_{\Gamma_{1}}(\Omega) \times \mathbf{R}_{\Gamma_{2}}(\Omega)$ satisfying $\left(\mathrm{M}-\omega_{n}\right) u_{n}=f$. Now our aim is to show that this sequence or at least a subsequence is converging to a solution $u$. By Lemma 5.1 we decompose

$$
u_{n}=\hat{u}_{n}+\tilde{u}_{n} \quad \text { and } \quad f=\hat{f}+\tilde{f}
$$

with

$$
\begin{align*}
& \hat{u}_{n}, \hat{f} \in \overline{\nabla \mathbf{H}_{\Gamma_{1}}^{1}(\Omega)} \times \overline{\nabla \mathbf{H}_{\Gamma_{2}}^{1}(\Omega)} \subset{ }_{0} \mathbf{R}_{\Gamma_{1}}(\Omega) \times{ }_{0} \mathbf{R}_{\Gamma_{2}}(\Omega), \\
& \tilde{u}_{n}, \tilde{f} \in\left(\mathbf{R}_{\Gamma_{1}}(\Omega) \cap \varepsilon^{-1}{ }_{0} \mathrm{D}_{\Gamma_{2}}(\Omega)\right) \times\left(\mathbf{R}_{\Gamma_{2}}(\Omega) \cap \mu^{-1}{ }_{0} \mathrm{D}_{\Gamma_{1}}(\Omega)\right) . \tag{5.1}
\end{align*}
$$

Inserting these (orthogonal) decompositions in the differential equation we end up with two equations

$$
-\omega_{n} \hat{u}_{n}=\hat{f} \quad \text { and } \quad\left(\mathrm{M}-\omega_{n}\right) \tilde{u}_{n}=\tilde{f}
$$

noting that the first one is trivial and implies $\mathrm{L}^{2}$-convergence of $\left(\hat{u}_{n}\right)_{n \in \mathbb{N}}$. For dealing with the second equation we need the following additional assumption on $\left(u_{n}\right)_{n \in \mathbb{N}}$, which we will prove in the end:

$$
\begin{equation*}
\forall t<-1 / 2 \quad \exists c \in(0, \infty) \quad \forall n \in \mathbb{N}: \quad\left\|u_{n}\right\|_{L_{t}^{2}(\Omega)} \leq c \tag{5.2}
\end{equation*}
$$

Let $\hat{t}<-1 / 2$ and $c \in(0, \infty)$ such that (5.2) holds. Then, by construction and $(5.1)_{2}$, the sequence $\left(\tilde{u}_{n}\right)_{n \in \mathbb{N}}$ is bounded in $\left(\mathbf{R}_{\hat{t}, \Gamma_{1}}(\Omega) \cap \varepsilon^{-1}{ }_{0} \mathrm{D}_{\hat{\mathrm{t}}, \Gamma_{2}}(\Omega)\right) \times\left(\mathbf{R}_{\hat{\mathrm{t}}, \Gamma_{2}}(\Omega) \cap \mu^{-1}{ }_{0} \mathrm{D}_{\hat{\mathrm{t}}, \Gamma_{1}}(\Omega)\right)$. Hence (Theorem 3.4 and Lemma 3.3), $\left(\tilde{u}_{n}\right)_{n \in \mathbb{N}}$ has a subsequence $\left(\tilde{u}_{\pi(n)}\right)_{n \in \mathbb{N}}$ converging in $L_{\tilde{t}}^{2}(\Omega)$ for some $\tilde{t}<\hat{t}$ and by (5.1) $)_{2}$ even in $\mathbf{R}_{\tilde{t}, \Gamma_{1}}(\Omega) \times \mathbf{R}_{\tilde{\mathrm{t}}, \Gamma_{2}}(\Omega)$. Consequently, the entire sequence $\left(u_{\pi(n)}\right)_{n \in \mathbb{N}}$ converges in $\mathbf{R}_{\tilde{t}}(\Omega)$ to some $u$ satisfying

$$
u \in \mathbf{R}_{\tilde{t}, \Gamma_{1}}(\Omega) \times \mathbf{R}_{\tilde{t}, \Gamma_{2}}(\Omega) \quad \text { and } \quad(\mathrm{M}-\omega) u=f
$$

Moreover, with Corollary 4.3 and Lemma 2.2 we obtain for $n \in \mathbb{N}$ and arbitrary $v \in \mathcal{N}_{\text {gen }}(\mathcal{M}-\omega)$

$$
\begin{aligned}
0=\langle f, v\rangle_{\mathrm{L}_{\Lambda}^{2}(\Omega)} & =\left\langle\left(\mathrm{M}-\omega_{\pi(n)}\right) u_{\pi(n)}, v\right\rangle_{\mathrm{L}_{\Lambda}^{2}(\Omega)} \\
& =\left\langle u_{\pi(n)},\left(\mathrm{M}-\bar{\omega}_{\pi(n)}\right) v\right\rangle_{\mathrm{L}_{\Lambda}^{2}(\Omega)}=\left(\omega-\omega_{\pi(n)}\right) \cdot\left\langle u_{\pi(n)}, v\right\rangle_{\mathrm{L}_{\Lambda}^{2}(\Omega)},
\end{aligned}
$$

hence $\left\langle u_{\pi(n)}, v\right\rangle_{\mathrm{L}_{\Lambda}^{2}(\Omega)}=0$ and as $\langle\cdot, v\rangle_{\mathrm{L}_{\Lambda}^{2}(\Omega)}$ is continuous on $\mathrm{L}_{\tilde{\mathrm{t}}}^{2}(\Omega) \times \mathrm{L}_{\tilde{\mathrm{t}}}^{2}(\Omega)$ by (1) we obtain

$$
\langle u, v\rangle_{\mathrm{L}_{\Lambda}^{2}(\Omega)}=\lim _{n \rightarrow \infty}\left\langle u_{\pi(n)}, v\right\rangle_{\mathrm{L}_{\Lambda}^{2}(\Omega)}=0 .
$$

Thus, up to now, we have constructed a vector field $u \in \mathcal{N}_{\text {gen }}(\mathcal{M}-\omega)^{\perp_{\Lambda}}$, which has the right boundary conditions and satisfies the differential equation. But for being a radiating solution, it still remains to show, that $u \in \mathbf{R}_{<-\frac{1}{2}}(\Omega)$ and enjoys the radiation condition (3.6). For that let $t<-1 / 2$. Then, by Lemma 4.4, there exist $c, \delta \in(0, \infty)$ and some $\check{t}>-1 / 2$, such that for $n \in \mathbb{N}$ large enough we obtain uniformly in $\sigma_{\pi(n)}, u_{\pi(n)}, f$ and $\tilde{r}>0$ :

$$
\left\|u_{\pi(n)}\right\|_{\mathbf{R}_{\mathrm{t}}(\Omega(\tilde{r}))}+\left\|\left(\Lambda_{0}+\sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u_{\pi(n)}\right\|_{\mathrm{L}_{\mathrm{t}}^{2}(\Omega(\tilde{r}))} \leq c \cdot\left\{\|f\|_{\mathrm{L}_{s}^{2}(\Omega)}+\left\|u_{\pi(n)}\right\|_{\mathrm{L}^{2}(\Omega(\delta))}\right\} .
$$

Sending $n \longrightarrow \infty$ and afterwards $\tilde{r} \longrightarrow \infty$ (monotone convergence) we obtain

$$
\begin{equation*}
\|u\|_{\mathbf{R}_{t}(\Omega)}+\left\|\left(\Lambda_{0}+\sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u\right\|_{L_{t}^{2}(\Omega)} \leq c \cdot\left\{\|f\|_{L_{s}^{2}(\Omega)}+\|u\|_{L^{2}(\Omega(\delta))}\right\}<\infty \tag{5.3}
\end{equation*}
$$

yielding

$$
u \in \mathbf{R}_{<-\frac{1}{2}}(\Omega) \quad \text { and } \quad\left(\Lambda_{0}+\sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u \in \mathrm{~L}_{>-\frac{1}{2}}^{2}(\Omega)
$$

This completes the proof of existence, if we can show (5.2). To this end we assume it to be wrong, i.e., there exists $t<-1 / 2$ and a sequence $\left(u_{n}\right)_{n \in \mathbb{N}} \subset \mathbf{R}_{\mathrm{t}, \Gamma_{1}}(\Omega) \times \mathbf{R}_{\mathrm{t}, \Gamma_{2}}(\Omega), u_{n}:=\mathcal{L}_{\omega_{n}} f$ with $\left\|u_{n}\right\|_{L_{\mathrm{t}}^{2}(\Omega)} \longrightarrow \infty$ for $n \longrightarrow \infty$. Defining

$$
\check{u}_{n}:=\left\|u_{n}\right\|_{\mathrm{L}_{\mathrm{t}}^{2}(\Omega)}^{-1} \cdot u_{n} \quad \text { and } \quad \check{f}_{n}:=\left\|u_{n}\right\|_{\mathrm{L}_{\mathrm{t}}^{2}(\Omega)}^{-1} \cdot f
$$

we have

$$
\left\|\check{u}_{n}\right\|_{L_{\mathrm{t}}^{2}(\Omega)}=1, \quad \check{f} \longrightarrow 0 \text { in } \mathrm{L}_{\mathrm{s}}^{2}(\Omega) \quad \text { and } \quad\left(\mathrm{M}-\omega_{n}\right) \check{u}_{n}=\check{f}_{n} .
$$

Then, repeating the arguments from above, we obtain some $\check{t}<t$ and a subsequence $\left(\check{u}_{\pi(n)}\right)_{n \in \mathbb{N}}$ converging in $L_{\check{t}}^{2}(\Omega)$ to some $\check{u} \in \mathcal{N}_{\text {gen }}(\mathcal{M}-\omega) \cap \mathcal{N}_{\text {gen }}(\mathcal{M}-\omega)^{\perp_{\Lambda}}$, hence $\check{u}=0$. But Lemma 4.4 ensures the existence of $c, \delta \in(0, \infty)$ (independent of $\sigma_{\pi(n)}, \check{u}_{\pi(n)}$ and $\left.\check{f}_{\pi(n)}\right)$ such that

$$
1=\left\|\check{u}_{\pi(n)}\right\|_{L_{t}^{2}(\Omega)} \leq c \cdot\left\{\left\|\check{f}_{\pi(n)}\right\|_{L_{s}^{2}(\Omega)}+\left\|\check{u}_{\pi(n)}\right\|_{L^{2}(\Omega(\delta))}\right\} \xrightarrow{n \rightarrow \infty} 0
$$

holds; a contradiction.
(5): Let $-t, s>1 / 2$. By (4) the solution operator

$$
\mathcal{L}_{\omega}: \underbrace{\mathrm{L}_{\mathrm{s}}^{2}(\Omega) \cap \mathcal{N}_{\mathrm{gen}}(\mathcal{M}-\omega)^{\perp_{\Lambda}}}_{=: \mathcal{D}\left(\mathcal{L}_{\omega}\right)} \longrightarrow \underbrace{\left(\mathbf{R}_{\mathrm{t}, \Gamma_{1}}(\Omega) \times \mathbf{R}_{\mathrm{t}, \Gamma_{2}}(\Omega)\right) \cap \mathcal{N}_{\mathrm{gen}}(\mathcal{M}-\omega)^{\perp_{\Lambda}}}_{=: \mathcal{R}\left(\mathcal{L}_{\omega}\right)}
$$

is well defined. Furthermore, due to the polynomial decay of eigensolutions, $\mathcal{D}\left(\mathcal{L}_{\omega}\right)$ is closed in $L_{s}^{2}(\Omega)$. Thus, the assertion follows from the closed graph theorem, if we can show that $\mathcal{L}_{\omega}$ is closed. Therefore, take $\left(f_{n}\right)_{n \in \mathbb{N}} \subset \mathcal{D}\left(\mathcal{L}_{\omega}\right)$ with

$$
f_{n} \longrightarrow f \text { in } \mathbf{L}_{\mathrm{s}}^{2}(\Omega) \quad \text { and } \quad u_{n}:=\mathcal{L}_{\omega} f_{n} \longrightarrow u \text { in } \mathbf{R}_{\mathrm{t}, \Gamma_{1}}(\Omega) \times \mathbf{R}_{\mathrm{t}, \Gamma_{2}}(\Omega)
$$

Then clearly $f \in \mathcal{D}\left(\mathcal{L}_{\omega}\right), u \in \mathcal{R}\left(\mathcal{L}_{\omega}\right)$ and as $(\mathrm{M}-\omega) u_{n}=f_{n}$, we obtain $(\mathrm{M}-\omega) u=f$. Now estimate (5.3) (along with monotone convergence) shows as before

$$
u \in \mathbf{R}_{<-\frac{1}{2}}(\Omega) \quad \text { and } \quad\left(\Lambda_{0}+\sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u \in \mathrm{~L}_{>-\frac{1}{2}}^{2}(\Omega)
$$

meaning $u$ is a radiating solution, i.e., $u=\mathcal{L}_{\omega} f$, which completes the proof.

Remark 5.2. During the discussion at AANMPDE10 (10th Workshop on Analysis and Advanced Numerical Methods for Partial Differential Equations ), M. Waurick and S. Trostorff pointed out, that it is sufficient to use weakly convergent subsequences for the construction of the (radiating) solution. This is in fact true (the radiation condition and regularity properties follow from Lemma (4.4) by the boundedness of the sequence and the weak lower semicontinuity of the norms ), but it should be noted, that Weck's local selection theorem is still needed to proof (5.2), since here norm convergence is indispensable in order to generate a contradiction. Anyway, we thank both for the vivid discussion and constructive criticism.

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## Appendix A. Technical Tools

Lemma A.1. Let $\Omega \subset \mathbb{R}^{3}$ be an arbitrary exterior domain and $s, t, \theta \in \mathbb{R}$ with $t<s$ and $\theta>0$. Then there exist constants $c, \delta \in(0, \infty)$ such that

$$
\|w\|_{L_{t}^{2}(\Omega)} \leq c \cdot\|w\|_{L^{2}(\Omega(\delta))}+\theta \cdot\|w\|_{L_{s}^{2}(\Omega)}
$$

holds for all $w \in \mathrm{~L}_{\mathrm{s}}^{2}(\Omega)$.
Proof. Let $\mathbb{R}^{3} \backslash \Omega \subset \mathrm{U}\left(r_{0}\right)$. For $\tilde{r} \geq r_{0}$ we obtain:

$$
\begin{aligned}
\|w\|_{\mathrm{L}_{\mathrm{t}}^{2}(\Omega)}^{2} & =\|w\|_{\mathrm{L}_{\mathrm{t}}^{2}(\Omega(\tilde{r}))}^{2}+\|w\|_{\mathrm{L}_{\mathrm{t}}^{2}(\check{\mathrm{U}}(\tilde{r}))}^{2} \\
& \leq\left(1+\tilde{r}^{2}\right)^{\max \{0, t\}} \cdot\|w\|_{\mathrm{L}^{2}(\Omega(\tilde{r}))}^{2}+\left(1+\tilde{r}^{2}\right)^{t-s} \cdot\|w\|_{\mathrm{L}_{s}^{2}(\check{\mathrm{U}}(\tilde{r}))}^{2} \\
& \leq\left(1+\tilde{r}^{2}\right)^{\max \{0, t\}} \cdot\|w\|_{\mathrm{L}^{2}(\Omega(\tilde{r}))}^{2}+\left(1+\tilde{r}^{2}\right)^{t-s} \cdot\|w\|_{\mathrm{L}_{s}^{2}(\Omega)}
\end{aligned}
$$

Since $t<s$ we can choose $\tilde{r}$ such that $\left(1+\tilde{r}^{2}\right)^{t-s} \leq \theta^{2}$, which completes the proof.

Lemma A.2. For $\tilde{r}>0$ and $f \in \mathrm{~L}^{1}\left(\mathbb{R}^{n}\right)$ we have

$$
\liminf _{r \rightarrow \infty} r \int_{\mathrm{S}(r)}|f| d \lambda_{s}^{n-1}=0
$$

Proof. Otherwise there exists $\hat{r}>0$ and $c>0$ such that

$$
\int_{\mathrm{S}(r)}|f| d \lambda_{s}^{n-1} \geq \frac{c}{r} \quad \forall r \geq \hat{r}
$$

and using Fubini's theorem we obtain

$$
\|f\|_{\mathrm{L}^{1}\left(\mathbb{R}^{n}\right)}^{2} \geq \int_{\check{\mathrm{U}}(\hat{r})}|f| d \lambda_{s}^{n-1}=\int_{\hat{r}}^{\infty} \int_{\mathrm{S}(r)}|f| d \lambda_{s}^{n-1} d r \geq c \cdot \int_{\hat{r}}^{\infty} \frac{1}{r} d r=\infty
$$

a contradiction.

Lemma A.3. Let $\Omega \subset \mathbb{R}^{3}$ be an exterior weak Lipschitz domain with boundary $\Gamma$ and weak Lipschitz boundary parts $\Gamma_{1}$ and $\Gamma_{2}:=\Gamma \backslash \bar{\Gamma}_{1}$. Furthermore let $\hat{r}, \tilde{r} \in \mathbb{R}_{+}$with $\tilde{r}>\hat{r}$ and $\mathbb{R}^{3} \backslash \Omega \subset \mathrm{U}(\hat{r})$ as well as $\phi \in \mathrm{C}^{0}([\hat{r}, \tilde{r}], \mathbb{C})$. If $u:=(E, H) \in \mathbf{R}_{\mathrm{t}, \Gamma_{1}}(\Omega) \times \mathbf{R}_{\mathrm{t}, \Gamma_{2}}(\Omega)$ for some $t \in \mathbb{R}$, we have

$$
\begin{equation*}
\left\langle\Psi \operatorname{Rot} u, \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\Omega(\tilde{r}))}+\left\langle\Psi u, \operatorname{Rot} \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\Omega(\tilde{r}))}=\left\langle\Phi \Xi u, \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\mathrm{G}(\hat{r}, \tilde{r}))} \tag{A.1}
\end{equation*}
$$

where $\Phi:=\phi \circ r, \Psi:=\psi \circ r$

$$
\psi:[0, \tilde{r}] \longrightarrow \mathbb{C}, s \longmapsto \int_{\max \{\hat{r}, s\}}^{\tilde{r}} \phi(\sigma) d \sigma \quad \text { and } \quad \Lambda_{0}=\left(\begin{array}{cc}
\varepsilon_{0} & 0 \\
0 & \mu_{0}
\end{array}\right), \varepsilon_{0}, \mu_{0} \in \mathbb{R}_{+} .
$$

Proof. As $\mathrm{C}_{\Gamma_{1}}^{\infty}(\Omega)$ resp. $\mathrm{C}_{\Gamma_{2}}^{\infty}(\Omega)$ are dense in $\mathbf{R}_{\mathrm{t}, \Gamma_{1}}(\Omega)$ resp. $\mathbf{R}_{\mathrm{t}, \Gamma_{2}}(\Omega)$ it is enough to show equation (A.1) for $u \in \mathrm{C}_{\Gamma_{1}}^{\infty}(\Omega) \times \mathrm{C}_{\Gamma_{2}}^{\infty}(\Omega)$. Using Gauss's divergence theorem we calculate

$$
\begin{aligned}
\left\langle\Phi \Xi u, \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\mathrm{G}(\hat{r}, \tilde{r}))} & =\int_{\hat{r}}^{\tilde{r}} \phi(r)\left\langle\Xi u, \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\mathrm{~S}(r))} d r \\
& =\int_{\hat{r}}^{\tilde{r}} \phi(r)\left(\mu_{0}\langle\xi \times E, H\rangle_{\mathrm{L}^{2}(\mathrm{~S}(r))}-\varepsilon_{0}\langle\xi \times H, E\rangle_{\mathrm{L}^{2}(\mathrm{~S}(r))}\right) d r \\
& =\int_{\hat{r}}^{\tilde{r}} \phi(r) \int_{\Omega(r)}\left(\mu_{0} \operatorname{div}(E \times \bar{H})-\varepsilon_{0} \operatorname{div}(H \times \bar{E})\right) d \lambda^{3} d r
\end{aligned}
$$

Note, that

$$
\begin{aligned}
\mu_{0} & \operatorname{div}(E \times \bar{H})-\varepsilon_{0} \operatorname{div}(H \times \bar{E}) \\
\quad & =\mu_{0}(\bar{H} \operatorname{rot} E-E \operatorname{rot} \bar{H})-\varepsilon_{0}(\bar{E} \operatorname{rot} H-H \operatorname{rot} \bar{E}) \\
& =\left(\left(\mu_{0} \bar{H}\right) \operatorname{rot} E-\left(\varepsilon_{0} \bar{E}\right) \operatorname{rot} H\right)-\left(H \operatorname{rot}\left(\varepsilon_{0} \bar{E}\right)-E \operatorname{rot}\left(\mu_{0} \bar{H}\right)\right) \\
& =\overline{\Lambda_{0} u} \cdot \operatorname{Rot} u+u \cdot \operatorname{Rot} \overline{\Lambda_{0} u}
\end{aligned}
$$

and additionally $\operatorname{supp} u=\operatorname{supp} E \cap \operatorname{supp} H$, meaning $\operatorname{dist}(\operatorname{supp} u, \partial \Omega)>0$ and we may extend $\left.u\right|_{\operatorname{supp} u}$ by zero to the whole of $\mathrm{U}(\tilde{r})$. Thus (by abuse of notation)

$$
\left\langle\Phi \Xi u, \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\mathrm{G}(\hat{r}, \tilde{r}))}=\int_{\hat{r}}^{\tilde{r}} \phi(r)\left(\left\langle\operatorname{Rot} u, \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\mathrm{U}(r))}-\left\langle u, \operatorname{Rot} \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\mathrm{U}(r))}\right) d r
$$

and using Fubini's theorem we obtain:

$$
\begin{aligned}
\left\langle\Phi \Xi u, \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\mathrm{G}(\hat{r}, \tilde{r}))} & =\int_{\hat{r}}^{\tilde{r}} \phi(r)\left(\left\langle\operatorname{Rot} u, \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\Omega(r))}-\left\langle u, \operatorname{Rot} \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\Omega(r))}\right) d r \\
& =\int_{\hat{r}}^{\tilde{r}} \phi(r) \int_{0}^{r}\left(\left\langle\operatorname{Rot} u, \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\mathrm{~S}(\sigma))}-\left\langle u, \operatorname{Rot} \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\mathrm{~S}(\sigma))}\right) d \sigma d r \\
& =\int_{0}^{\tilde{r}} \int_{\max \{\hat{r}, \sigma\}}^{\tilde{r}} \phi(r)\left(\left\langle\operatorname{Rot} u, \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\mathrm{~S}(\sigma))}-\left\langle u, \operatorname{Rot} \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\mathrm{~S}(\sigma))}\right) d r d \sigma \\
& =\int_{0}^{\tilde{r}} \psi(\sigma)\left(\left\langle\operatorname{Rot} u, \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\mathrm{~S}(\sigma))}-\left\langle u, \operatorname{Rot} \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\mathrm{~S}(\sigma))}\right) d \sigma \\
& =\left\langle\Psi \operatorname{Rot} u, \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\Omega(\tilde{r}))}+\left\langle\Psi u, \operatorname{Rot} \Lambda_{0} u\right\rangle_{\mathrm{L}^{2}(\Omega(\tilde{r}))}
\end{aligned}
$$

We end this section with a Lemma, which will be needed to proof the polynomial decay and a-prioriestimate for the Helmholtz equation and can be shown by elementary partial integration.

Lemma A.4. Let $w \in \mathbf{H}_{\mathrm{loc}}^{2}\left(\mathbb{R}^{n}\right), 0 \notin \operatorname{supp} w, m \in \mathbb{R}$ and $\tilde{r}>0$. Then we have:
(1) $\operatorname{Re} \int_{U(\tilde{r})} r^{m+1} \Delta w \partial_{r} \bar{w}$

$$
=\frac{1}{2} \int_{\mathrm{U}(\tilde{r})} r^{m}\left((n+m-2)|\nabla w|^{2}-2 m\left|\partial_{r} w\right|^{2}\right)+\int_{\mathrm{S}(\tilde{r})} r^{m+1}\left(\left|\partial_{r} w\right|^{2}-\frac{1}{2}|\nabla w|^{2}\right),
$$

(2) $\operatorname{Re} \int_{\mathrm{U}(\tilde{r})} r^{m} \Delta w \bar{w}$

$$
=-\int_{\mathrm{U}(\tilde{r})} r^{m}\left(|\nabla w|^{2}-\frac{m}{2}(n+m-2) r^{-2}|w|^{2}\right)+\int_{\mathrm{S}(\tilde{r})} r^{m}\left(\operatorname{Re}\left(\partial_{r} w \bar{w}\right)-\frac{m}{2} r^{-1}|w|^{2}\right),
$$

(3) $\operatorname{Im} \int_{\mathrm{U}(\tilde{r})} r^{m} \Delta w \bar{w}=-m \int_{\mathrm{U}(\tilde{r})} r^{m-1} \operatorname{Im}\left(\partial_{r} w \bar{w}\right)+\frac{1}{2} \int_{\mathrm{S}(\tilde{r})} r^{m+1}|w|^{2}$,
(4) $\operatorname{Re} \int_{\mathrm{U}(\tilde{r})} r^{m+1} w \partial_{r} \bar{w}=-\frac{1}{2} \int_{\mathrm{U}(\tilde{r})} r^{m}(n+m)|w|^{2}+\frac{1}{2} \int_{\mathrm{S}(\tilde{r})} r^{m+1}|w|^{2}$,
where $\partial_{r}:=\xi \cdot \nabla$.

## Appendix B. Polynomial Decay and A-Priori Estimate for the Helmholtz Equation

In this section we present well known results for the Helmholtz equation, which we will use to achieve similar results for Maxwell's equations. We start with a regularity result (cf. [27, Lemma 4]) and the polynomial decay (cf. [27, Lemma 5]).

Lemma B.1. Let $t \in \mathbb{R}$. If $w \in \mathrm{~L}_{\mathrm{t}}^{2}\left(\mathbb{R}^{n}\right)$ and $\Delta w \in \mathrm{~L}_{\mathrm{t}}^{2}\left(\mathbb{R}^{n}\right)$ we have $w \in \mathbf{H}_{\mathrm{t}}^{2}\left(\mathbb{R}^{n}\right)$ and

$$
\|w\|_{\mathbf{H}_{t}^{2}\left(\mathbb{R}^{n}\right)} \leq c \cdot\left\{\|\Delta w\|_{L_{t}^{2}\left(\mathbb{R}^{n}\right)}+\|w\|_{L_{t}^{2}\left(\mathbb{R}^{n}\right)}\right\}
$$

with $c \in(0, \infty)$ independent of $w$ and $\Delta w$.
Proof. For $t=0$ we have $w, \Delta w \in \mathrm{~L}^{2}\left(\mathbb{R}^{n}\right)$ and using Fourier-Transformation we obtain

$$
\begin{align*}
& \|\Delta w\|_{\mathrm{L}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\|w\|_{\mathrm{L}^{2}\left(\mathbb{R}^{n}\right)}^{2} \\
& \quad=\left\|r^{2} \mathcal{F}(w)\right\|_{\mathrm{L}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\|\mathcal{F}(w)\|_{\mathrm{L}^{2}\left(\mathbb{R}^{n}\right)}^{2}  \tag{B.1}\\
& \quad=\int_{\mathbb{R}^{n}}\left(r^{4}+1\right)|\mathcal{F}(w)|^{2} \geq \frac{1}{2} \cdot\left\|\left(1+r^{2}\right) \mathcal{F}(w)\right\|_{\mathrm{L}^{2}\left(\mathbb{R}^{n}\right)}^{2},
\end{align*}
$$

yielding $w \in \mathbf{H}^{2}\left(\mathbb{R}^{n}\right)$ and the desired estimate. So let us switch to $t \neq 0$. Then, using a well known result concerning inner regularity (e.g., [3, Chapter VII, $\S 3.2$, Theorem 1]), we already have $w \in \mathbf{H}_{\text {loc }}^{2}\left(\mathbb{R}^{n}\right)$. Now let $\tilde{r}>1$ and define $\eta_{\tilde{r}} \in \dot{C}^{\infty}\left(\mathbb{R}^{n}\right)$ by $\eta_{\tilde{r}}(x):=\rho^{t} \eta(r(x) / \tilde{r})$. Then $\eta_{\tilde{r}} w \in \mathbf{H}^{2}\left(\mathbb{R}^{n}\right)$,

$$
\left|\nabla \eta_{\tilde{r}}\right| \leq c \cdot \rho^{t-1} \text { with } c=c(t)>0
$$

and

$$
\begin{aligned}
& \left\langle\nabla\left(\eta_{\tilde{r}} w\right), \nabla\left(\eta_{\tilde{r}} w\right)\right\rangle_{\mathrm{L}^{2}\left(\mathbb{R}^{n}\right)} \\
& =\operatorname{Re}\left\langle\nabla w, \nabla\left(\eta_{\tilde{r}}^{2} w\right)\right\rangle_{\mathrm{L}^{2}\left(\mathbb{R}^{n}\right)}+\left\|\left(\nabla \eta_{\tilde{r}}\right) w\right\|_{\mathrm{L}^{2}\left(\mathbb{R}^{n}\right)}^{2} \\
& \leq c \cdot\left\{\left\|\eta_{\tilde{r}} \Delta w\right\|_{\mathrm{L}^{2}\left(\mathbb{R}^{n}\right)}\left\|\eta_{\tilde{r}} w\right\|_{\mathrm{L}^{2}\left(\mathbb{R}^{n}\right)}+\|w\|_{\mathrm{L}_{t-1}^{2}\left(\mathbb{R}^{n}\right)}^{2}\right\} \\
& \quad \leq c \cdot\left\{\|\Delta w\|_{\mathrm{L}_{t}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\|w\|_{\mathrm{L}_{t}^{2}\left(\mathbb{R}^{n}\right)}^{2}\right\},
\end{aligned}
$$

with $c=c(n, t) \in(0, \infty)$, hence

$$
\|\nabla w\|_{\mathrm{L}_{\mathrm{t}}^{2}(\mathrm{~B}(\tilde{r}))}=\left\|\nabla\left(\eta_{\tilde{r}} w\right)-\left(\nabla \eta_{\tilde{r}}\right) w\right\|_{\mathrm{L}^{2}\left(\mathbb{R}^{n}\right)} \leq c(n, t) \cdot\left\{\|\Delta w\|_{\mathrm{L}_{\mathrm{t}}^{2}\left(\mathbb{R}^{n}\right)}+\|w\|_{\mathrm{L}_{\mathrm{t}}^{2}\left(\mathbb{R}^{n}\right)}\right\}
$$

Sending $\tilde{r} \longrightarrow \infty$ ( monotone convergence) shows $w \in \mathbf{H}_{\mathrm{t}}^{1}\left(\mathbb{R}^{n}\right)$ and

$$
\begin{equation*}
\|w\|_{\mathbf{H}_{t}^{1}\left(\mathbb{R}^{n}\right)} \leq c(n, t) \cdot\left\{\|\Delta w\|_{L_{t}^{2}\left(\mathbb{R}^{n}\right)}+\|w\|_{L_{t}^{2}\left(\mathbb{R}^{n}\right)}\right\} \tag{B.2}
\end{equation*}
$$

Moreover,

$$
\Delta\left(\rho^{t} w\right)=t\left(n+(t-2) \frac{r^{2}}{1+r^{2}}\right) \rho^{t-2} w+2 r \rho^{t-2} \partial_{r} w+\rho^{t} \Delta w
$$

such that with (B.2) we obtain

$$
\begin{equation*}
\left\|\Delta\left(\rho^{t} w\right)\right\|_{L^{2}\left(\mathbb{R}^{n}\right)} \leq c \cdot\left\{\|\Delta w\|_{L_{t}^{2}\left(\mathbb{R}^{n}\right)}+\|w\|_{L_{t}^{2}\left(\mathbb{R}^{n}\right)}\right\} \tag{B.3}
\end{equation*}
$$

with $c \in(0, \infty)$ independent of $w$ and $\Delta w$. Hence $\Delta\left(\rho^{t} w\right) \in \mathrm{L}^{2}\left(\mathbb{R}^{n}\right)$ and we may apply the first case. This shows $\rho^{t} w \in \mathbf{H}^{2}\left(\mathbb{R}^{n}\right)$ and using (B.1), (B.2) and (B.3), we obtain (uniformly w.r.t. $w$ and $\Delta w$ )

$$
\begin{aligned}
& \|w\|_{\mathbf{H}_{t}^{2}\left(\mathbb{R}^{n}\right)} \\
& \leq c \cdot\left\{\left\|\rho^{t} w\right\|_{\mathbf{H}^{2}\left(\mathbb{R}^{n}\right)}+\left\|\left(\nabla \rho^{t}\right) \nabla w\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}+\left\|\left(\nabla \rho^{t}\right) w\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}+\sum_{|\alpha|=2}\left\|\left(\partial^{\alpha} \rho^{t}\right) w\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}\right\} \\
& \quad \leq c \cdot\left\{\left\|\Delta\left(\rho^{t} w\right)\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}+\left\|\rho^{t} w\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}+\|\nabla w\|_{L_{t-1}^{2}\left(\mathbb{R}^{n}\right)}+\|w\|_{L_{\mathrm{t}-1}^{2}\left(\mathbb{R}^{n}\right)}\right\} \\
& \quad \leq c \cdot\left\{\|\Delta w\|_{L_{\mathrm{t}}^{2}\left(\mathbb{R}^{n}\right)}+\|w\|_{L_{t}^{2}\left(\mathbb{R}^{n}\right)}\right\}
\end{aligned}
$$

yielding $w \in \mathbf{H}_{\mathrm{t}}^{2}\left(\mathbb{R}^{n}\right)$ and the required estimate.

Lemma B. 2 (Polynomial decay). Let $J \Subset \mathbb{R} \backslash\{0\}$ some interval, $\gamma \in J$ and $s, t \in \mathbb{R}$ with $t>-1 / 2$ and $t \leq s$. If $w \in \mathrm{~L}_{\mathrm{t}}^{2}\left(\mathbb{R}^{n}\right)$ and $g:=\left(\Delta+\gamma^{2}\right) w \in \mathrm{~L}_{\mathrm{s}+1}^{2}\left(\mathbb{R}^{n}\right)$ it holds

$$
w \in \mathbf{H}_{\mathbf{s}}^{2}\left(\mathbb{R}^{n}\right) \quad \text { and } \quad\|w\|_{\mathbf{H}_{\mathbf{s}}^{2}\left(\mathbb{R}^{n}\right)} \leq c \cdot\left\{\|g\|_{L_{\mathrm{s}+1}^{2}\left(\mathbb{R}^{n}\right)}+\|w\|_{L_{\mathrm{s}-1}^{2}\left(\mathbb{R}^{n}\right)}\right\}
$$

with $c=c(n, s, J) \in(0, \infty)$ not depending on $\gamma, g$ or $w$.
Proof. The assertion follows directly from Lemma B.1, if we can show

$$
w \in \mathrm{~L}_{\mathrm{s}}^{2}\left(\mathbb{R}^{n}\right) \quad \text { with } \quad\|w\|_{\mathrm{L}_{\mathrm{s}}^{2}\left(\mathbb{R}^{n}\right)} \leq c \cdot\left\{\|g\|_{\mathrm{L}_{\mathrm{s}+1}^{2}\left(\mathbb{R}^{n}\right)}+\|w\|_{\mathrm{L}_{\mathrm{s}-1}^{2}\left(\mathbb{R}^{n}\right)}\right\}
$$

Therefore let $v:=\check{\chi} w$, where $\check{\chi} \in \mathrm{C}^{\infty}\left(\mathbb{R}^{n}\right)$ with $\check{\chi}=1$ on $\check{\mathrm{U}}(1)$ and vanishing in a neighbourhood of the origin. By assumption we already have $w \in \mathbf{H}_{\mathrm{t}}^{2}\left(\mathbb{R}^{n}\right)\left(\right.$ cf. Lemma B.1), hence $v \in \mathbf{H}_{\text {loc }}^{2}\left(\mathbb{R}^{n}\right)$ and we may apply the partial integration rules from Lemma A. 4 to

$$
\begin{aligned}
& \operatorname{Re} \int_{\mathrm{G}(\hat{r}, \tilde{r})}\left(\Delta w+\gamma^{2} w\right)\left(r^{2 t+1} \partial_{r} \bar{w}+\beta r^{2 t} \bar{w}\right) \\
& \quad=\operatorname{Re} \int_{\mathrm{G}(\hat{r}, \tilde{r})}\left(\Delta v+\gamma^{2} v\right)\left(r^{2 t+1} \partial_{r} \bar{v}+\beta r^{2 t} \bar{v}\right)=\ldots
\end{aligned}
$$

with $\tilde{r}>\hat{r} \geq 1$ and

$$
\beta:=\max \{(n-1) / 2, t+(n-1) / 2\} .
$$

After some rearrangements this leads to

$$
\begin{align*}
& \int_{\mathrm{G}(\hat{r}, \tilde{r})} r^{2 t}\left(\{\beta-(n+2 t-2) / 2\}|\nabla w|^{2}+\{(n+2 t) / 2-\beta\} \gamma^{2}|w|^{2}\right) \\
& +2 t \int_{\mathrm{G}(\hat{r}, \tilde{r})} r^{2 t}\left|\partial_{r} w\right|^{2}+\int_{\mathrm{S}(\tilde{r})} \tilde{r}^{2 t+1}|\nabla w|^{2} \\
& =-\operatorname{Re} \int_{\mathrm{G}(\hat{r}, \tilde{r})}\left(\Delta w+\gamma^{2} w\right)\left(r^{2 t+1} \partial_{r} \bar{w}+\beta r^{2 t} \bar{w}\right)+t(n+2 t-2) \beta \int_{\mathrm{G}(\hat{r}, \tilde{r})} r^{2 t-2}|w|^{2} \\
& \quad+\int_{\mathrm{S}(\hat{r})} \hat{r}^{2 t+1}\left(\beta t \hat{r}^{-2}|w|^{2}-\beta \hat{r}^{-1} \operatorname{Re}\left(\partial_{r} w \bar{w}\right)-\left|\partial_{r} w\right|^{2}\right)  \tag{B.4}\\
& \quad+\int_{\mathrm{S}(\tilde{r})} \tilde{r}^{2 t+1}\left(\left|\partial_{r} w\right|^{2}+\beta \tilde{r}^{-1} \operatorname{Re}\left(\partial_{r} w \bar{w}\right)-\beta t \tilde{r}^{-2}|w|^{2}\right) \\
& \quad+\frac{1}{2} \int_{\mathrm{S}(\tilde{r})} \tilde{r}^{2 t+1}\left(|\nabla w|^{2}+\gamma^{2}|w|^{2}\right)+\frac{1}{2} \int_{\mathrm{S}(\hat{r})} \hat{r}^{2 t+1}\left(|\nabla w|^{2}-\gamma^{2}|w|^{2}\right)
\end{align*}
$$

Let us first have a look on the left hand side of this equation. For $t \geq 0$ (i.e., $\beta=t+(n-1) / 2)$ we skip the second and third integral to obtain

$$
\begin{aligned}
& \int_{\mathrm{G}(\hat{r}, \tilde{r})} r^{2 t}\left(\{\beta-(n+2 t-2) / 2\}|\nabla w|^{2}+\{(n+2 t) / 2-\beta\} \gamma^{2}|w|^{2}\right) \\
& \quad+2 t \int_{\mathrm{G}(\hat{r}, \tilde{r})} r^{2 t}\left|\partial_{r} w\right|^{2}+\int_{\mathrm{S}(\tilde{r})} \tilde{r}^{2 t+1}|\nabla w|^{2} \\
& \quad \geq \frac{1}{2} \int_{\mathrm{G}(\hat{r}, \tilde{r})} r^{2 t}\left(\{2 \beta-(n+2 t-2)\}|\nabla w|^{2}+\{(n+2 t)-2 \beta\} \gamma^{2}|w|^{2}\right) \\
& \quad=\frac{1}{2} \int_{\mathrm{G}(\hat{r}, \tilde{r})} r^{2 t}\left(|\nabla w|^{2}+\gamma^{2}|w|^{2}\right)
\end{aligned}
$$

while in the case of $t<0$ (i.e., $\beta=(n-1) / 2)$ we just skip the third integral and end up with

$$
\begin{aligned}
& \int_{\mathrm{G}(\hat{r}, \tilde{r})} r^{2 t}\left(\{\beta-(n+2 t-2) / 2\}|\nabla w|^{2}+\{(n+2 t) / 2-\beta\} \gamma^{2}|w|^{2}\right) \\
& \quad+2 t \int_{\mathrm{G}(\hat{r}, \tilde{r})} r^{2 t}\left|\partial_{r} w\right|^{2}+\int_{\mathrm{S}(\tilde{r})} \tilde{r}^{2 t+1}|\nabla w|^{2} \\
& \quad \geq \int_{\mathrm{G}(\hat{r}, \tilde{r})} r^{2 t}\left(\{\beta-(n+2 t-2) / 2+2 t\}|\nabla w|^{2}+\{(n+2 t) / 2-\beta\} \gamma^{2}|w|^{2}\right) \\
& \quad=\left\{\frac{1}{2}+t\right\} \int_{\mathrm{G}(\hat{r}, \tilde{r})} r^{2 t}\left(|\nabla w|^{2}+\gamma^{2}|w|^{2}\right),
\end{aligned}
$$

since $\left|\partial_{r} w\right| \leq|\nabla w|$. Thus for arbitrary $t \in \mathbb{R}$ the left hand side of (B.4) can be estimated from below by

$$
\min \left\{\frac{1}{2}, \frac{1}{2}+t\right\} \int_{\mathrm{G}(\hat{r}, \tilde{r})} r^{2 t}\left(|\nabla w|^{2}+\gamma^{2}|w|^{2}\right)
$$

For the right hand side we have ( $\tilde{r}>1$ )

$$
\begin{aligned}
& \int_{\mathrm{S}(\tilde{r})} \tilde{r}^{2 t+1}\left(\left|\partial_{r} w\right|^{2}+\beta \tilde{r}^{-1} \operatorname{Re}\left(\partial_{r} w \bar{w}\right)-\beta t \tilde{r}^{-2}|w|^{2}\right) \\
& \quad \leq \int_{\mathrm{S}(\tilde{r})} \tilde{r}^{2 t+1}\left(\left|\partial_{r} w\right|^{2}+\beta\left|\partial_{r} w \bar{w}\right|+\beta|t||w|^{2}\right) \leq c \cdot \int_{\mathrm{S}(\tilde{r})} \tilde{r}^{2 t+1}\left(|\nabla w|^{2}+|w|^{2}\right)
\end{aligned}
$$

as well as

$$
\begin{aligned}
& \int_{\mathrm{S}(\hat{r})} \hat{r}^{2 t+1}\left(\beta t \hat{r}^{-2}|w|^{2}-\beta \hat{r}^{-1} \operatorname{Re}\left(\partial_{r} w \bar{w}\right)-\left|\partial_{r} w\right|^{2}\right) \\
& \quad \leq \int_{\mathrm{S}(\tilde{r})} \tilde{r}^{2 t+1}\left(\beta|t| \hat{r}^{-2}|w|^{2}+\beta \hat{r}^{-1}\left|\partial_{r} w \bar{w}\right|\right) \leq c \cdot \int_{\mathrm{S}(\tilde{r})} \tilde{r}^{2 t+1}\left(\hat{r}^{-2}|w|^{2}+|\nabla w|^{2}\right)
\end{aligned}
$$

such that equation (B.4) becomes

$$
\begin{aligned}
\min \{ & \left.\frac{1}{2}, \frac{1}{2}+t\right\} \int_{\mathrm{G}(\hat{r}, \tilde{r})} r^{2 t}\left(|\nabla w|^{2}+\gamma^{2}|w|^{2}\right) \\
\leq & \int_{\mathrm{G}(\hat{r}, \tilde{r})} r^{t+1}|g|\left(r^{t}|\nabla w|+\beta r^{t-1}|w|\right)+\beta|t(n+2 t-2)| \int_{\mathrm{G}(\hat{r}, \tilde{r})} r^{2 t-2}|w|^{2} \\
& +c(n, t) \cdot\left\{\int_{\mathrm{S}(\hat{r})} \hat{r}^{2 t+1}\left(\hat{r}^{-2}|w|^{2}+|\nabla w|^{2}-\gamma^{2}|w|^{2}\right)+\int_{\mathrm{S}(\tilde{r})} \tilde{r}^{2 t+1}\left(|\nabla w|^{2}+|w|^{2}\right)\right\} .
\end{aligned}
$$

By assumption we have $w \in \mathbf{H}_{\mathrm{t}}^{2}\left(\mathbb{R}^{n}\right)$, such that according to Lemma A. 2 the lower limit for $\tilde{r} \longrightarrow \infty$ of the last boundary integral vanishes. Hence we may replace $\mathrm{G}(\hat{r}, \tilde{r})$ by $\check{\mathrm{U}}(\hat{r})$ and in addition use Young's
inequality to obtain

$$
\begin{align*}
& \left\|r^{t} \nabla w\right\|_{\mathrm{L}^{2}(\check{\mathrm{U}}(\hat{r}))}^{2}+\gamma^{2}\left\|r^{t} w\right\|_{\mathrm{L}^{2}(\check{\mathrm{U}}(\hat{r}))}^{2} \\
& \quad \leq c(n, t) \cdot\left\{\left\|r^{t+1} g\right\|_{\mathrm{L}^{2}(\check{\mathrm{U}}(\hat{r}))}^{2}+\left\|r^{t-1} w\right\|_{\mathrm{L}^{2}(\check{\mathrm{U}}(\hat{r}))}^{2}+\int_{\mathrm{S}(\hat{r})} \hat{r}^{2 t+1}\left(|\nabla w|^{2}-\gamma^{2}|w|^{2}+\hat{r}^{-2}|w|^{2}\right)\right\} \mathrm{B} . \\
& \quad \leq c(n, t) \cdot\left\{\|g\|_{\mathrm{L}_{\mathrm{t}+1}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\|w\|_{\mathrm{L}_{\mathrm{t}-1}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\int_{\mathrm{S}(\hat{r})} \hat{r}^{2 t+1}\left(|\nabla w|^{2}-\gamma^{2}|w|^{2}+\hat{r}^{-2}|w|^{2}\right)\right\} .
\end{align*}
$$

Now suppose first that $s=t$. Then the assertion simply follows by choosing $\hat{r}:=1$ as the trace theorem bounds the surface integral by $\|w\|_{\mathbf{H}^{2}(\mathrm{U}(1))}^{2}$ and therefore (using Lemma B.1)

$$
\begin{aligned}
&\|w\|_{\mathbf{H}_{t}^{1}\left(\mathbb{R}^{n}\right)} \leq c(n, s, J) \cdot\left\{\|g\|_{L_{t+1}^{2}\left(\mathbb{R}^{n}\right)}+\|w\|_{L_{t-1}^{2}\left(\mathbb{R}^{n}\right)}+\|w\|_{\mathbf{H}^{2}(\mathrm{U}(2))}\right\} \\
& \leq c(n, s, J) \cdot\left\{\|g\|_{\mathrm{L}_{t+1}^{2}\left(\mathbb{R}^{n}\right)}+\|w\|_{\mathrm{L}_{\mathrm{t}-1}^{2}\left(\mathbb{R}^{n}\right)}+\|w\|_{\mathbf{H}_{\mathrm{t}-1}^{2}\left(\mathbb{R}^{n}\right)}\right\} \\
& \leq c(n, s, J) \cdot\left\{\|g\|_{\mathrm{L}_{\mathrm{t}+1}^{2}\left(\mathbb{R}^{n}\right)}+\|w\|_{\mathrm{L}_{\mathrm{t}-1}^{2}\left(\mathbb{R}^{n}\right)}+\|\Delta w\|_{\mathrm{L}_{-1}^{2}\left(\mathbb{R}^{n}\right)}\right\} \\
& \leq c(n, s, J) \cdot\left\{\|g\|_{\mathrm{L}_{t+1}^{2}\left(\mathbb{R}^{n}\right)}+\|w\|_{\mathrm{L}_{t-1}^{2}\left(\mathbb{R}^{n}\right)}\right\} .
\end{aligned}
$$

For the case $w \notin \mathrm{~L}_{\mathrm{s}}^{2}\left(\mathbb{R}^{n}\right)$ let $\hat{s}:=\sup \left\{m \in \mathbb{R} \mid u \in \mathrm{~L}_{\mathrm{m}}^{2}\left(\mathbb{R}^{n}\right)\right\}$. Then, w.l.o.g. ${ }^{\text {iv }}$, we may assume

$$
\hat{s}-1 / 2<t<\hat{s}<s \leq t+1 / 2
$$

hence $\delta:=1-2(s-t) \in(0,1)$. Multiplying (B.5) with $\hat{r}^{-\delta}$ and integrating from 1 to some $\check{r}>1$ leads to:

$$
\begin{align*}
& \int_{1}^{\check{r}} \hat{r}^{-\delta} \int_{\check{\mathrm{U}}(\hat{r})} r^{2 t}\left(|\nabla w|^{2}+\gamma^{2}|w|^{2}\right) d \lambda^{n} d \hat{r} \\
& \quad \leq c(n, t) \cdot\left\{\int_{1}^{\check{r}} \hat{r}^{-\delta} \int_{\check{\mathrm{U}}(\hat{r})} r^{2 t+2}|g|^{2}+r^{2 t-2}|w|^{2} d \lambda^{n} d \hat{r}\right.  \tag{B.6}\\
&\left.+\int_{\mathrm{G}(1, \check{r})} r^{2 t+1-\delta}\left(|\nabla w|^{2}-\gamma^{2}|w|^{2}+r^{-2}|w|^{2}\right) d \lambda^{n}\right\}
\end{align*}
$$

By Fubini's theorem we have for arbitrary $h \in \mathrm{~L}^{1}\left(\mathbb{R}^{n}\right)$

$$
\begin{aligned}
\int_{1}^{\check{r}} & \hat{r}^{-\delta} \int_{\check{\mathrm{U}}(\hat{r})} h d \lambda^{n} d \hat{r} \\
& =\int_{1}^{\check{r}} \int_{\hat{r}}^{\infty} \int_{\mathrm{S}(\sigma)} \hat{r}^{-\delta} h d \lambda_{s}^{n-1} d \sigma d \hat{r} \\
& =\int_{1}^{\infty} \int_{1}^{\min \{\sigma, \check{r}\}} \hat{r}^{-\delta} \int_{\mathrm{S}(\sigma)} h d \lambda_{s}^{n-1} d \hat{r} d \sigma \\
& =\int_{1}^{\infty}(1-\delta)^{-1} \min \left\{\sigma^{1-\delta}-1, \check{r}^{1-\delta}-1\right\} \int_{\mathrm{S}(\sigma)} h d \lambda_{s}^{n-1} d \sigma \\
& =\int_{1}^{\infty} \int_{\mathrm{S}(\sigma)} \underbrace{(1-\delta)^{-1} \min \left\{r^{1-\delta}-1, \check{r}^{1-\delta}-1\right\}}_{=: \theta_{\tilde{r}}} h d \lambda_{s}^{n-1} d \sigma=\int_{\check{\mathrm{U}}(1)} \theta_{\check{r}} h d \lambda^{n}
\end{aligned}
$$

[^3]such that (B.6) becomes ( note that $\theta_{\check{r}} \leq(1-\delta)^{-1} \cdot r^{1-\delta}$ and $1-\delta=2(s-t)$ )
\[

$$
\begin{aligned}
& \int_{\check{\mathrm{U}}(1)} \theta_{\check{r}} r^{2 t}\left(|\nabla w|^{2}+\gamma^{2}|w|^{2}\right) \\
& \left.\quad \leq c(n, t) \cdot\left\{\int_{\check{\mathrm{U}}(1)} \theta_{\check{r}}\left(r^{2 t+2}|g|^{2}+r^{2 t-2}|w|^{2}\right)+\int_{\mathrm{G}(1, \check{r})} r^{2 t+1-\delta}\left(|\nabla w|^{2}-\gamma^{2}|w|^{2}+r^{-2}|w|^{2}\right)\right\} \text { B.7 }\right) \\
& \quad \leq c(n, s) \cdot\left\{\|g\|_{\mathrm{L}_{\mathrm{s}+1}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\|w\|_{\mathrm{L}_{\mathrm{s}-1}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\int_{\mathrm{G}(1, \check{r})} r^{2 s}\left(|\nabla w|^{2}-\gamma^{2}|w|^{2}\right) d \lambda^{n}\right\} .
\end{aligned}
$$
\]

Finally, look at

$$
\operatorname{Re} \int_{\mathrm{G}(1, \check{r})} r^{2 s} g \bar{w}=\operatorname{Re} \int_{\mathrm{G}(1, \check{r})} r^{2 s} g \bar{v}
$$

Applying Lemma A. 3 we obtain (after some rearrangements)

$$
\begin{aligned}
\int_{\mathrm{G}(1, \check{r})} & r^{2 s}\left(|\nabla w|^{2}-\gamma^{2}|w|^{2}\right) \\
= & -\operatorname{Re} \int_{\mathrm{G}(1, \check{r})} r^{2 s} g \bar{w}+s(n+2 s-2) \int_{\mathrm{G}(1, \check{r})} r^{2 s-2}|w|^{2} \\
& +\int_{\mathrm{S}(\check{r})} \check{r}^{2 s}\left(\operatorname{Re}\left(\partial_{r} w \bar{w}\right)-s \check{r}^{-1}|w|^{2}\right)-\int_{S(1)}\left(\operatorname{Re}\left(\partial_{r} w \bar{w}\right)-s|w|^{2}\right) \\
& \leq c(n, s) \cdot\left\{\int_{\mathrm{G}(1, \check{r})}\left(r^{2 s+2}|g|^{2}+r^{2 s-2}|w|^{2}\right)+\int_{S(1)}\left(|\nabla w|^{2}+|w|^{2}\right)+\int_{\mathrm{S}(\check{r})} \check{r}^{2 s}\left(|\nabla w|^{2}+|w|^{2}\right)\right\},
\end{aligned}
$$

hence (using the trace theorem and Lemma B.1)

$$
\begin{align*}
& \int_{\mathrm{G}(1, \check{r})} r^{2 s}\left(|\nabla w|^{2}-\gamma^{2}|w|^{2}\right) \\
& \quad \leq c(n, s) \cdot\left\{\|g\|_{\mathrm{L}_{\mathrm{s}+1}\left(\mathbb{R}^{n}\right)}^{2}+\|w\|_{\mathrm{L}_{\mathrm{s}-1}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\|w\|_{\mathbf{H}^{2}(\mathrm{U}(1))}^{2}+\int_{\mathrm{S}(\check{r})} \check{r}^{2 s}\left(|\nabla w|^{2}+|w|^{2}\right)\right\}  \tag{B.8}\\
& \quad \leq c(n, s, J) \cdot\left\{\|g\|_{\mathrm{L}_{\mathrm{s}+1}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\|w\|_{\mathrm{L}_{\mathrm{s}-1}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\int_{\mathrm{S}(\check{r})} \check{r}^{2 s}\left(|\nabla w|^{2}+|w|^{2}\right)\right\}
\end{align*}
$$

and inserting (B.8) into (B.7) we end up with

$$
\begin{aligned}
\int_{\breve{\mathrm{U}}(1)} \theta_{\check{r}} & r^{2 t}\left(|\nabla w|^{2}+\gamma^{2}|w|^{2}\right) \\
& \leq c(n, s, J) \cdot\left\{\|g\|_{L_{\mathrm{s}+1}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\|w\|_{\mathrm{L}_{\mathrm{s}-1}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\int_{\mathrm{S}(\check{r})} \check{r}^{2 s}\left(|\nabla w|^{2}+|w|^{2}\right)\right\} .
\end{aligned}
$$

Again the lower limit for $\check{r} \longrightarrow \infty$ of the boundary integral vanishes (cf. Lemma A. 2 and observe that $w \in \mathbf{H}_{\mathrm{s}-\frac{1}{2}}^{2}\left(\mathbb{R}^{n}\right)$, since $0<s-t \leq 1 / 2$ by assumption), such that passing to the limit on a suitable subsequence we obtain

$$
\begin{gathered}
\|w\|_{\mathrm{L}_{\mathrm{s}}^{2}\left(\mathbb{R}^{n}\right)}^{2} \leq c(n, s, J) \cdot\left\{\int_{\check{\mathrm{U}}(1)}(1-\delta)^{-1} r^{2 t+1-\delta}\left(|\nabla w|^{2}+\gamma^{2}|w|^{2}\right)+\|w\|_{\mathrm{L}_{\mathrm{s}}^{2}(\mathrm{U}(1))}^{2}\right\} \\
\leq c(n, s, J) \cdot\left\{\lim _{\check{r} \rightarrow \infty} \int_{\check{\mathrm{U}}(1)} \theta_{\check{r}} r^{2 t}\left(|\nabla w|^{2}+\gamma^{2}|w|^{2}\right)+\|w\|_{\mathrm{L}_{\mathrm{s}-1}^{2}(\mathrm{U}(1))}^{2}\right\} \\
\leq c(n, s, J) \cdot\left\{\|g\|_{\mathrm{L}_{\mathrm{s}+1}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\|w\|_{\mathrm{L}_{\mathrm{s}-1}^{2}\left(\mathbb{R}^{n}\right)}^{2}\right\},
\end{gathered}
$$

showing $w \in \mathrm{~L}_{\mathrm{s}}^{2}\left(\mathbb{R}^{n}\right)$ and the required estimate.

Lemma B. 3 (A-priori estimate). Let $n \in \mathbb{N}, t<-1 / 2,1 / 2<s<1$ and $J \Subset \mathbb{R} \backslash\{0\}$ an interval. Then there exists $c, \delta \in(0, \infty)$, such that for all $\beta \in \mathbb{C}_{+}$with $\beta^{2}=\nu^{2}+i \nu \tau, \nu \in J, \tau \in(0,1]$ and $g \in \mathrm{~L}_{\mathbf{s}}^{2}\left(\mathbb{R}^{n}\right)$

$$
\begin{align*}
&\left\|(\Delta+\beta)^{-1} g\right\|_{L_{\mathrm{t}}^{2}\left(\mathbb{R}^{n}\right)}+\| \exp ( -i \nu r)(\Delta+\beta)^{-1} g \|_{\mathrm{H}_{\mathrm{s}-2}^{1}\left(\mathbb{R}^{n}\right)} \\
& \leq c \cdot\left\{\|g\|_{\mathrm{L}_{\mathrm{s}}^{2}\left(\mathbb{R}^{n}\right)}+\left\|(\Delta+\beta)^{-1} g\right\|_{\mathrm{L}^{2}(\Omega(\delta))}\right\} \tag{B.9}
\end{align*}
$$

holds.
Ikebe and Saito [6] proved this estimate for the space dimension $n=3$ and with $t=-s$, which already shows the result also for any $t<-1 / 2$ as the norms depend monotonic on the parameters $s$ and $t$. For arbitrary space dimensions we follow the proof of Vogelsang [21, Satz 4].

Proof. First of all, observe that

$$
\Delta: \mathbf{H}^{2}\left(\mathbb{R}^{n}\right) \subset \mathrm{L}^{2}\left(\mathbb{R}^{n}\right) \longrightarrow \mathrm{L}^{2}\left(\mathbb{R}^{n}\right), w \longmapsto \Delta w
$$

is self-adjoint and therefore $w:=(\Delta+\beta)^{-1} g \in \mathbf{H}^{2}\left(\mathbb{R}^{n}\right)$ well defined. Moreover, as mentioned above, due to the monotonic dependence of the norms on the parameters $s$ and $t$, it is enough to concentrate on the case $t=-s$. With $w_{e}:=\exp (-i \nu r) w$ and $g_{e}:=\exp (-i \nu r) g$, we have $w_{e} \in \mathbf{H}^{2}(\Omega)$ and

$$
\Delta w_{e}+i \nu\left(\tau w_{e}+\frac{n-1}{r} w_{e}+2 \partial_{r} w_{e}\right)=g_{e} .
$$

Applying Lemma A. 4 to

$$
\operatorname{Re} \int_{\mathrm{G}(1, \tilde{r})} g_{e}\left(r^{2 s-1} \partial_{r} \bar{w}_{e}+\frac{n-1}{2} r^{2 s-2} \bar{w}_{e}+\frac{\tau}{2} r^{2 s-1} \bar{w}_{e}\right)=\ldots
$$

with $\tilde{r}>1$ and using the same techniques as in the proof of Lemma B. 2 we obtain

$$
\begin{aligned}
& \frac{1}{2} \int_{\mathrm{G}(1, \tilde{r})} r^{2 s-2}\left((4 s-4)\left|\partial_{r} w_{e}\right|^{2}-(2 s-3)\left|\nabla w_{e}\right|^{2}\right)+\frac{1}{2} \int_{\mathrm{G}(1, \tilde{r})} r^{2 s-1} \tau\left|\nabla w_{e}\right|^{2} \\
& =-\operatorname{Re} \int_{\mathrm{G}(1, \tilde{r})} g_{e}\left(r^{2 s-1} \partial_{r} \bar{w}_{e}+\frac{n-1}{2} r^{2 s-2} \bar{w}_{e}+\frac{\tau}{2} r^{2 s-1} \bar{w}_{e}\right) \\
& \quad+\frac{n-1}{2}(s-1)(n+2 s-4) \int_{\mathrm{G}(1, \tilde{r})} r^{2 s-4}\left|w_{e}\right|^{2}+\frac{\tau}{4}(2 s-1)(n+2 s-3) \int_{\mathrm{G}(1, \tilde{r})} r^{2 s-3}\left|w_{e}\right|^{2} \\
& \quad+\frac{1}{2} \int_{\mathrm{S}(\tilde{r})} \tilde{r}^{2 s-1}\left(2\left|\partial_{r} w_{e}\right|^{2}+\tau \operatorname{Re}\left(\partial_{r} w_{e} \bar{w}_{e}\right)-\left|\nabla w_{e}\right|^{2}-\frac{(2 s-1)}{\tilde{r}} \tau\left|w_{e}\right|^{2}\right) \\
& \quad+\frac{(n-1)}{2} \int_{\mathrm{S}(\tilde{r})} \tilde{r}^{2 s-2}\left(\operatorname{Re}\left(\partial_{r} w_{e} \bar{w}_{e}\right)-\frac{s-1}{\tilde{r}}\left|w_{e}\right|^{2}\right) \\
& \quad-\frac{1}{2} \int_{\mathrm{S}(1)}\left(2\left|\partial_{r} w_{e}\right|^{2}+\tau \operatorname{Re}\left(\partial_{r} w_{e} \bar{w}_{e}\right)-\left|\nabla w_{e}\right|^{2}-(2 s-1) \tau\left|w_{e}\right|^{2}\right) \\
& \quad-\frac{(n-1)}{2} \int_{\mathrm{S}(1)}\left(\operatorname{Re}\left(\partial_{r} w_{e} \bar{w}_{e}\right)-(s-1)\left|w_{e}\right|^{2}\right)
\end{aligned}
$$

Since $4 s-4<0$ and $\left|\partial_{r} w_{e}\right| \leq\left|\nabla w_{e}\right|$, the left hand side can be estimated from below

$$
\begin{aligned}
& \frac{1}{2} \int_{\mathrm{G}(1, \tilde{r})} r^{2 s-2}\left((4 s-4)\left|\partial_{r} w_{e}\right|^{2}-(2 s-3)\left|\nabla w_{e}\right|^{2}\right)+\frac{1}{2} \int_{\mathrm{G}(1, \tilde{r})} r^{2 s-1} \tau\left|\nabla w_{e}\right|^{2} \\
& \quad \geq \frac{1}{2} \int_{\mathrm{G}(1, \tilde{r})} r^{2 s-2}((4 s-4)-(2 s-3))\left|\nabla w_{e}\right|^{2}=\left(s-\frac{1}{2}\right) \int_{\mathrm{G}(1, \tilde{r})} r^{2 s-2}\left|\nabla w_{e}\right|^{2}
\end{aligned}
$$

while for the right hand side we obtain

$$
\begin{aligned}
& -\operatorname{Re} \int_{\mathrm{G}(1, \tilde{r})} g_{e}\left(r^{2 s-1} \partial_{r} \bar{w}_{e}+\frac{n-1}{2} r^{2 s-2} \bar{w}_{e}+\frac{\tau}{2} r^{2 s-1} \bar{w}_{e}\right)+\ldots \\
& \leq \int_{\mathrm{G}(1, \tilde{r})} r^{s}\left|g_{e}\right|\left(r^{s-1}\left|\nabla w_{e}\right|+\frac{n-1}{2} r^{s-2}\left|w_{e}\right|+\frac{\tau}{2} r^{s-1}\left|w_{e}\right|\right) \\
& +c \cdot\left\{\int_{\mathrm{G}(1, \tilde{r})} r^{2 s-4}\left|w_{e}\right|^{2}+\tau \int_{\mathrm{G}(1, \tilde{r})} r^{s-2}\left|w_{e}\right| r^{s-1}\left|w_{e}\right|\right. \\
& \left.+\int_{\mathrm{S}(1)}\left(\left|\nabla w_{e}\right|^{2}+\left|\partial_{r} w_{e} \bar{w}_{e}\right|+\left|w_{e}\right|^{2}\right)+\int_{\mathrm{S}(\tilde{r})} \tilde{r}^{2 s-1}\left(\left|\nabla w_{e}\right|^{2}+\left|\partial_{r} w_{e} \bar{w}_{e}\right|+\left|w_{e}\right|^{2}\right)\right\},
\end{aligned}
$$

yielding

$$
\begin{aligned}
& \left(s-\frac{1}{2}\right) \int_{\mathrm{G}(1, \tilde{r})} r^{2 s-2}\left|\nabla w_{e}\right|^{2} \\
& \quad \leq \int_{\mathrm{G}(1, \tilde{r})} r^{s}\left|g_{e}\right|\left(r^{s-1}\left|\nabla w_{e}\right|+\frac{n-1}{2} r^{s-2}\left|w_{e}\right|+\frac{\tau}{2} r^{s-1}\left|w_{e}\right|\right)+c \cdot\left\{\int_{\mathrm{G}(1, \tilde{r})} r^{2 s-4}\left|w_{e}\right|^{2}\right. \\
& \quad \\
& \left.\quad+\tau \int_{\mathrm{G}(1, \tilde{r})} r^{2 s-2}\left|w_{e}\right|+\int_{\mathrm{S}(1)}\left(\left|\nabla w_{e}\right|^{2}+\left|w_{e}\right|^{2}\right)+\int_{\mathrm{S}(\tilde{r})} \tilde{r}^{2 s-1}\left(\left|\nabla w_{e}\right|^{2}+\left|w_{e}\right|^{2}\right)\right\}
\end{aligned}
$$

In the above, as well as in the sequel, $c \in(0, \infty)$ denotes a generic constant independent of $\nu, \tau, w$ and $g$. According to Lemma A. 2 the lower limit for $\tilde{r} \longrightarrow \infty$ of the last boundary integral vanishes. Thus we may omit it and replace $\mathrm{G}(1, \tilde{r})$ by $\check{\mathrm{U}}(1)$, such that using Young's inequality we end up with

$$
\begin{aligned}
& \left\|r^{s-1} \nabla w_{e}\right\|_{\mathrm{L}^{2}(\check{\mathrm{U}}(1))}^{2} \\
& \quad \leq c \cdot\left\{\left\|r^{s} g_{e}\right\|_{\mathrm{L}^{2}(\check{\mathrm{U}}(1))}^{2}+\tau\left\|r^{s-1} w_{e}\right\|_{\mathrm{L}^{2}(\check{\mathrm{U}}(1))}^{2}+\left\|r^{s-2} w_{e}\right\|_{\mathrm{L}^{2}(\check{\mathrm{U}}(1))}^{2}+\int_{\mathrm{S}(1)}\left(\left|\nabla w_{e}\right|^{2}+\left|w_{e}\right|^{2}\right)\right\} \\
& \quad \leq c \cdot\left\{\left\|g_{e}\right\|_{\mathrm{L}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\tau\left\|w_{e}\right\|_{\mathrm{L}_{\mathrm{s}-1}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\left\|w_{e}\right\|_{\mathrm{L}_{\mathrm{s}-2}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\int_{\mathrm{S}(1)}\left(\left|\nabla w_{e}\right|^{2}+\left|w_{e}\right|^{2}\right)\right\} .
\end{aligned}
$$

In addition the surface integral is bounded by $\left\|w_{e}\right\|_{\mathbf{H}^{2}(\mathrm{U}(1))}^{2}$ ( trace theorem) and Lemma B. 1 yields,

$$
\left\|w_{e}\right\|_{\mathbf{H}^{2}(\mathrm{U}(2))} \leq c \cdot\left\|w_{e}\right\|_{\mathbf{H}_{-\mathrm{s}}^{2}\left(\mathbb{R}^{n}\right)} \leq c \cdot\left\{\left\|g_{e}\right\|_{\mathrm{L}_{\mathbf{s}}^{2}\left(\mathbb{R}^{n}\right)}+\left\|w_{e}\right\|_{\mathrm{L}_{-\mathrm{s}}^{2}\left(\mathbb{R}^{n}\right)}\right\}
$$

showing

$$
\begin{aligned}
\left\|\nabla w_{e}\right\|_{\mathrm{L}_{\mathrm{s}-1}^{2}\left(\mathbb{R}^{n}\right)}^{2} \leq c \cdot\left\{\|g\|_{\mathrm{L}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\tau\left\|w_{e}\right\|_{\mathrm{L}_{\mathrm{s}-1}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\left\|w_{e}\right\|_{\mathrm{L}_{\mathrm{s}-2}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\left\|w_{e}\right\|_{\mathbf{H}^{2}(\mathrm{U}(1))}^{2}\right. \\
\leq c \cdot\left\{\|g\|_{\mathrm{L}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\tau\|w\|_{\mathrm{L}_{\mathrm{s}-1}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\|w\|_{\mathrm{L}_{\mathrm{s}-2}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\|w\|_{\mathrm{L}_{-\mathrm{s}}^{2}\left(\mathbb{R}^{n}\right)}^{2}\right\},
\end{aligned}
$$

Now, the differential equation implies

$$
\|g\|_{\mathbf{L}^{2}\left(\mathbb{R}^{n}\right)}\|w\|_{\mathbf{L}^{2}\left(\mathbb{R}^{n}\right)} \geq\left|\operatorname{Im}\langle g, w\rangle_{\mathbf{L}^{2}\left(\mathbb{R}^{n}\right)}\right|=\left|i \nu \tau\langle w, w\rangle_{\mathbf{L}^{2}\left(\mathbb{R}^{n}\right)}\right|=\tau|\nu|\|w\|_{\mathbf{L}^{2}\left(\mathbb{R}^{n}\right)}^{2}
$$

hence $(-s>s-2)$

$$
\begin{align*}
\|\exp (-i \nu r) w\|_{\mathrm{H}_{\mathrm{s}-2}^{1}\left(\mathbb{R}^{n}\right)} \leq c \cdot\left\{\left\|w_{e}\right\|_{\mathrm{L}_{\mathrm{s}-2}^{2}\left(\mathbb{R}^{n}\right)}+\left\|\nabla w_{e}\right\|_{\mathrm{L}_{\mathrm{s}-1}^{2}\left(\mathbb{R}^{n}\right)}\right\} \\
\leq c \cdot\left\{\|g\|_{L_{s}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\tau\|w\|_{\mathrm{L}_{s-1}^{2}\left(\mathbb{R}^{n}\right)}^{2}+\|w\|_{\mathrm{L}_{-\mathrm{s}}^{2}\left(\mathbb{R}^{n}\right)}^{2}\right\}  \tag{B.10}\\
\leq c \cdot\left\{\|g\|_{\mathrm{L}_{\mathrm{s}}^{2}\left(\mathbb{R}^{n}\right)}+\|w\|_{\mathrm{L}_{-\mathrm{s}}^{2}\left(\mathbb{R}^{n}\right)}\right\},
\end{align*}
$$

and it remains to estimate $\|w\|_{L_{-s}^{2}\left(\mathbb{R}^{n}\right)}$. For that we calculate

$$
\operatorname{Im} \int_{\mathrm{G}(1, \tilde{r})} g_{e} \bar{w}_{e}=\operatorname{Im} \int_{\mathrm{G}(1, \tilde{r})} \Delta w_{e} \bar{w}_{e}+\int_{\mathrm{G}(1, \tilde{r})} \nu\left(\tau w_{e}+\frac{n-1}{r} w_{e}\right) \bar{w}_{e}+2 \nu \operatorname{Re} \int_{\mathrm{G}(1, \tilde{r})} \partial_{r} w_{e} \bar{w}_{e}=\ldots
$$

using Lemma A. 4 and obtain

$$
\begin{aligned}
& \nu \int_{\mathrm{G}(1, \tilde{r})} r^{-2 s}\left((2 s-1)\left|w_{e}\right|^{2}+\tau r\left|w_{e}\right|^{2}\right) \\
&= \operatorname{Im} \int_{\mathrm{G}(1, \tilde{r})} r^{1-2 s} g_{e} \bar{w}_{e}-(2 s-1) \int_{\mathrm{G}(1, \tilde{r})} r^{-2 s} \operatorname{Im}\left(\partial_{r} w_{e} \bar{w}_{e}\right) \\
&+\int_{\mathrm{S}(\tilde{r})} r^{1-2 s}\left(\tau\left|w_{e}\right|^{2}+\operatorname{Im}\left(\partial_{r} w_{e} \bar{w}_{e}\right)\right)-\int_{\mathrm{S}(1)}\left(\tau\left|w_{e}\right|^{2}+\operatorname{Im}\left(\partial_{r} w_{e} \bar{w}_{e}\right)\right) \\
& \leq \int_{\mathrm{G}(1, \tilde{r})} r^{s}\left|g_{e}\right| r^{1-3 s}\left|w_{e}\right|+(2 s-1) \int_{\mathrm{G}(1, \tilde{r})} r^{s-1}\left|\partial_{r} w_{e}\right| r^{1-3 s}\left|w_{e}\right| \\
&+c \cdot\left\{\int_{\mathrm{S}(\tilde{r})} r^{1-2 s}\left(\left|w_{e}\right|^{2}+\left|\partial_{r} w_{e} \bar{w}_{e}\right|\right)+\int_{\mathrm{S}(1)}\left(\left|w_{e}\right|^{2}+\left|\partial_{r} w_{e} \bar{w}_{e}\right|\right)\right\} \\
& \leq( \left.\left\|r^{s} g_{e}\right\|_{\mathrm{L}^{2}(\mathrm{G}(1, \tilde{r}))}+(2 s-1)\left\|r^{s-1} \nabla w_{e}\right\|_{\mathrm{L}^{2}(\mathrm{G}(1, \tilde{r}))}\right) \cdot\left\|r^{1-3 s} w_{e}\right\|_{\mathrm{L}^{2}(\mathrm{G}(1, \tilde{r}))} \\
&+c \cdot\left\{\int_{\mathrm{S}(1)}\left(\left|w_{e}\right|^{2}+\left|\nabla w_{e}\right|^{2}\right)+\int_{\mathrm{S}(\tilde{r})} \tilde{r}^{1-2 s}\left(\left|w_{e}\right|^{2}+\left|\nabla w_{e}\right|^{2}\right)\right\} .
\end{aligned}
$$

As before the lower limit for $\tilde{r} \longrightarrow \infty$ of the last boundary integral vanishes (cf. Lemma A. 2 and observe that $\left.w_{e} \in \mathbf{H}^{2}(\Omega), s>0\right)$, such that we may omit it and replace $\mathrm{G}(1, \tilde{r})$ by $\check{\mathrm{U}}(1)$, yielding (with (B.10))

$$
\begin{aligned}
& \left\|r^{-s} w_{e}\right\|_{\mathrm{L}^{2}(\check{\mathrm{U}}(1))}^{2} \\
& \leq c \cdot\left\{\left(\left\|r^{s} g_{e}\right\|_{\mathrm{L}^{2}(\check{\mathrm{U}}(1))}+\left\|r^{s-1} \nabla w_{e}\right\|_{\mathrm{L}^{2}(\breve{\mathrm{U}}(1))}\right) \cdot\left\|r^{1-3 s} w_{e}\right\|_{\mathrm{L}^{2}(\check{\mathrm{U}}(1))}+\int_{\mathrm{S}(1)}\left(\left|w_{e}\right|^{2}+\left|\nabla w_{e}\right|^{2}\right)\right\} \\
& \quad \leq c \cdot\left\{\left(\left\|g_{e}\right\|_{\mathrm{L}_{\mathrm{s}}^{2}\left(\mathbb{R}^{n}\right)}+\left\|\nabla w_{e}\right\|_{\mathrm{L}_{\mathrm{s}-1}^{2}\left(\mathbb{R}^{n}\right)}\right) \cdot\left\|w_{e}\right\|_{L_{1-3 s}^{2}\left(\mathbb{R}^{n}\right)}+\int_{\mathrm{S}(1)}\left(\left|w_{e}\right|^{2}+\left|\nabla w_{e}\right|^{2}\right)\right\} .
\end{aligned}
$$

Moreover, the surface integral is bounded by $\left\|w_{e}\right\|_{\mathbf{H}^{2}(\mathrm{U}(1))}^{2}$ ( trace theorem) and with (B.10) we obtain

$$
\begin{aligned}
\left\|w_{e}\right\|_{L_{-s}^{2}\left(\mathbb{R}^{n}\right)}^{2} & \leq c \cdot\left\{\left(\left\|g_{e}\right\|_{\left.L_{s}^{2}\left(\mathbb{R}^{n}\right)\right)}+\left\|\nabla w_{e}\right\|_{L_{s-1}^{2}\left(\mathbb{R}^{n}\right)}\right) \cdot\left\|w_{e}\right\|_{L_{1-3 s}^{2}\left(\mathbb{R}^{n}\right)}+\left\|w_{e}\right\|_{\mathbf{H}^{2}(\mathrm{U}(2))}^{2}\right\} \\
& \leq c \cdot\left\{\left(\left\|g_{e}\right\|_{L_{s}^{2}\left(\mathbb{R}^{n}\right)}+\left\|w_{e}\right\|_{L_{-s}^{2}\left(\mathbb{R}^{n}\right)}\right) \cdot\left\|w_{e}\right\|_{L_{1-3 s}^{2}\left(\mathbb{R}^{n}\right)}+\left\|w_{e}\right\|_{\mathbf{H}^{2}(\mathrm{U}(2))}^{2}\right\}
\end{aligned}
$$

hence ( Young's inequality)

$$
\left\|w_{e}\right\|_{\mathrm{L}_{-\mathrm{s}}^{2}\left(\mathbb{R}^{n}\right)}^{2} \leq c \cdot\left\{\left\|g_{e}\right\|_{\mathrm{L}_{\mathrm{s}}^{2}\left(\mathbb{R}^{n}\right)}+\left\|w_{e}\right\|_{\mathrm{L}_{1-3 \mathrm{~s}}^{2}\left(\mathbb{R}^{n}\right)}+\left\|w_{e}\right\|_{\mathbf{H}^{2}(\mathrm{U}(2))}^{2}\right\}
$$

Now, using once again Lemma B. 1 we arrive at

$$
\left\|w_{e}\right\|_{L_{-s}^{2}\left(\mathbb{R}^{n}\right)}^{2} \leq c \cdot\left\{\left\|g_{e}\right\|_{L_{s}^{2}\left(\mathbb{R}^{n}\right)}+\left\|w_{e}\right\|_{L_{1-3 s}^{2}\left(\mathbb{R}^{n}\right)}\right\}
$$

which together with (B.10) and Lemma A. 1 implies

$$
\begin{equation*}
\|w\|_{\mathrm{L}_{-\mathrm{s}}^{2}\left(\mathbb{R}^{n}\right)}+\|\exp (-i \nu r) w\|_{\mathrm{H}_{\mathrm{s}-2}^{1}\left(\mathbb{R}^{n}\right)} \leq c \cdot\left\{\|g\|_{\mathrm{L}_{\mathrm{s}}^{2}\left(\mathbb{R}^{n}\right)}+\|w\|_{\mathrm{L}^{2}(\Omega(\delta))}\right\} \tag{B.11}
\end{equation*}
$$

with $c, \delta>0$ independent of $\nu, \tau, w$ and $g$.

## Appendix C. Proofs in the Case of the Time-Harmonic Maxwell Equations

This section deals with the proofs of the decomposition lemma, the polynomial decay and the a-priori estimate we omitted in the main text of this paper.

Proof of Lemma 4.1. We start with $u=\eta u+\check{\eta} u$, noting that $\check{\eta} u \in \mathbf{R}_{\mathrm{t}}$. Moreover

$$
\operatorname{Rot} \check{\eta} u=\mathrm{C}_{\operatorname{Rot}, \check{\eta}} u+\check{\eta} \operatorname{Rot} u=\mathrm{C}_{\operatorname{Rot}, \check{\eta}} u-i \check{\eta} \Lambda f-i \omega \check{\eta} \Lambda u
$$

and we have

$$
\left(\operatorname{Rot}+i \omega \Lambda_{0}\right) \check{\eta} u=\left(\mathrm{C}_{\mathrm{Rot}, \check{\eta}}-i \omega \check{\eta} \hat{\Lambda}\right) u-i \check{\eta} \Lambda f=f_{1} \in \mathrm{~L}_{\mathrm{s}}^{2}
$$

since supp $\nabla \check{\eta}$ is compact and $t+\kappa \geq s$. According to [26, Theorem 4]

$$
f_{1}=f_{\mathrm{R}}+f_{\mathrm{D}}+f_{\mathcal{S}} \in{ }_{0} \mathrm{R}_{\mathrm{s}} \dot{+}_{0} \mathrm{D}_{\mathrm{s}} \dot{+} \mathcal{S}_{s}
$$

holds and we obtain

$$
i \omega \check{\eta} \Lambda_{0} u=f_{1}-\operatorname{Rot} \check{\eta} u=f_{\mathrm{D}}-\operatorname{Rot} \check{\eta} u+f_{\mathrm{R}}+f_{\mathcal{S}} .
$$

Defining

$$
\begin{aligned}
& \circ u_{1}:=-\frac{i}{\omega} \Lambda_{0}^{-1}\left(f_{\mathrm{R}}+f_{\mathcal{S}}\right) \in \mathrm{R}_{\mathrm{s}} \\
& \circ \tilde{u}:=\check{\eta} u-u_{1}=\frac{i}{\omega} \Lambda_{0}^{-1}\left(\operatorname{Rot} \check{\eta} u-f_{\mathrm{D}}\right) \in \mathbf{R}_{\mathrm{t}} \cap_{0} \mathrm{D}_{\mathrm{t}}
\end{aligned}
$$

[8, Lemma 4.2] shows $\tilde{u} \in H_{t}^{1}$ and we have

$$
\left(\operatorname{Rot}+i \omega \Lambda_{0}\right) \tilde{u}=\operatorname{Rot}\left(\check{\eta} u-u_{1}\right)+i \omega \Lambda_{0} \tilde{u}=f_{\mathrm{D}}+\frac{i}{\omega} \widetilde{\Lambda}_{0}^{-1} \operatorname{Rot} f_{\mathcal{S}}=f_{2} \in{ }_{0} \mathrm{D}_{\mathrm{s}} .
$$

Next we solve $(\operatorname{Rot}+1) u_{2}=f_{2}$. Using Fourier-transformation we look at

$$
\hat{u}_{2}:=\left(1+r^{2}\right)^{-1}(1-i r \Xi) \mathcal{F}\left(f_{2}\right)
$$

Since $s>1 / 2$ and $f_{2} \in \mathrm{~L}_{\mathrm{s}}^{2}$, we obtain $\hat{u} \in \mathrm{~L}_{1}^{2}$, hence $u_{2}:=\mathcal{F}^{-1}\left(\hat{u}_{2}\right) \in \mathbf{H}^{1}$. Moreover, $\mathcal{F}\left(\mathcal{F}\left(f_{2}\right)\right)=$ $\mathcal{P}\left(f_{2}\right) \in \mathrm{L}_{\mathrm{s}}^{2}$ ( $\mathcal{P}$ : parity operator ) yielding $\mathcal{F}\left(f_{2}\right) \in \mathbf{H}^{\mathrm{s}}$ and as product of an $\mathbf{H}^{\mathrm{s}}$-field with bounded $\mathrm{C}^{\infty}$ functions, $\hat{u} \in \mathbf{H}^{\mathrm{s}}$ (cf. [30, Lemma 3.2]), hence $u_{2} \in \mathrm{~L}_{\mathrm{s}}^{2}$. In addition a straight forward calculation shows $\mathcal{F}\left((\operatorname{Rot}+1) u_{2}\right)=\mathcal{F}\left(f_{2}\right)$, which by [8, Lemma 4.2] implies

$$
(\operatorname{Rot}+1) u_{2}=f_{2} \quad \text { and } \quad u_{2} \in \mathbf{H}_{\mathrm{s}}^{1} \cap_{0} \mathrm{D}_{\mathrm{s}} .
$$

Then $(t \leq s)$

$$
u_{3}:=\tilde{u}-u_{2} \in \mathbf{H}_{\mathrm{t}}^{1} \cap_{0} \mathrm{D}_{\mathrm{t}}
$$

satisfies

$$
\begin{aligned}
& \left(\operatorname{Rot}+i \omega \Lambda_{0}\right) u_{3} \\
& \qquad=\left(\operatorname{Rot}+i \omega \Lambda_{0}\right) \tilde{u}-\left(\operatorname{Rot}+i \omega \Lambda_{0}\right) u_{2} \\
& \quad=f_{2}-(\operatorname{Rot}+1) u_{2}+\left(1-i \omega \Lambda_{0}\right) u_{2}=\left(1-i \omega \Lambda_{0}\right) u_{2} \in \mathbf{H}_{\mathrm{s}}^{1} \cap_{0} \mathrm{D}_{\mathrm{s}}
\end{aligned}
$$

and using once more [8, Lemma 4.2] we get

$$
u_{3} \in \mathbf{H}_{\mathrm{t}}^{2} \cap_{0} \mathrm{D}_{\mathrm{t}} .
$$

Finally

$$
\begin{aligned}
\Delta u_{3} & =\operatorname{Rot}\left(\operatorname{Rot} u_{3}\right) \\
& =\left(1-i \omega \widetilde{\Lambda}_{0}\right) \operatorname{Rot} u_{2}-i \omega \widetilde{\Lambda}_{0} \operatorname{Rot} u_{3} \\
& =\left(1-i \omega \widetilde{\Lambda}_{0}\right)\left(f_{2}-u_{2}\right)-i \omega \widetilde{\Lambda}_{0}\left(\left(1-i \omega \Lambda_{0}\right) u_{2}-i \omega \Lambda_{0} u_{3}\right) \\
& =\left(1-i \omega \widetilde{\Lambda}_{0}\right) f_{2}-\left(1+\omega^{2} \varepsilon_{0} \mu_{0}\right) u_{2}-\omega^{2} \varepsilon_{0} \mu_{0} u_{3}
\end{aligned}
$$

holds, hence

$$
\left(\Delta+\omega^{2} \varepsilon_{0} \mu_{0}\right) u_{3}=\left(1-i \omega \widetilde{\Lambda}_{0}\right) f_{2}-\left(1+\omega^{2} \varepsilon_{0} \mu_{0}\right) u_{2}
$$

The asserted estimates follow by straight forward calculations using [8, Lemma 4.2] and the continuity of the projections from $\mathrm{L}_{\mathrm{s}}^{2}$ into ${ }_{0} \mathrm{R}_{\mathrm{s}},{ }_{0} \mathrm{D}_{\mathrm{s}}$ and $\mathcal{S}_{s}$.

Proof of Lemma 4.2. As for $t \geq s-1$ there is nothing to proof, we concentrate on the case

$$
u \in \mathbf{R}_{\mathrm{t}}(\Omega) \quad \text { with } \quad-1 / 2<t<s-1
$$

Therefore assume first that in addition

$$
s-\kappa<t \Longrightarrow t<s<t+\kappa
$$

Then we may apply Lemma 4.1 and decompose the field $u$ in

$$
u=\eta u+u_{1}+u_{2}+u_{3}
$$

with $\eta u+u_{1}+u_{2} \in \mathbf{R}_{\mathrm{s}}(\Omega)$ and $u_{3} \in \mathbf{H}_{\mathrm{t}}^{2}$ satisfying $\left(\Delta+\omega^{2} \varepsilon_{0} \mu_{0}\right) u_{3} \in \mathrm{~L}_{\mathrm{s}}^{2}$. Thus the polynomial decay for Helmholtz equation (cf. Appendix, Section B, Lemma B.2) shows

$$
u_{3} \in \mathbf{H}_{\mathrm{s}-1}^{2} \quad \text { and } \quad\left\|u_{3}\right\|_{\mathbf{H}_{\mathrm{s}-1}^{2}} \leq c \cdot\left\{\left\|\left(\Delta+\omega^{2} \varepsilon_{0} \mu_{0}\right) u_{3}\right\|_{\mathrm{L}_{\mathrm{s}}^{2}}+\left\|u_{3}\right\|_{\mathrm{L}_{\mathrm{s}-2}^{2}}\right\}
$$

$c=c(s, J)>0$, yielding $u=\eta u+u_{1}+u_{2}+u_{3} \in \mathbf{R}_{\mathrm{s}-1}(\Omega)$. Moreover, using the estimates of Lemma 4.1 we obtain uniformly with respect to $\omega, u$ and $f$

$$
\begin{aligned}
& \|u\|_{\mathbf{R}_{s-1}(\Omega)} \leq c \cdot\left\{\|f\|_{L_{s}^{2}(\Omega)}+\|u\|_{L_{s-\kappa}^{2}(\Omega)}+\left\|u_{3}\right\|_{L_{s-1}^{2}}\right\} \\
& \leq c \cdot\left\{\|f\|_{L_{s}^{2}(\Omega)}+\|u\|_{\mathrm{L}_{s-\kappa}^{2}(\Omega)}+\left\|\left(\Delta+\omega^{2} \varepsilon_{0} \mu_{0}\right) u_{3}\right\|_{L_{\mathrm{s}}^{2}}+\left\|u_{3}\right\|_{L_{\mathrm{s}-2}^{2}}\right\} \\
& \leq c \cdot\left\{\|f\|_{L_{s}^{2}(\Omega)}+\|u\|_{L_{s-m}^{2}(\Omega)}\right\},
\end{aligned}
$$

where $m:=\min \{\kappa, 2\}$ and applying Lemma A. 1 we end up with:

$$
\|u\|_{\mathbf{R}_{\mathrm{s}-1}(\Omega)} \leq c \cdot\left\{\|f\|_{\mathrm{L}_{\mathrm{s}}^{2}(\Omega)}+\|u\|_{\mathrm{L}^{2}(\Omega(\delta))}\right\}
$$

for $c, \delta \in(0, \infty)$ independent of $\omega, u$ and $f$. So let's switch to the case

$$
t \leq s-\kappa \Longrightarrow t+\kappa \leq s
$$

Here, the idea is to approach $s$ by overlapping intervals to which the first case is applicable. For that, we choose some $\hat{k} \in \mathbb{N}$, such that with $\gamma:=(\kappa-1) / 2>0$ we have

$$
t+\kappa+(\hat{k}-1) \cdot \gamma \leq s \leq t+\kappa+\hat{k} \cdot \gamma
$$

and for $k=0,1, \ldots, \hat{k}$ we define

$$
t_{k}:=t+k \cdot \gamma \quad \text { as well as } \quad s_{k}:=t_{k+1}+1=t_{k}+(\kappa+1) / 2
$$

Then ( as $\kappa>1$ )

$$
t_{k+1}<s_{k}=t_{k+1}+1=t+\kappa+(k-1) \cdot \gamma \leq s \quad \text { and } \quad t_{k}<t_{k+1}+1=s_{k}=t_{k}+(\kappa+1) / 2<t_{k}+\kappa
$$

such that we can successively apply the first case, ending up with $u \in \mathbf{R}_{\mathrm{s}_{\hat{k}}-1}(\Omega)$. If $s=s_{\hat{k}}$ we are done. Otherwise we choose $t_{\hat{k}+1}:=s_{\hat{k}}-1$ and apply the first case once more, since

$$
t_{\hat{k}+1}<s_{\hat{k}}<s \leq t+\kappa+\hat{k} \cdot \gamma=t_{\hat{k}+1}+\kappa
$$

Either way we obtain $u \in \mathbf{R}_{\mathrm{s}-1}(\Omega)$ and now the estimate follows as in the first case.

Proof of Lemma 4.4. Without loss of generality we assume $s \in(1 / 2,1)$. Then $s \in \mathbb{R} \backslash \mathbb{I}$ with $0<s<\kappa$ and we can apply Lemma 4.1 ( with $t=0$ ) to decompose $u:=\mathcal{L}_{\omega} f \in \mathbf{R}_{\Gamma_{1}}(\Omega)$ into

$$
u=\eta u+u_{1}+u_{2}+u_{3}
$$

with $u_{3} \in \mathbf{H}^{2}$ solving

$$
\left(\Delta+\omega^{2} \varepsilon_{0} \mu_{0}\right) u_{3}=\left(1-i \omega \widetilde{\Lambda}_{0}\right) f_{2}-\left(1+\omega^{2} \varepsilon_{0} \mu_{0}\right) u_{2}=: f_{3} \in \mathrm{~L}_{\mathrm{s}}^{2}
$$

where $f_{2}$ is defined as in Lemma 4.1. Moreover, the estimates from Lemma 4.1 along with

$$
\left(\operatorname{Rot}-i \omega \sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u=-i \Lambda f-i \omega\left(\Lambda_{0}+\sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u-i \omega \hat{\Lambda} u
$$

yield

$$
\begin{align*}
& \|u\|_{\mathbf{R}_{t}(\Omega)}+\left\|\left(\Lambda_{0}+\sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u\right\|_{L_{s}^{2}-1}(\Omega) \\
& \quad \leq c \cdot\left\{\|u\|_{\mathbf{R}_{\mathrm{t}}(\Omega)}+\left\|\left(\operatorname{Rot}-i \omega \sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u\right\|_{\mathrm{L}_{\mathrm{s}-1}^{2}(\Omega)}+\|f\|_{\mathrm{L}_{\mathrm{s}}^{2}(\Omega)}+\|u\|_{\mathrm{L}_{\mathrm{s}-\kappa}^{2}(\Omega)}\right\}  \tag{C.1}\\
& \quad \leq c \cdot\left\{\left\|u_{3}\right\|_{\mathrm{L}_{\mathrm{t}}^{2}}+\left\|\left(\operatorname{Rot}-i \omega \sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u_{3}\right\|_{\mathrm{L}_{\mathrm{s}-1}^{2}}+\|f\|_{\mathrm{L}_{\mathrm{s}}^{2}(\Omega)}+\|u\|_{\mathrm{L}_{\mathrm{s}-\kappa}^{2}(\Omega)}\right\}
\end{align*}
$$

with $c=c(s, t, J)>0$. Due to the monotonicity of the norms with respect to $t$ and $s$, we may assume $t$ and $s$ to be close enough to $-1 / 2$ resp. $1 / 2$ such that $1<s-t<\kappa$ holds. Hence, the assertion follows by (C.1) and Lemma A.1, if we can show

$$
\left\|u_{3}\right\|_{L_{t}^{2}}+\left\|\left(\operatorname{Rot}-i \omega \sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u_{3}\right\|_{L_{s-1}^{2}} \leq c \cdot\left\{\|f\|_{L_{s}^{2}(\Omega)}+\|u\|_{L_{s-\kappa}^{2}(\Omega)}\right\}
$$

with $c \in(0, \infty)$ independent of $\omega, u$ and $f$. Therefore note thatthe self-adjointness of the laplacian $\Delta: \mathbf{H}^{2} \subset \mathrm{~L}^{2} \longrightarrow \mathrm{~L}^{2}$ yields $\left(\Delta+\omega^{2} \varepsilon_{0} \mu_{0}\right)^{-1} f_{3}=u_{3}$ and applying Lemma B. 3 componentwise, we obtain

$$
\left\|u_{3}\right\|_{\mathrm{L}_{\mathrm{t}}^{2}}+\left\|\exp \left(-i \lambda \sqrt{\varepsilon_{0} \mu_{0}} r\right) u_{3}\right\|_{\mathrm{H}_{\mathrm{s}-2}^{1}} \leq c \cdot\left\{\left\|f_{3}\right\|_{\mathrm{L}_{\mathrm{s}}^{2}}+\left\|u_{3}\right\|_{\mathrm{L}^{2}(\Omega(\delta))}\right\} .
$$

With $\operatorname{Rot}\left(\exp \left(-i \lambda \sqrt{\varepsilon_{0} \mu_{0}} r\right) u_{3}\right)=\exp \left(-i \lambda \sqrt{\varepsilon_{0} \mu_{0}} r\right)\left(\operatorname{Rot}-i \lambda \sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u_{3}$ this leads to

$$
\begin{align*}
& \left\|u_{3}\right\|_{L_{t}^{2}}+\left\|\left(\operatorname{Rot}-i \lambda \sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u_{3}\right\|_{L_{s-1}^{2}} \\
& \leq\left\|u_{3}\right\|_{L_{\mathrm{t}}^{2}}+\left\|\operatorname{Rot}\left(\exp \left(-i \lambda \sqrt{\varepsilon_{0} \mu_{0}} r\right) u_{3}\right)\right\|_{L_{s-1}^{2}}  \tag{C.2}\\
& \quad \leq\left\|u_{3}\right\|_{L_{t}^{2}}+\left\|\exp \left(-i \lambda \sqrt{\varepsilon_{0} \mu_{0}} r\right) u_{3}\right\|_{\mathrm{H}_{\mathrm{s}-2}^{1}} \leq c \cdot\left\{\left\|f_{3}\right\|_{\mathrm{L}_{\mathrm{s}}^{2}}+\left\|u_{3}\right\|_{\mathrm{L}^{2}(\Omega(\delta))}\right\}
\end{align*}
$$

where $c>0$ is not depending on $\omega, u_{3}$ and $f_{3}$. But we would like to estimate $\left(\operatorname{Rot}-i \omega \sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u_{3}$. For that we need some additional arguments, starting with the observation that

$$
\omega=|\lambda|\left(1+(\sigma / \lambda)^{2}\right)^{1 / 4} \cdot\left\{\begin{array}{cc}
\exp (i \varphi / 2) & \text { for } \lambda>0 \\
\exp (i(\varphi / 2+\pi)) & \text { for } \lambda<0
\end{array} \text { with } \varphi:=\arctan (\sigma / \lambda) \in\left(-\frac{\pi}{2}, \frac{\pi}{2}\right)\right.
$$

hence $|\operatorname{Re}(\omega)| \geq \sqrt{2} / 2 \cdot|\lambda|$. Thus $|\omega+\lambda| \geq \sqrt{3 / 2} \cdot|\lambda|$ and we have

$$
|\omega-\lambda|^{2}=\left|\frac{\omega^{2}-\lambda^{2}}{\omega+\lambda}\right|^{2}=\left|\frac{i \sigma \lambda}{\omega+\lambda}\right|^{2} \leq \frac{2}{3} \cdot \sigma^{2} .
$$

From this and the resolvent estimate

$$
\left\|f_{3}\right\|_{\mathrm{L}^{2}}=\left\|\left(\Delta+\omega^{2} \varepsilon_{0} \mu_{0}\right) u_{3}\right\|_{\mathrm{L}^{2}} \geq\left|\operatorname{Im}\left(\omega^{2} \varepsilon_{0} \mu_{0}\right)\right| \cdot\left\|u_{3}\right\|_{\mathrm{L}^{2}}=\varepsilon_{0} \mu_{0} \sigma|\lambda| \cdot\left\|u_{3}\right\|_{\mathrm{L}^{2}} .
$$

we obtain $(s>1 / 2)$

$$
\begin{aligned}
& \left\|\left(\operatorname{Rot}-i \omega \sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u_{3}\right\|_{L_{\mathrm{s}-1}^{2}} \\
& \quad \leq\left\|\left(\operatorname{Rot}-i \lambda \sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u_{3}\right\|_{\mathrm{L}_{\mathrm{s}-1}^{2}}+\left\|(\omega-\lambda) \sqrt{\varepsilon_{0} \mu_{0}} \Xi u_{3}\right\|_{\mathrm{L}_{\mathrm{s}-1}^{2}} \\
& \quad \leq\left\|\left(\operatorname{Rot}-i \lambda \sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u_{3}\right\|_{\mathrm{L}_{\mathrm{s}-1}^{2}}+c \cdot|\lambda|^{-1}\left\|f_{3}\right\|_{\mathrm{L}_{\mathrm{s}}^{2}},
\end{aligned}
$$

such that with (C.2) and the estimates from Lemma 4.1 uniformly with respect to $\omega, u$ and $f$ :

$$
\left\|u_{3}\right\|_{L_{\mathrm{t}}^{2}}+\left\|\left(\operatorname{Rot}-i \omega \sqrt{\varepsilon_{0} \mu_{0}} \Xi\right) u_{3}\right\|_{L_{\mathrm{s}-1}^{2}} \leq c \cdot\left\{\|f\|_{L_{\mathrm{s}}^{2}(\Omega)}+\|u\|_{\mathrm{L}_{\mathrm{s}-\kappa}^{2}(\Omega)}\right\} .
$$

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[^1]:    ${ }^{i}$ Although the right hand sides 0 and $r^{-2}$ are $\mathrm{L}^{2}(\breve{\mathrm{U}}(1))$-functions, we have $E=\xi / r \notin \mathrm{~L}^{2}(\breve{\mathrm{U}}(1))$, but $E \in \mathrm{~L}_{-1}^{2}(\check{\mathrm{U}}(1))$

[^2]:    ${ }^{\text {ii }}$ Here and in what follows, $c$ resp. $c(\ldots)$ always denotes some constant changing from step to step.

[^3]:    ${ }^{\text {iv }}$ Otherwise we replace $s$ and $t$ by $t_{k}:=t+k / 4$ resp. $s_{k}:=t_{k+1}, k=0,1,2, \ldots$ and obtain the assertion after finitely many steps of the type $t_{k}<s_{k} \leq t_{k}+1 / 2$ (cf. Appendix, Section C, Proof of Lemma 4.1, ).

