Scattering of time-harmonic electromagnetic plane waves by perfectly conducting diffraction gratings

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submitted: March 16, 2012

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No. 1694
Berlin 2012

2010 Mathematics Subject Classification. 78A45, 35J20, 35R30, 78A46.

Key words and phrases. Electromagnetic scattering, diffraction gratings, variational approach, mortar technique, non-uniqueness.

The authors would like to thank their colleague J. Elschner for pointing out the research articles [15, 21] which stimulate this paper. The first author gratefully acknowledges the support by the German Research Foundation (DFG) under Grant No. EL 584/1-2.
Abstract

Consider scattering of time-harmonic electromagnetic plane waves by a doubly periodic surface in $\mathbb{R}^3$. The medium above the surface is supposed to be homogeneous and isotropic with a constant dielectric coefficient, while below is a perfectly conducting material. This paper is concerned with the existence of quasiperiodic solutions for any frequency of incidence. Based on an equivalent variational formulation established by the mortar technique of Nitsche, we verify the existence of solutions for a broad class of incident waves including plane waves, under the assumption that the grating profile is a Lipschitz biperiodic surface. Our solvability result covers the resonance case where a Rayleigh frequency is allowed. Non-uniqueness examples are also presented in the resonance case and the TE or TM polarization case for classical gratings.

1 Introduction

Consider a time-harmonic electromagnetic plane wave incident from above to a biperiodic surface $\tilde{\Gamma}$ in $\mathbb{R}^3$. Here a biperiodic (or doubly periodic) surface means a continuous surface which is $\Lambda_1$-periodic in $x_1$, $\Lambda_2$-periodic in $x_2$ and bounded in $x_3$. This biperiodic surface divides $\mathbb{R}^3$ into two regions. The dielectric coefficient in the upper region $\tilde{\Omega}$ is supposed to be a fixed positive constant, while the medium below $\tilde{\Gamma}$ is a perfect conductor. Such structures are also called crossed diffraction gratings in the engineering and physics literature. They have many important applications in micro optics and semiconductor industry. This paper is concerned with a new existence result for the scattering problem for fixed incident frequency $\omega \in \mathbb{R}^+$. 

There are many references on the scattering of electromagnetic waves by general inhomogeneous diffraction gratings in $\mathbb{R}^3$. First rigorous results on existence and uniqueness are obtained by Chen & Friedman [8], Nédélec & Starling [23] using integral equation methods. In [1], Abboud introduces a variational formulation in a truncated periodic cell involving a nonlocal boundary (Dirichlet-to-Neumann) operator for a transparent boundary condition. This variational problem is of saddle point type and the existence and uniqueness follow from the Fredholm alternative. In the case of a constant magnetic permeability in $\mathbb{R}^3$, Abboud’s arguments have been adapted to isotropic biperiodic inhomogeneous medium by Dobson [10], Bao [5], Bao and Dobson [4], Schmidt [27] and to anisotropic optical materials by Schmidt [26], with new variational formulations that only involve the magnetic field. These variational formulations are proved to be strongly elliptic over a certain quasiperiodic Sobolev space and appear to be well adapted for the analytic and numerical treatment of quite general diffraction structures. It is proved that for all but possibly a discrete set of frequencies accumulating at infinity there always exists a unique quasiperiodic solution of locally finite energy. Moreover, uniqueness for any frequency can be guaranteed if an absorbing (lossy) material is included into the grating (see [18, 26]).
a non-absorbing (lossless) inhomogeneous material satisfies a certain non-trap condition (see [6] in the cases of TE and TM polarizations).

The existence and uniqueness results mentioned above apply to our scattering problem, since they are already implicitly contained in Abboud [1] and Ammari [2]. Some researchers seem to believe that the quasiperiodic solution in $H_{loc}(\text{curl}, \Omega)$ is unique for all frequencies, provided the perfectly grating profile is given by the graph of a $C^2$-smooth periodic function and the Rayleigh frequencies are excluded. However, in this paper we will present a counterexample (see example 2.5 in Section 2) to show that uniqueness does not hold. This counterexample is constructed in the TM polarization case, where the perfectly conducting boundary value problem of the curl-curl equation is reduced to the Neumann boundary value problem of the two-dimensional scalar Helmholtz equation. This reduction enables us to construct non-uniqueness examples to the Maxwell system, relying on the existence of non-trivial solutions for the reduced homogeneous Neumann problem established in [21]. Non-uniqueness examples in the resonance or non-graph case are also presented in this paper.

Since a grating profile is a special case of a rough surface, these non-uniqueness examples reported in Section 2 can be also viewed as counterexamples to the electromagnetic scattering by perfectly conducting rough surfaces. Concerning the variational approach applied to electromagnetic rough surface scattering problems modeled by the full Maxwell system, we refer to the recent publications by Li, Wu & Zheng [20] where existence and uniqueness is established for an incident magnetic or electric dipole in a lossy medium, and to Haddar & Lechleiter [16] in the more challenging case of a penetrable dielectric layer. As far as we know, the well-posedness of electromagnetic scattering by perfectly conducting rough surfaces or biperiodic structures in a homogeneous non-absorbing (lossless) medium is still an open problem.

Our aim of this paper is to prove the following existence result to the scattering problem: for any fixed incident wavenumber $k > 0$ there always exists a quasiperiodic solution in $H_{loc}(\text{curl}, \tilde{\Omega})$ for a broad class of incident waves including plane waves, whenever the grating profile is a Lipschitz biperiodic surface. This result is rather general, because the grating profile is not necessarily the graph of a smooth periodic function and a Rayleigh frequency is allowed. The non-graph gratings have many practical applications in diffractive optics and in optimal design of complicated grating structures. As an example, we mention the binary gratings which are composed of only a finite number of horizontal and vertical segments (see e.g. [12]). Note further that Rayleigh frequencies are always excluded in many references on electromagnetic scattering by biperiodic structures. This is mainly due either to the definition of the quasiperiodic fundamental solution to the Helmholtz equation needed in integral equation methods ([8, 23]) or to the explicit formula for the Dirichlet-to-Neumann map of the transparent boundary condition (see e.g. [2, 10, 4] or (2.8) below). Note that both, the quasiperiodic fundamental solution and the Dirichlet-to-Neumann map, are well defined only if Rayleigh frequencies are excluded.

To prove the existence of quasiperiodic solutions in $H_{loc}(\text{curl}, \tilde{\Omega})$ for any frequency, we need a replacement of the Dirichlet-to-Neumann (D-to-N) map imposed on the artificial boundary $\tilde{\Gamma}_b$ above the grating surface. Motivated by the variational formulations proposed in [19, 25] using the mortar technique combined with Nitsche’s method (see Nitsche [24] and Sternberg [28]), we employ a consistent coupling of the electric field $E$ below $\tilde{\Gamma}_b$ and the scattered field $E^+$ above $\tilde{\Gamma}_b$. This way the necessary transmission conditions are fulfilled on $\tilde{\Gamma}_b$, so that $E$
belongs to $H_{loc}(\text{curl}, \tilde{\Omega})$. Combined with the Rayleigh series expansion for $E^+$, this coupling enables us to establish an equivalent variational formulation for the pair $(E, E^+)$. We show the Fredholmness of the operator generated by the corresponding sesquilinear form, and then prove the existence of quasiperiodic solutions for any frequency by applying the Fredholm theory.

This paper provides a theoretical justification of the modified Nitsche’s method applied to electromagnetic scattering problems for periodic structures. It is expected that our argument can be extended to more general inhomogeneous diffraction gratings as considered in [19, 25]. Since the D-to-N map is not involved in our variational formulation, the approximation of the transparent boundary operator employed in [5] can be avoided. Finally, note that the presented variational approach could be a basis for the numerical analysis of an FEM method (cf. [22]).

The remaining part is organized as follows. The boundary value problem (BVP) and the needed Sobolev spaces are rigorously defined in Section 2. Our main result Theorem 2.1 on the existence of solutions and some non-uniqueness examples are also presented in this section. In Section 3, we propose a variational formulation based on the method of Nitsche and prove its equivalence to (BVP). The Fredholmness of the operator generated by the corresponding sesquilinear form will be established in Section 4. Finally we prove our main Theorem 2.1 in Section 5 by applying the Fredholm alternative.

2 Mathematical formulations and non-uniqueness examples

Consider the scattering of an electromagnetic plane wave by a perfectly conducting grating profile in an isotropic homogeneous lossless medium. Recall that the symbol $\tilde{\Gamma}$ denotes the grating profile which is $(\Lambda_1, \Lambda_2)$-periodic in $(x_1, x_2)$ and that $\tilde{\Omega}$ denotes the region above $\tilde{\Gamma}$. Suppose that a time-harmonic incident electromagnetic plane wave $E^{in}$ (time dependence $e^{-i\omega t}$) given by

\[ E^{in} := q \exp(ikx \cdot \hat{\theta}) = q \exp \left( i(x' \cdot \alpha - \beta x_3) \right), \quad i := \sqrt{-1} \tag{2.1} \]

is incident to the grating from above. Here $k := \omega \sqrt{\epsilon \mu}$ is the positive wavenumber in terms of the angular frequency $\omega$, the electric permittivity $\epsilon$ and the magnetic permeability $\mu$, which are assumed to be positive constants everywhere in $\tilde{\Omega}$. The symbol $\hat{\theta}$ denotes the direction of incidence

\[
\hat{\theta} := (\sin \theta_1 \cos \theta_2, \sin \theta_1 \sin \theta_2, -\cos \theta_1) \in S^2,
\]

\[
S^2 := \left\{ x = (x_1, x_2, x_3) \in \mathbb{R}^3 : ||x|| = 1 \right\}
\]

with the incident angles $\theta_1 \in [0, \pi/2), \theta_2 \in [0, 2\pi)$. Throughout the paper, the symbol $(\cdot)^\top$ denotes the transpose of a row vector in $\mathbb{C}^2$ or $\mathbb{C}^3$. In (2.1), the vector $q = (q_1, q_2, q_3)^\top \in S^2$ stands for the direction of polarization satisfying $q \perp \hat{\theta}$, and

\[
x' := (x_1, x_2) \in \mathbb{R}^2, \quad \alpha = (\alpha_1, \alpha_2)^\top := k(\sin \theta_1 \cos \theta_2, \sin \theta_1 \sin \theta_2) \in \mathbb{R}^2, \quad \beta := k \cos \theta_1.
\]

Since the substrate below $\tilde{\Gamma}$ is a perfect conductor, the total electric field $E$, which can be decomposed as the sum of the incident field $E^{in}$ and the scattered field $E^{sc}$, satisfies the
where $\nu = (\nu_1, \nu_2, \nu_3)^\top \in \mathbb{S}^2$ is the unit normal on $\tilde{\Gamma}$ pointing into the exterior of $\tilde{\Omega}$. The total electric field $E$ fulfills the reduced time-harmonic curl-curl equation
\begin{equation}
\text{curl} \, \text{curl} \, E - k^2 E = 0 \quad \text{in} \quad \tilde{\Omega}. \tag{2.3}
\end{equation}
Since the grating profile is biperiodic, we require $E$ to be $\alpha$-quasiperiodic in the sense that
\begin{align*}
E(x_1 + \Lambda_1, x_2, x_3) &= \exp(i\Lambda_1 \alpha_1)E(x_1, x_2, x_3), \quad (x_1, x_2, x_3)^\top \in \tilde{\Omega}, \tag{2.4}
E(x_1, x_2 + \Lambda_2, x_3) &= \exp(i\Lambda_2 \alpha_2)E(x_1, x_2, x_3), \quad (x_1, x_2, x_3)^\top \in \tilde{\Omega}.
\end{align*}
We also impose a radiation condition in the $x_3$-direction by assuming that the scattered field $E^{sc}$ is composed of bounded outgoing plane waves:
\begin{equation}
E^{sc}(x) = \sum_{n \in \mathbb{Z}^2} E_n \exp \left( i(\alpha_n \cdot x' + \beta_n x_3) \right) \quad \text{for} \quad x_3 > \Gamma_{\max} := \max_{x \in \Gamma} \{x_3\}, \quad E_n \perp (\alpha_n, \beta_n)^\top, \tag{2.5}
\end{equation}
where $\alpha_n := (\alpha_n^{(1)}, \alpha_n^{(2)})^\top \in \mathbb{R}^2$, with $\alpha_n^{(j)} = \alpha_j + 2\pi n_j / \Lambda_j, \quad j = 1, 2$ for $n = (n_1, n_2)^\top \in \mathbb{Z}^2$, and
$$
\beta_n = \beta_n(k, \alpha) := \begin{cases} \sqrt{k^2 - |\alpha_n|^2} & \text{if} \ |\alpha_n| \leq k, \\
 i \sqrt{|\alpha_n|^2 - k^2} & \text{if} \ |\alpha_n| > k.
\end{cases}
$$
For the constant coefficient vector $E_n = (E_n^{(1)}, E_n^{(2)}, E_n^{(3)})^\top \in \mathbb{C}^3$, the relation $E_n \perp (\alpha_n, \beta_n)^\top$ means that $E_n^{(1)} \alpha_n^{(1)} + E_n^{(2)} \alpha_n^{(2)} + E_n^{(3)} \beta_n = 0$. The series in (2.5), which is also referred to as the Rayleigh series expansion, is the radiation condition we are going to use in the following sections. The constant vectors $E_n$ are called the Rayleigh coefficients. Since $\beta_n$ are real-valued only for the finitely many indices $n$ from the set $\{ n \in \mathbb{Z}^2 : |\alpha_n| \leq k^2 \}$, we observe that only a finite number of plane waves in (2.5) propagate into the far field, while the remaining part consists of evanescent (or surface) waves decaying exponentially as $x_3 \to +\infty$. Thus, the above expansion converges uniformly with all derivatives in the half plane $\{ x_3 > a \}$ for any $a > \Gamma_{\max}$.

It is assumed throughout this paper that the grating profile $\tilde{\Gamma}$ is a Lipschitz biperiodic surface in $\mathbb{R}^3$, which is not necessarily the graph of a biperiodic function. Since the unbounded domain $\tilde{\Omega}$ is $(\Lambda_1, \Lambda_2)$-periodic in $x'$ and the incident and scattered fields are both quasiperiodic, we can reduce the scattering problem to a single periodic cell $\Omega$. To this end, we introduce the following notation:
\begin{align*}
\Gamma &:= \left\{ (x_1, x_2, x_3)^\top \in \tilde{\Gamma} : 0 < x_j < \Lambda_j, j = 1, 2 \right\}, \\
\Omega &:= \left\{ (x_1, x_2, x_3)^\top \in \tilde{\Omega} : 0 < x_j < \Lambda_j, j = 1, 2 \right\}, \\
\tilde{\Gamma}_b &:= \left\{ (x_1, x_2, x_3)^\top : x_3 = b \right\}, \\
\Gamma_b &:= \left\{ (x_1, x_2, x_3)^\top \in \tilde{\Gamma}_b : 0 < x_j < \Lambda_j, j = 1, 2 \right\}, \\
\Omega_b &:= \left\{ x \in \Omega : x_3 < b \right\}
\end{align*}
for some \( b \) with \( b > \Gamma_{\text{max}} \). We next introduce some scalar and vector valued \( \alpha \)-quasiperiodic Sobolev spaces. Let \( H^s(\tilde{\Gamma}_b) \) be the complex valued \( L^2 \)-based Sobolev spaces of order \( s \) in \( \tilde{\Gamma}_b \). Write

\[
H_{\text{loc}}(\text{curl }, \Omega) := \left\{ G : \chi G, \text{curl } (\chi G) \in L^2(\Omega), \forall \chi \in C^\infty_0(\mathbb{R}^3) \right\},
\]

\[
H^s_{\text{loc}}(\tilde{\Gamma}_b) := \left\{ G : \chi G \in H^s(\tilde{\Gamma}_b), \forall \chi \in C^\infty_0(\tilde{\Gamma}_b) \right\},
\]

\[
H^s_{t,\text{loc}}(\tilde{\Gamma}_b) := \left\{ G \in H^s_{\text{loc}}(\tilde{\Gamma}_b)^3 : e_3 \cdot G = 0 \right\}, \quad e_3 := (0, 0, 1)^T,
\]

\[
H^s_{t,\text{loc}}(\text{Div }, \tilde{\Gamma}_b) := \left\{ G : G \in H^s_{t,\text{loc}}(\tilde{\Gamma}_b), \text{Div } G \in H^s_{\text{loc}}(\tilde{\Gamma}_b) \right\},
\]

\[
H^s_{t,\text{loc}}(\text{Curl }, \tilde{\Gamma}_b) := \left\{ G : G \in H^s_{t,\text{loc}}(\tilde{\Gamma}_b), \text{Curl } G \in H^s_{\text{loc}}(\tilde{\Gamma}_b) \right\},
\]

and

\[
H(\text{curl }, \Omega_b) := \left\{ G|_{\Omega_b} : G \in H_{\text{loc}}(\text{curl }, \tilde{\Omega}), G \text{ is } \alpha \text{-quasiperiodic} \right\},
\]

\[
H^s(\Gamma_b) := \left\{ G|_{\Gamma_b} : G \in H^s_{t,\text{loc}}(\tilde{\Gamma}_b)^3, G \text{ is } \alpha \text{-quasiperiodic} \right\},
\]

\[
H^s(\text{Div }, \Gamma_b) := \left\{ G|_{\Gamma_b} : G \in H^s_{t,\text{loc}}(\text{Div }, \tilde{\Gamma}_b)^3, G \text{ is } \alpha \text{-quasiperiodic} \right\},
\]

\[
H^s(\text{Curl }, \Gamma_b) := \left\{ G|_{\Gamma_b} : G \in H^s_{t,\text{loc}}(\text{Curl }, \tilde{\Gamma}_b)^3, G \text{ is } \alpha \text{-quasiperiodic} \right\},
\]

where \( \text{Div } (\cdot) \) and \( \text{Curl } (\cdot) \) stand for the surface divergence and the surface scalar rotational operators, respectively. Note that, for \( E(x') \in H^s_t(\Gamma_b), s \in \mathbb{R}, \) we have the Fourier series expansion

\[
E(x') = \sum_{n \in \mathbb{Z}^2} E_n \exp(i\alpha_n \cdot x'),
\]

with

\[
E_n := (A_1 A_2)^{-1} \int_0^{A_1} \int_0^{A_2} E(x') \exp(-i\alpha_n \cdot x') dx_1 dx_2.
\]

Then, the spaces \( H^s_t(\Gamma_b), H^s_t(\text{Div }, \Gamma_b) \) and \( H^s_t(\text{Curl }, \Gamma_b) \) can be equipped with the following equivalent Sobolev norms

\[
||E||_{H^s_t(\Gamma_b)} = \left( \sum_{n \in \mathbb{Z}^2} (1 + |\alpha_n|^2)^s \right)^{1/2},
\]

\[
||E||_{H^s_t(\text{Div }, \Gamma_b)} = \left( \sum_{n \in \mathbb{Z}^2} (1 + |\alpha_n|^2)^s \left( |E_n|^2 + |E_n \cdot \alpha_n|^2 \right) \right)^{1/2},
\]

\[
||E||_{H^s_t(\text{Curl }, \Gamma_b)} = \left( \sum_{n \in \mathbb{Z}^2} (1 + |\alpha_n|^2)^s \left( |E_n|^2 + |E_n \times \alpha_n|^2 \right) \right)^{1/2}.
\]
Recall that the space dual to $H_t^s(\text{curl}, \Gamma_b)$ w.r.t. the $L^2$-scalar product is

$$H_t^s(\text{Div}, \Gamma_b)' = H_t^{-s-1}(\text{Curl}, \Gamma_b),$$

and that, for $s = -1/2$,

$$H_t^{-1/2}(\text{Div}, \Gamma_b) = \left\{ e_3 \times E|_{\Gamma_b} : E \in H(\text{curl}, \Omega_b) \right\},$$

$$H_t^{-1/2}(\text{Curl}, \Gamma_b) = \left\{ (e_3 \times E|_{\Gamma_b}) \times e_3 : E \in H(\text{curl}, \Omega_b) \right\}.$$

Further, the trace mappings from $H(\text{curl}, \Omega_b)$ to the tangential spaces $H_t^{-1/2}(\text{Div}, \Gamma_b)$ and $H_t^{-1/2}(\text{Curl}, \Gamma_b)$ are continuous and surjective (see [7, 22] and the references there). Finally, define our variational space

$$X = X_b := \left\{ E : \Omega_b \to \mathbb{C}^3 : E \in H(\text{curl}, \Omega_b), \nu \times E|_{\Gamma} = 0 \right\}$$

endowed with the norm

$$||E||_X := ||E||_{H(\text{curl}, \Omega_b)} = \left( ||E||_{L^2(\Omega_b)}^2 + ||\text{curl } E||_{L^2(\Omega_b)}^2 \right)^{1/2}.$$

The boundary value problem for our scattering problem can be stated as follows. Let the grating profile $\Gamma$ and the number $b > \Gamma_{\text{max}}$ be fixed.

**BVP:** Given an incident electric field $E^{\text{in}}$, determine the total field $E = E^{\text{in}} + E^{\text{sc}} \in X$ such that $E$ satisfies the curl-curl equation (2.3) over $\Omega_b$ in a distributional sense and that $E^{\text{sc}}$ admits a Rayleigh expansion (2.5) valid for any $\Gamma_{\text{max}} < x \leq b$.

Note that any $E^{\text{sc}}$ satisfying (2.5) in the strip $\Gamma_{\text{max}} < x \leq b$ can be extended to the upper half space by (2.5) (see Remark 3.4 in Section 3). Below is our main result to (BVP) for a broad class of incident waves.

**Theorem 2.1.** Assume that the incident electric wave takes the form

$$E^{\text{in}}_{\text{gen}} := \sum_{n: \beta_n > 0} Q_n \exp \left( \alpha_n \cdot x' - \beta_n x_3 \right),$$

where $Q_n \in \mathbb{C}^3$ satisfies $Q_n \perp (\alpha_n, -\beta_n)^\top$. Then the problem (BVP) admits at least one solution for any $k \in \mathbb{R}^+$. Moreover, the part of the solution reflected into the upper half space is unique, i.e., the Rayleigh coefficients of the plane wave modes propagating into the upper half space (namely, those with $\beta_n > 0$) are unique.

Note that $E^{\text{in}}$ of (2.1) is of the form (2.6), where $Q_n = q$ for $n = (0, 0)^\top$ and $Q_n = (0, 0, 0)^\top$ else. We do not exclude “resonances” in Theorem 2.1, that is, the set defined by

$$\Upsilon := \left\{ n \in \mathbb{Z}^2 : \beta_n(k, \alpha) = 0 \right\}$$

is allowed to be nonempty. An incident angular frequency $\omega$ with $\Upsilon \neq \emptyset$ is called Rayleigh frequency. Note that the set of Rayleigh frequencies depends on $\Lambda_1$ and $\Lambda_2$ but not on the shape of $\Gamma$.
Remark 2.2. It seems to be known that, for all wavenumbers except those from a sequence \( k_j \in \mathbb{R}^+ \), \( k_j \to +\infty \), the problem (BVP) admits a unique solution. To see this, one may consider the variational formulation

\[
\int_{\Omega_b} [\text{curl} \ E \cdot \text{curl} \ \varphi - k^2 E \cdot \varphi] \, dx - \int_{\Gamma_b} \mathcal{R}(e_3 \times E) \cdot (e_3 \times \varphi) \, ds
= \int_{\Gamma_b} [(\text{curl} \ E^{in})_T - \mathcal{R}(e_3 \times E^{in})] \cdot (e_3 \times \varphi) \, ds
\tag{2.7}
\]

for all \( \varphi \in X \), where \( (\cdot)_T := \nu \times (\cdot)|_{\Gamma_b} \times \nu \), and \( \mathcal{R} : H^{-1/2}_t(Div, \Gamma_b) \to H^{-1/2}_t(Curl, \Gamma_b) \) is the Dirichlet-to-Neumann map defined by

\[
(\mathcal{R} \tilde{E})(x') = -\sum_{n \in \mathbb{Z}^2} \frac{1}{i \beta_n} \left[ k^2 \tilde{E}_n - (\alpha_n \cdot \tilde{E}_n) \alpha_n \right] \exp(i \alpha_n \cdot x'),
\tag{2.8}
\]

for \( \tilde{E}(x') = \sum_{n \in \mathbb{Z}^2} \tilde{E}_n \exp(i \alpha_n \cdot x') \in H^{-1/2}_t(Div, \Gamma_b) \); see [1, 2]. Note that the operator \( \mathcal{R} \) maps \( e_3 \times E^{sc} \) to \( (\text{curl} \ E^{sc})_T \) on \( \Gamma_b \) and that Rayleigh frequencies must be excluded in (2.8). An alternative Dirichlet-to-Neumann operator for the magnetic field is given in [4, 5, 10].

It is seen from Lemma 6.1 in Section 6 that the variational formulation is uniquely solvable for all frequencies \( k \in (0, k_0) \) with \( k_0 > 0 \) being sufficiently small. This combined with the analytic Fredholm theory (see e.g. [9, Theorem 8.26] or [14, Theorem I.5.1]) leads to the existence and uniqueness for all \( k \in \mathbb{R}^+ \setminus D \), where \( D \supseteq \Upsilon \) is a discrete set with the only accumulating point at infinity. Since such a solvability result is contained in many references on diffraction gratings, we skip the details and refer to [26, 11, 12] for the applications of the analytic Fredholm theory in periodic structures. However, it follows from the examples below that the uniqueness to (BVP) does not hold in general, even if \( \Gamma \) is a smooth graph and Rayleigh frequencies are excluded.

The proof of Theorem 2.1 will be carried out in Section 5 using an equivalent formulation which covers the resonance case. Next we present some non-uniqueness examples to (BVP) by constructing non-trivial solutions to the homogeneous scattering problem \( (E^{in} = 0) \). Now suppose that the periodicities \( \Lambda_1 \) and \( \Lambda_2 \) are fixed for the remainder of this paper.

Example 2.3. For any fixed Rayleigh frequency \( \omega \), there exists a biperiodic surface \( \tilde{\Gamma} \) such that the solutions to (BVP) are non-unique.

Indeed, the grating profile defined by \( \tilde{\Gamma} := \{ x_3 = 0 \} \) is such an example. Set

\[
E^{sc}(x) = e_3 \sum_{n: \beta_n = 0} C_n \exp(i \alpha_n \cdot x'), \quad C_n \in \mathbb{C}.
\]

Then \( E^{sc} \) is \( \alpha \)-quasiperiodic and satisfies the curl-curl equation (2.3), the Rayleigh expansion condition (2.5) as well as the boundary condition (2.2).

In the following examples, the branch of the square root is chosen such that its imaginary part is non-negative, i.e., \( \sqrt{\alpha} = i \sqrt{-\alpha} \) if \( \alpha < 0 \).
Example 2.4. There exists a non-Rayleigh frequency \( \omega \) and a non-graph grating profile \( \tilde{\Gamma} \) such that the solutions to (BVP) are non-unique.

Restrict the search for examples to gratings which remain invariant in \( x_2 \) direction. We seek a special solution of the form \( E^{sc}(x) = (0, u^{sc}(x_1, x_3), 0)^\top \), where the scalar function \( u^{sc} \) fulfills

\[
(\Delta + k^2) u^{sc} = 0 \quad \text{in} \quad \tilde{\Omega}_0 := \tilde{\Omega} \cap \{ x_2 = 0 \}, \quad u^{sc} = 0 \quad \text{on} \quad \tilde{\Gamma} \cap \{ x_2 = 0 \},
\]

\[
u \times u^{sc} = 0 \quad \text{on} \quad \tilde{\Gamma} \cap \{ x_2 = 0 \},
\]

\[
u \cdot u^{sc} = C_{n_1} \exp \left( \left( i [\alpha_n^{(1)} x_1 + (k^2 - |\alpha_n^{(1)}|^2) / 2 x_3] \right) \right), \quad (x_1, x_3)^\top \in \tilde{\Omega}_0, \quad x_3 > \Gamma_{\text{max}},
\]

with \( C_{n_1} \in \mathbb{C} \). Recall that \( n = (n_1, n_2)^\top \in \mathbb{Z}^2 \) and \( \alpha_n^{(j)} \) denotes the \( j \)-th component of \( \alpha_n \in \mathbb{R}^2 \), \( j = 1, 2 \). In fact, the previous Dirichlet boundary value problem is the TE polarization of (BVP). Non-trivial solutions to the above problem do exist for the \( \Lambda_1 \)-periodic non-graph grating profile constructed in [15] with \( \Lambda_1 = 2\pi \). Thus, the solution \( E^{sc} \), which is independent of \( x_2 \) and transversal to the \((x_1, x_3)\)-plane, is an \( \alpha \)-quasiperiodic solution to the homogeneous scattering problem (BVP) with \( \alpha = (\alpha_1, 0) \).

Example 2.5. There exists a non-Rayleigh frequency \( \omega \) and a grating \( \tilde{\Gamma} \) represented as the graph of a profile function such that the solutions to (BVP) are non-unique.

Again restrict to gratings invariant in \( x_2 \) direction and consider gratings such that \( \tilde{\Gamma} \cap \{ x_2 = 0 \} \) can be represented as a smooth function \( x_3 = f(x_1) \) of period \( \Lambda_1 = 2\pi \). We seek a special magnetic field \( H^{sc} \) of the form

\[
H^{sc}(x) = \frac{1}{ik} \text{curl } E^{sc}(x) = (0, u^{sc}(x_1, x_3), 0).
\]

Since \( H^{sc} \) should satisfy the curl-curl equation (2.3) in \( \tilde{\Omega} \) with the boundary condition \( \nu \times H^{sc} = 0 \) on \( \tilde{\Gamma} \), we only need to find a non-trivial scalar function \( u^{sc} \) such that

\[
(\Delta + k^2) u^{sc} = 0 \quad \text{in} \quad x_3 > f(x_1), \quad \frac{\partial u^{sc}}{\partial n} = 0 \quad \text{on} \quad x_3 = f(x_1),
\]

\[
u \times u^{sc} = 0 \quad \text{on} \quad \tilde{\Gamma} \cap \{ x_2 = 0 \},
\]

\[
u \cdot u^{sc} = C_{n_1} \exp \left( \left( i [\alpha_n^{(1)} x_1 + \sqrt{k^2 - |\alpha_n^{(1)}|^2} x_3] \right) \right), \quad C_{n_1} \in \mathbb{C}, \quad x_3 > \max_{x_1} f(x_1),
\]

where \( n \in \mathbb{R}^2 \) denotes the normal to the one-dimensional curve \( x_3 = f(x_1) \) in the \((x_1, x_3)\)-plane. This case is just the TM polarization of (BVP). It follows from [21] that exponentially decaying solutions (surface waves) to the above Neumann boundary value problem exist for a broad class of grating profiles that are given by the graphs of smooth functions. Thus, we obtain a TM polarized solution \( H^{sc} \) which is transversal to the \((x_1, x_3)\)-plane, and have constructed a non-trivial solution

\[
E^{sc}(x) = -\frac{1}{ik} \text{curl } H^{sc}(x)
\]

\[
= \frac{1}{k} \sum_{n_1 \in \mathbb{Z}, n_2 = 0} C_{n_1} \left( \sqrt{k^2 - |\alpha_n^{(1)}|^2}, 0, -\alpha_n^{(1)} \right)^\top \exp \left( i [\alpha_n^{(1)} x_1 + \sqrt{k^2 - |\alpha_n^{(1)}|^2} x_3] \right)
\]

to the homogenous problem of (BVP).
Note that the last two examples in the non-resonance case are obtained only if the grating surface \(\tilde{\Gamma}\) remains constant in \(x_2\)-direction. Similar non-trivial solutions can be constructed for biperiodic structures only varying in \(x_1\)-direction. However, we do not have a corresponding example for the diffraction gratings that vary in two orthogonal directions. It remains an interesting question that under what kind of geometry conditions imposed on \(\tilde{\Gamma}\) the uniqueness to (BVP) holds. Although there is no uniqueness in the general case, we can prove the existence of solutions to (BVP) for any wavenumber \(k \in \mathbb{R}^+\). This will be done in the subsequent sections.

### 3 An equivalent variational formulation

The goal of this section is to propose a variational formulation equivalent to (BVP). We begin with the fact that any column vector \(E_n \in \mathbb{C}^3\) satisfying \((\alpha_n, \beta_n)^\top \perp E_n\) for some \(n = (n_1, n_2)^\top \in \mathbb{Z}^2\) can be represented as a linear combination of two vectors \(E_{n,1}, E_{n,2} \in \mathbb{C}^3\):

\[
E_n = C_{n,0} E_{n,0} + C_{n,1} E_{n,1}, \quad C_{n,0}, C_{n,1} \in \mathbb{C},
\]

where

\[
E_{n,0} := \begin{cases} 
(-\alpha_n^{(2)}, \alpha_n^{(1)}, 0)^\top /|\alpha_n| \in \mathbb{S}^2, & \text{if } |\alpha_n| \neq 0, \\
(0, 1, 0)^\top, & \text{else},
\end{cases}
\]  

(3.9)

\[
E_{n,1} := \begin{cases} 
\frac{|\alpha_n|}{h_n} (\alpha_n, \beta_n)^\top \times E_{n,0} = (-\alpha_n^{(1)} \beta_n, -\alpha_n^{(2)} \beta_n, |\alpha_n|^2)^\top / h_n, & \text{if } |\alpha_n| \neq 0, \\
(-1, 0, 0)^\top, & \text{else},
\end{cases}
\]  

(3.10)

with \(h_n := |\alpha_n| \sqrt{|\alpha_n|^2 + |\beta_n|^2}\). Obviously, it holds that \((\alpha_n, \beta_n)^\top \perp E_{n,l}, |E_{n,l}| = 1\) for \(l = 0, 1, n \in \mathbb{Z}^2\). One can observe further that \(E_{n,1} \in \mathbb{S}^2\) if \(\beta_n \in \mathbb{R}\), and that \(E_{n,1} = e_3\) if \(\beta_n = 0\). The above decomposition of \(E_n\) allows us to rewrite the Rayleigh expansion (2.5) as

\[
E^{sc}(x) = \sum_{n \in \mathbb{Z}^2, l=1,2} C_{n,l} U_{n,l}(x), \quad U_{n,l} := E_{n,l} \exp \left(i[\alpha_n \cdot x' + \beta_n x_3]\right), \quad C_{n,l} \in \mathbb{C}
\]  

(3.11)

for \(x_3 > \Gamma_{\text{max}}\) (see also [25, Section 2.5]). Define the layer \(\Omega_b^+\) above \(\Gamma_b\) of height one by (see Figure 1)

\[
\Omega_b^+ := \left\{ x \in \mathbb{R}^3 : 0 < x_j < \Lambda_j, j = 1, 2, \ b < x_3 < b + 1 \right\},
\]

and the Sobolev spaces \(Y_l\) by

\[
Y_l := \left\{ U \in H(\text{curl}, \Omega_b^+) : U(x) = \sum_{n \in \mathbb{Z}^2} C_{n,l} U_{n,l}(x), \ C_{n,l} \in \mathbb{C} \right\}, \quad l = 0, 1.
\]

Then we see that the function \(E^+(x) := E^{sc}|_{\Omega_b^+}\) belongs to the space \(Y := Y_0 \oplus Y_1\), and any function in \(Y\) can be analytically extended to the whole half-space \(\{x_3 > \Gamma_{\text{max}}\}\). Hence, the following problem is equivalent to (BVP):

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Figure 1: The geometry of the scattering problem.

(BVP'): Given an incident electric field $E^{\text{in}}$, find $(E, E^+) \in \mathbb{H} := X \times Y$ such that $E$ satisfies the curl-curl equation (2.3) in a distributional sense and the transmission conditions

$$e_3 \times (E - E^{\text{in}} - E^+) = 0, \quad e_3 \times \text{curl} (E - E^{\text{in}} - E^+) = 0 \quad \text{on} \quad \Gamma_b. \quad (3.12)$$

Motivated by the arguments in [25, Section 3.2] and the variational formulation in [19], we propose a new variational formulation that is equivalent to (BVP'). For $(E, E^+), (V, V^+) \in \mathbb{H}$, define the sesquilinear form $a(\cdot, \cdot) : \mathbb{H} \times \mathbb{H} \to \mathbb{C}$ by

$$a((E, E^+), (V, V^+)) = \int_{\Omega_b} \left\{ \text{curl} E \cdot \text{curl} V - k^2 E \cdot \overline{V} \right\} dx - \int_{\Gamma_b} \text{curl} E^+ \cdot e_3 \times \overline{V} ds$$

$$+ \int_{\Gamma_b} e_3 \times (E - E^+) \cdot \text{curl} \overline{V^+} ds$$

$$+ \eta \sum_{n \in \mathcal{T}} \int_{\Gamma_b} e_3 \times (E - E^+) \cdot (e_3 \times U_{n,0}) ds \int_{\Gamma_b} e_3 \times (V - V^+) \cdot (e_3 \times \overline{U_{n,0}}) ds, \quad (3.13)$$

where $\eta > 0$ is a constant factor for mortaring and is normally chosen as a multiple of the reciprocal mesh size (see [19]). Our variational formulation is to find $(E, E^+) \in \mathbb{H}$ such that

$$a((E, E^+), (V, V^+)) = -a((0, E^{\text{in}}), (V, V^+)) \quad \text{for all} \quad (V, V^+) \in \mathbb{H}. \quad (3.14)$$

Note that terms like $\int_{\Gamma_b} \text{curl} E^+ \cdot e_3 \times V ds$ are bounded. Indeed, since $E^+$ is the solution of the curl-curl equation, we get $\text{curl} E^+ \in H(\text{curl}, \Omega^+_b)$ and $(\text{curl} E^+)_{|\Gamma_b} \in H^{-1/2}(\text{Curl}, \Gamma_b)$. Further, note that the third term on the right-hand side of (3.13) has the opposite sign than the corresponding term in [19]. Moreover, the integrals with factor $\eta$ in (3.13) correspond to the following term involved in the variational equation established in [19]:

$$\eta \int_{\Gamma_b} e_3 \times (E - E^+) \cdot e_3 \times (V - V^+) ds. \quad (3.15)$$
The expression (3.15) is not meaningful for general \((E, E^+), (V, V^+) \in \mathbb{H}\), since both \(\nu \times (E - E^+)\) and \(\nu \times (V - V^+)\) belong to \(H_t^{-1/2}(\text{Div}, \Gamma_b)\). Integals like \(\eta \int_{\Gamma_b} \nu \times u \cdot \nu \times \mathbf{r} ds\) in the mortar approach make sense for finite element methods, where \(u\) and \(v\) are finite element functions and \(\eta\) tends to zero with the meshsize. The idea employed in [25] is to replace the integral (3.15) by the Galerkin approximation

\[
\sum_{n,l:\|n\|^2 < N} \eta \int_{\Gamma_b} e_3 \times (E - E^+) \cdot e_3 \times U_{n,l} ds = \int_{\Gamma_b} e_3 \times (V - V^+) \cdot e_3 \times U_{n,l} ds \tag{3.16}
\]

\[
+ \eta \sum_{n: \beta_n = 0} \int_{\Gamma_b} e_3 \times (E - E^+) \cdot U_{n,0} ds = \int_{\Gamma_b} e_3 \times (V - V^+) \cdot U_{n,0} ds, \tag{3.17}
\]

with a sufficiently large number \(N > 0\). It is also mentioned in [25] that the summation in (3.16) and (3.17) can even be restricted to all \(n \in \mathbb{Z}^2\) with \(\beta_n = 0\). In the present paper, we only use the terms of (3.16) with \(\beta_n = 0\), which are the last terms in (3.13).

To prove the equivalence of (3.14) and (BVP'), we need two lemmas.

**Lemma 3.1.**

(i) We have \(\text{curl } U_{n,l} = i(-1)^l U_{n,1-l} \sqrt{|\alpha_n|^2 + |\beta_n|^2} |^{2l} k^l, l = 0, 1\).

(ii) It holds that

\[
e_3 \times U_{n,l} |_{\Gamma_b} = \begin{cases} (-\alpha_n/|\alpha_n|, 0)^T \exp(i\alpha_n \cdot x'), & \text{if } n \in \mathbb{Y}, l = 0, \\
(0, 0, 0)^T, & \text{if } n \in \mathbb{Y}, l = 1, \\
(-1)^l [(e_3 \times U_{n,1-l}) \times e_3] (\cos \theta_n)^{2l-1}, & \text{if } n \notin \mathbb{Y},
\end{cases}
\]

where \(\cos \theta_n := \beta_n/\sqrt{|\beta_n|^2 + |\alpha_n|^2}\).

(iii) The following set is an \(L^2\)-orthogonal basis of the space \(H_t^{-1/2}(\Gamma_b)\):

\[
\{e_3 \times U_{n,l} |_{\Gamma_b} : n \notin \mathbb{Y}, l = 1, 2\} \cup \{e_3 \times U_{n,0} |_{\Gamma_b} : n \in \mathbb{Y}\} \cup \{U_{n,0} |_{\Gamma_b} : n \in \mathbb{Y}\}.
\]

**Proof.** Lemma 3.1 (i) and (ii) can be proved directly using the definitions of \(U_{n,l}\) in (3.11). To prove the third assertion, we define the set

\[
\Pi_n = \begin{cases} \{e_3 \times E_{n,0}, e_3 \times E_{n,1}\}, & \text{if } \beta_n \neq 0, \\
\{e_3 \times E_{n,0}, E_{n,0}\}, & \text{if } \beta_n = 0,
\end{cases}
\]

with \(E_{n,l} \in \mathbb{C}^3\) given in (3.9) and (3.10). Then Lemma 3.1 (iii) simply follows from the definition of \(H_t^{-1/2}(\Gamma_b)\) and the fact that \(\Pi_n\) is a basis of the set \(\{a = (a_1, a_2, 0)^T : a_1, a_2 \in \mathbb{C}\}\) for any \(n \in \mathbb{Z}^2\).

In the following sections we make the convention that a summation over the index \(l\) is always from zero to one.
Lemma 3.2. For any two absolutely convergent Rayleigh series expansion \( U \) and \( V \) defined in a neighborhood of \( \Gamma_b \), there holds

\[
\int_{\Gamma_b} (\text{curl} \, U) \cdot e_3 \times V \, ds = \int_{\Gamma_b} e_3 \times U \cdot (\text{curl} \, [V]_{\text{mo}}) \, ds,
\]

where \([\cdot]_{\text{mo}}\) is a modification operator defined by

\[
\left[ \sum_{n \in \mathbb{Z}^2, l} C_{n,l} U_{n,l} \right]_{\text{mo}} := - \sum_{l,n; \beta_n > 0} C_{n,l} U_{n,l} + \sum_{l,n; \beta_n \in \mathbb{R}} C_{n,l} U_{n,l}.
\]

Proof. See [25, Lemma 3.1].

We are now going to prove

Lemma 3.3. The variational formulation (3.14) and the problem (BVP') are equivalent.

Proof. (i) Assume that \((E, E^+) \in \mathbb{H}\) is a solution to (BVP'). Applying Green’s first vector theorem to the region \(\Omega\) gives

\[
0 = \int_{\Omega_b} \{\text{curl} \, \text{curl} \, E - k^2 E\} \cdot V \, dx = \int_{\Omega_b} \{\text{curl} \, E \cdot \text{curl} \, V - k^2 E \cdot V\} \, dx - \int_{\Gamma_b} e_3 \times V \cdot \text{curl} \, E \, ds
\]

for any \(V \in X\). Note that the integral over \(\Gamma\) vanishes due to the perfectly conducting boundary condition \(\nu \times V = 0\) on \(\Gamma\), and that the integrals over the vertical parts of \(\partial \Omega_b\) cancel because of the \(\alpha\)-quasiperiodicity of \(V\) and \(E\) in \(\Omega_b\). This implies that

\[
\int_{\Omega_b} \{\text{curl} \, E \cdot \text{curl} \, V - k^2 E \cdot V\} \, dx = \int_{\Gamma_b} e_3 \times V \cdot \text{curl} \, E \, ds, \quad \forall \ V \in X. \quad (3.18)
\]

Making use of the identity (3.18) and the transmission conditions in (3.12), we derive from the definition of the sesquilinear form \(a(\cdot, \cdot)\) that

\[
a\left((E, E^+ + E^{\text{in}}), (V, V^+)\right) = \int_{\Gamma_b} \text{curl} \, (E - E^+ - E^{\text{in}}) \cdot e_3 \times V \, ds + \int_{\Gamma_b} e_3 \times (E - E^+ - E^{\text{in}}) \cdot \text{curl} \, V^+ \, ds + \eta \sum_{n \in \Upsilon} \int_{\Gamma_b} e_3 \times (E - E^+ - E^{\text{in}}) \cdot (e_3 \times \overline{U}_{n,0}) \, ds \int_{\Gamma_b} e_3 \times (V - V^+) \cdot (e_3 \times \overline{U}_{n,0}) \, ds
\]

\[
= 0 \quad (3.19)
\]

for any \((V, V^+) \in \mathbb{H}\), i.e., the pair \((E, E^+)\) is a solution to (3.14).
(ii) Suppose that \((E, E^+) \in \mathbb{H}\) is a solution to the variational formulation (3.14). Choose \(V \in X\) with a compact support in the interior of \(\Omega_b\) (i.e. \(\text{Supp}(V) \subset \text{Int} \Omega_b\)) and choose \(V^+ \equiv 0\) in \(Y\). Then,

\[
0 = a\left((E, E^+ + E^{in}), (V, 0)\right) = \int_{\Omega_b} \{\text{curl} E \cdot \text{curl} \nabla - k^2 E \cdot \nabla\} dx
= \int_{\Omega_b} (\text{curl} \, \text{curl} E - k^2 E) \cdot \nabla dx. \tag{3.20}
\]

This implies that \(\text{curl} \, \text{curl} E - k^2 E = 0\) in \(\Omega_b\). It only remains to prove that \(E\) and \(E^+\) satisfy the transmission conditions in (3.12).

Analogously to part (i), multiplying \(V \in X\) to the \(\text{curl-curl}\) equation \(\text{curl} \, \text{curl} E - k^2 E = 0\) in \(\Omega_b\) and then using integration by parts yields the identity (3.18). Combining this identity with the variational formulation (3.14) gives rise to the equality (3.19) for all \((V, V^+) \in \mathbb{H}\). By Lemma 3.1 (ii) and (iii), we consider the Fourier expansion

\[
(E - E^+ - E^{in})_T = \sum_{l, n \in \Upsilon} C_{n,l}(U_{n,l})_T + \sum_{n \in \Upsilon} \left[C_{n,0} U_{n,0} + C_{n,1} e_3 \times U_{n,0}\right] \quad \text{on } \Gamma_b. \tag{3.21}
\]

It then suffices to prove that \(C_{n,l} = 0\) for all \(n \in \mathbb{Z}^2, l = 0, 1\). Indeed, \((E - E^+ - E^{in})_T = 0\) on \(\Gamma_b\) together with (3.19) for all \(V \in X\) would lead to \((\text{curl} \, (E - E^+ - E^{in}))_T = 0\) on \(\Gamma_b\).

First we take \(V \equiv 0\) and \(V^+ = U_{n,1}\) for some \(n \in \Upsilon\) in (3.19). Applying Lemma 3.1 (i) to \(U_{n,1}\) gives the identity \(\text{curl} \, V^+ = -ikU_{n,0}\), and then, using \(e_3 \times U_{n,1} = 0\) for \(n \in \Upsilon\) (see Lemma 3.1 (ii)), we derive from (3.19) that

\[
\int_{\Gamma_b} e_3 \times (E - E^+ - E^{in}) \cdot \nabla U_{n,0} \, ds = 0 \quad \text{if } n \in \Upsilon. \tag{3.22}
\]

Together with (3.21), this implies that \(C_{n,1} = 0\) for \(n \in \Upsilon\).

Next, inserting (3.22) into (3.19) with \(V \equiv 0\) and using Lemma 3.2, we have

\[
0 = \int_{\Gamma_b} e_3 \times (E - E^+ - E^{in}) \cdot \text{curl} \, \nabla^+ \, ds
- \eta \sum_{n \in \Upsilon} \int_{\Gamma_b} (E - E^+ - E^{in})_T \cdot \nabla U_{n,0} \, ds \int_{\Gamma_b} V^+_T \cdot \nabla U_{n,0} \, ds
= \int_{\Gamma_b} \text{curl} \, [(E - E^+ - E^{in})|_{\Gamma_b}]_{mo} \cdot e_3 \times (\nabla^+) \, ds
- \eta \sum_{n \in \Upsilon} \int_{\Gamma_b} (E - E^+ - E^{in})_T \cdot \nabla U_{n,0} \, ds \int_{\Gamma_b} V^+_T \cdot \nabla U_{n,0} \, ds \tag{3.23}
\]

for all \(V^+ \in Y\), where the quantity

\[
[(E - E^+ - E^{in})|_{\Gamma_b}]_{mo} := - \sum_{l, n, \beta_n > 0} C_{n,l} U_{n,l} + \sum_{l, n, \beta_n \in \mathbb{R}} C_{n,l} U_{n,l} \quad \text{on } \Gamma_b
\]
is obtained by firstly extending the series expansion (3.21) to a neighborhood of \( \Gamma_b \) and then applying the modification operator \([ \cdot ]_{mo}\) defined in Lemma 3.2. From Lemma 3.1 (\(i\)) and (\(ii\)), it follows that on \( \Gamma_b \),

\[
\left\{ \text{curl} \left[ (E - E^+ - E^{in})|_{\Gamma_b} \right] \right\}_{mo}
= \sum_{n,l; \beta_n > 0} i(-1)^{l+1} k C_{n,l}(U_{n,1-l})_{\Gamma_b} + \sum_{n,l; \beta_n \notin \mathbb{R}} i(-1)^{l} k^{2l} C_{n,l} \sqrt{|\alpha_n|^2 + |\beta_n|^2} (U_{n,1-l})_{\Gamma_b}
= \sum_{n,l; \beta_n > 0} -ik C_{n,l} (\cos \theta_n)^{1-2l} e_3 \times U_{n,l} + \sum_{n,l; \beta_n \notin \mathbb{R}} ik^{2l} C_{n,l} (\beta_n)^{1-2l} e_3 \times U_{n,l}.
\]

Inserting (3.24) into (3.23) and choosing \( V^+ = U_{n,0} \) for some \( n \notin \Upsilon \), we derive \( C_{n,l} = 0 \). Analogously, the choice of \( V^+ = U_{n,0} \) for some \( n \in \Upsilon \) leads to \( C_{n,0} = 0 \). The proof is thus completed.

**Remark 3.4.** In the non-resonance case, i.e. \( \Upsilon = \emptyset \), the variational formulations (3.14) and (2.7) are equivalent. In fact, if \((E, E^+)\) is a solution to the problem (3.14), then by Lemma 3.3, the transmission conditions in (3.12) hold. Hence, we obtain

\[
0 = a\left( (E, E^+ + E^{in}), (V, V^+) \right) = \int_{\Omega_b} \text{curl} E \cdot \text{curl} \nabla - k^2 E \cdot \nabla dx - \int_{\Gamma_b} \text{curl} (E^+ + E^{in}) \cdot e_3 \times \nabla ds
= \int_{\Omega_b} \text{curl} E \cdot \text{curl} \nabla - k^2 E \cdot \nabla dx - \int_{\Gamma_b} \mathcal{R}(e_3 \times E) \cdot e_3 \times ds
+ \int_{\Gamma_b} \left[ \mathcal{R}(e_3 \times E^{in}) - \text{curl} E^{in} \right] \cdot e_3 \times ds,
\]

which is equivalent to the variational formulation (2.7) involving the Dirichlet-to-Neumann map \( \mathcal{R} \). Note that in the last step of the previous identity we have used the identity

\[
(\text{curl} E^+)_{\Gamma_b} = \mathcal{R}(e_3 \times E^+) = \mathcal{R}(e_3 \times E) - \mathcal{R}(e_3 \times E^{in}) \text{ on } \Gamma_b.
\]

On the other hand, supposing that \( E \in H(\text{curl}, \Omega_b) \) is a solution to (2.7), we extend the scattered field \( E^{sc} := E - E^{in} \) from \( \Omega_b \) to \( x_3 > b \) by the Rayleigh expansion (2.5). Assume that the coefficients \( A_n \) are given by

\[
e_3 \times E^{sc}|_{\Gamma_b^-} = e_3 \times (E - E^{in})|_{\Gamma_b^-} = \sum_{n \in \mathbb{Z}^2} A_n e^{i\alpha_n x'} \in H^{-1/2}_+(\text{Div}, \Gamma_b), \quad A_n \in \mathbb{C}.
\]

Here and in the following, the symbol \( (\cdot)|_{\Gamma_b^-} \) resp. \( (\cdot)|_{\Gamma_b^+} \) denotes the trace obtained from below resp. above \( \Gamma_b \). It follows from the variational formulation (2.7) that \( e_3 \times \text{curl} E^{sc} \times e_3|_{\Gamma_b^-} = \mathcal{R}(e_3 \times E^{sc}|_{\Gamma_b^-}) \). The extension of the series in (3.25) to the upper half space \( x_3 > b \) in form of the Rayleigh expansion (2.5) is

\[
E^{sc}(x) = \sum_{n \in \mathbb{Z}^2} \left[ A_n \times e_3 + \beta_n^{-1}(e_3 \times A_n) \cdot \alpha_n e_3 \right] e^{i\alpha_n x' + i\beta_n (x_3 - b)}, \quad x_3 > b.
\]
Then,

\[ e_3 \times E^{sc}|_{\Gamma_b}^- = e_3 \times E^{sc}|_{\Gamma_b}^+; \quad e_3 \times \text{curl} \, E^{sc} \times e_3|_{\Gamma_b}^+ = \mathcal{R}(e_3 \times E^{sc}|_{\Gamma_b}^+) \]

Setting \( E^+ = E^{sc} \) in \( \Omega_b^+ \), we conclude that \( (E, E^+) \) satisfies the transmission conditions (3.12) and thus is a solution of (3.14).

### 4 Analysis of the variational formulation (3.14)

Since the sesquilinear form \( a(\cdot, \cdot) \) defined in Section 3 is bounded on \( \mathbb{H} \), it obviously generates a continuous linear operator \( A : \mathbb{H} \to \mathbb{H}' \) satisfying

\[ a\left( (E, E^+), (V, V^+) \right) = \left\langle A(E, E^+), (V, V^+) \right\rangle_{\Omega_b^+} \]  (4.26)

Here \( \mathbb{H}' \) denotes the dual of the space \( \mathbb{H} \) with respect to the duality \( \langle \cdot, \cdot \rangle_{\Omega_b^+} \) extending the scalar product in \( L^2(\Omega_b^+) \). The aim of this section is to prove

**Theorem 4.1.** The operator \( A \) defined by (4.26) is a Fredholm operator with index zero.

To prove Theorem 4.1, we need several auxiliary lemmas. We first prove a periodic analogue of the Hodge-decomposition of \( X \), following the argument in [22, Theorem 4.3]. See also [1, 2, 3, 17] for other Hodge-decompositions of the Sobolev spaces in periodic structures. Define

\[
X_1 := \{ \nabla p : p \in H^1(\Omega_b), p = 0 \text{ on } \Gamma \}, \quad X_0 := \{ E_0 \in X : \int_{\Omega_b} \nabla p \cdot \overline{E_0} \, dx = 0 \text{ for all } \nabla p \in X_1 \}.
\]

**Lemma 4.2.** We have \( X = X_0 \oplus X_1 \) with the subspaces \( X_0 \) and \( X_1 \) orthogonal in \( L^2(\Omega_b)^3 \times L^2(\Omega_b^+)^3 \). Moreover, \( \text{div} \, E_0 = 0 \) and \( (e_3 \cdot E_0)|_{\Gamma_b} = 0 \) for any \( E_0 \in X_0 \), and the space \( X_0 \) is compactly embedded into \( L^2(\Omega_b)^3 \).

**Proof.** Define the bilinear form \( b(\cdot, \cdot) : X_1 \times X_1 \to \mathbb{C} \) by

\[ b(E, V) := \int_{\Omega_b} \{ \text{curl} \, E \cdot \text{curl} \, V + E \cdot \overline{V} \} \, dx, \quad E, V \in X. \]

Then, for \( \nabla p \in X_1 \), it holds that

\[ b(\nabla p, \nabla p) = ||\nabla p||_{L^2(\Omega_b)}^2 = ||\nabla p||_X^2. \]

Thus, for every \( E \in X \) there exists a unique solution \( \nabla p \in X_1 \) such that

\[ b(\nabla p, \nabla \xi) = b(E, \nabla \xi), \quad \forall \ nabla \xi \in X_1. \]  (4.27)
Let $E_0 := E - \nabla p$. Using integration by parts and the quasiperiodicity of $E_0$ and $\xi$ in $\Omega_b$, it follows from (4.27) that

$$0 = \int_{\Omega_b} E_0 \cdot \nabla \xi \; dx = -\int_{\Omega_b} \xi \; \text{div} E_0 \; dx + \int_{\Gamma_b} \xi \; e_3 \cdot E_0 \; ds,$$

for any $\nabla \xi \in X_1$.

This implies that $X = X_1 + X_0$ and $\text{div} E_0 = 0$, $(e_3 \cdot E_0)_{|\Gamma_b} = 0$. On the other hand, if $\nabla q \in X_0 \cap X_1$, then the definition of $X_0$ implies that $\int_{\Omega_b} \nabla p \cdot \nabla q \; dx = 0$. Setting $p = q$, we get $\nabla q = 0$, i.e., $X_0 \cap X_1 = \emptyset$. Finally, the compact imbedding of $X_0$ into $L^2(\Omega_b)^3$ follows from [22, Corollary 3.49] (see also [3, Lemma 3.2]).

By Lemma 4.2 and the definitions of $Y_l$, we can decompose our space $\mathbb{H}$ into four subspaces. For $(E, E^+), (V, V^+) \in \mathbb{H}$, we may assume that

$$E = \nabla p + E_0, \quad E^+ = E_0^+ + E_1^+, \quad \text{where} \quad \nabla p \in X_1, \quad E_0 \in X_0, \quad E_l^+ \in Y_l, \quad l = 1, 2,$$

$$V = \nabla \xi + V_0, \quad V^+ = V_0^+ + V_1^+, \quad \text{where} \quad \nabla \xi \in X_1, \quad V_0 \in X_0, \quad V_l^+ \in Y_l, \quad l = 1, 2.$$

For the convenience to analyze the form $a$, we define several sesquilinear forms $a_j$ with $j = 1, 2, \cdots, 6$. Let

$$a_1(\nabla p, \nabla \xi) := k^2 \int_{\Omega_b} \nabla p \cdot \nabla \xi \; dx, \quad \forall \nabla p, \nabla \xi \in X_1,$$

$$a_2(E_0, V_0) := \int_{\Omega_b} \{ \text{curl} E_0 \cdot \text{curl} V_0 - k^2 E_0 \cdot V_0 \} \; dx, \quad \forall E_0, V_0 \in X_0,$$

$$a_3(E_0^+, V_0^+) := \int_{\Gamma_b} e_3 \times E_0^+ \cdot \text{curl} V_0^+ \; ds, \quad \forall E_0^+, V_0^+ \in Y_0,$$

$$a_4(E_1^+, V_1^+) := \int_{\Gamma_b} e_3 \times E_1^+ \cdot \text{curl} V_1^+ \; ds, \quad \forall E_1^+, V_1^+ \in Y_1,$$

and let

$$a_5((E, E^+), (V, V^+)) := \int_{\Gamma_b} e_3 \times E \cdot \text{curl} V^+ \; ds,$$

$$a_6((E, E^+), (V, V^+))$$

$$:= \eta \sum_{n \in \mathbb{Y}} \left\{ \int_{\Gamma_b} e_3 \times (E - E^+) \cdot (e_3 \times \overline{U}_{n,0}) \; ds \int_{\Gamma_b} e_3 \times (V - V^+) \cdot (e_3 \times \overline{U}_{n,0}) \; ds \right\}$$

for any $(E, E^+), (V, V^+) \in \mathbb{H}$. For brevity we write

$$a_5((E, E^+), (V, V^+)) = a_5(E, V^+). \quad (4.28)$$

**Lemma 4.3.** For any $\nabla \xi \in X_1$ and $V_0^+ \in Y_0$, we have $a_5(\nabla \xi, V_0^+) = 0$. 

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Proof. From the definition of $Y$ and $Y_0$ we conclude that $Y_0$ is the subspace of all vector functions $V^+ \in Y$ with $e_3 \cdot V^+ = 0$. Therefore it suffices to prove
\[
\int_{\Gamma_b} [e_3 \times \nabla \xi] \cdot \text{curl} \, \overline{V}^+ \, ds = k^2 \int_{\Gamma_b} e_3 \cdot V^+ \, \xi \, ds. \tag{4.29}
\]
Note that the right-hand side of (4.29) is a continuous functional of $V^+$ and $\xi$. Indeed, from
\[V^+ \in L^2(\Omega_b^+) \text{ and } 0 = \nabla \cdot V^+ \in L^2(\Omega_b^+),\]
we conclude $e_3 \cdot V^+ \in H^{-1/2}(\Gamma_b)$, and $\xi \in H^{1/2}(\Gamma_b)$ follows from $\xi \in H^1(\Omega_b)$. Knowing the continuity, it suffices to prove (4.29) for a dense subset, e.g., for a truncated Rayleigh expansion $V^+$ and smooth $\xi$. We conclude
\[
\int_{\Gamma_b} [e_3 \times \nabla \xi] \cdot \text{curl} \, \overline{V}^+ \, ds = -\int_{\Gamma} [\nu \times \nabla \xi] \cdot \text{curl} \, \overline{V}^+ \, ds + \int_{\Omega_b} [\text{curl} \, \nabla \xi] \cdot \text{curl} \, \overline{V}^+ \, ds
\]
\[= k^2 \int_{\Omega_b} [\nabla \xi] \cdot \text{curl} \, \overline{V}^+ \, ds,
\]
where we have used that the tangential derivative $\nu \times \nabla \xi$ of the function $\xi$ with $\xi|_\Gamma = 0$ vanishes. Using $\nabla \cdot V^+ = 0$, we continue
\[
\int_{\Omega_b} [\nabla \xi] \cdot \text{curl} \, \overline{V}^+ \, ds = k^2 \int_{\Omega_b} \nabla \cdot [\xi \overline{V}^+] \, ds = k^2 \int_{\Gamma_b} \xi e_3 \cdot \overline{V}^+ \, ds + k^2 \int_{\Gamma} \xi \nu \cdot \overline{V}^+ \, ds
\]
\[= k^2 \int_{\Gamma_b} \xi e_3 \cdot \overline{V}^+ \, ds
\]
and the proof is completed. \qed

Using the Lemmas 3.1, 4.2 and 4.3, it follows from the definition of $a$ that (see Table 1)
\[
a((E, E^+), (V, V^+)) = a((\nabla p + E_0, E_0^+ + E_1^+), (\nabla \xi + V_0, V_0^+ + V_1^+))
\]
\[= \int_{\Omega} \{ \text{curl} \, E_0 \cdot \text{curl} \, \overline{V}_0 - k^2 E_0 \cdot \overline{V}_0 - k^2 \nabla p \cdot \nabla \xi \} \, dx - \int_{\Gamma_b} \text{curl} \, E_0^+ \cdot e_3 \times \overline{V}_0 \, dx
\]
\[- \int_{\Gamma_b} \text{curl} \, E_0^+ \cdot e_3 \times \nabla \xi \, dx - \int_{\Gamma_b} \text{curl} \, E_0^+ \cdot e_3 \times \nabla \xi \, dx + \int_{\Gamma_b} e_3 \times E_0 \cdot \text{curl} \, \overline{V}_0^+ \, dx
\]
\[+ \int_{\Gamma_b} e_3 \times E_0 \cdot \text{curl} \, \overline{V}_1^+ \, dx + \int_{\Gamma_b} e_3 \times \nabla p \cdot \text{curl} \, \overline{V}_1^+ \, dx - \int_{\Gamma_b} e_3 \times E_0^+ \cdot \text{curl} \, \overline{V}_0^+ \, dx
\]
\[- \int_{\Gamma_b} e_3 \times E_1^+ \cdot \text{curl} \, \overline{V}_1^+ \, dx + \, a_0((E, E^+), (V, V^+))
\]
\[= -a_1(\nabla p, \nabla \xi) + a_2(E_0, V_0) - a_3(E_0^+ + V_0^+) - a_4(E_1^+, V_1^+) + a_5(E_0, V_0^+)
\]
\[-a_5(V_0, E_0^+) + a_5(E_0, V_1^+) - a_5(V_0, E_1^+) + a_5(\nabla p, V_1^+) - a_5(\nabla \xi, E_1^+)
\]
\[+ a_6((E, E^+), (V, V^+)) . \tag{4.30}
\]

Definition 4.4. A bounded sesquilinear form $l(\cdot, \cdot)$ given on some Hilbert space $X$ is called strongly elliptic if there exists a compact form $\tilde{l}(\cdot, \cdot)$ and a constant $c > 0$ such that
\[
\text{Re} \, l(u, u) \geq c \|u\|_X^2 - \tilde{l}(u, u), \quad \forall \, u \in X.
\]
| $\mathbb{H}_0$ | $X_0(E_0)$ | $a_2(E_0, V_0)$ | $a_5(E_0, V_0^+)$ | $0$ | $a_5(E_0, V_1^+)$ |
| | | | | | |
| $Y_0(E_0^+)$ | $-a_5(V_0, E_0^-)$ | $-a_5(E_0^+, V_0^+)$ | $0$ | $0$ |

| $\mathbb{H}_1$ | $X_1(\nabla p)$ | $0$ | $0$ | $-a_1(\nabla p, \nabla \xi)$ | $a_5(\nabla p, V_1^+)$ |
| | $Y_1(E_1^+)$ | $-a_5(V_0, E_1^-)$ | $0$ | $-a_5(\nabla \xi, E_1^-)$ | $-a_4(E_1^+, V_1^+)$ |

Table 1: The diagram for the sesquilinear form $a - a_6$ over $\mathbb{H} = X \times Y$.

Obviously, $a_1$ is coercive on $X_1$ and by Lemma 4.2 the sesquilinear form $a_2$ is strongly elliptic over $X_0$. In addition, $a_6$ is a compact form over $\mathbb{H}$, since it corresponds to a finite rank operator over $\mathbb{H}$. To demonstrate the Fredholm property of the sesquilinear form $a$, we now need to study the other forms $a_3, a_4$ and $a_5$.

**Lemma 4.5.** There exist compact forms $\tilde{a}_3 : Y_0 \times Y_0 \to \mathbb{C}$ and $\tilde{a}_4 : Y_1 \times Y_1 \to \mathbb{C}$ such that

\[
- \text{Re } a_3(\cdot, \cdot) \geq C \| \cdot \|^2_{H(\nabla, \Omega_b^+)} - \tilde{a}_3(\cdot, \cdot),
\]

\[
- \text{Re } a_4(\cdot, \cdot) \geq C \| \cdot \|^2_{H(\nabla, \Omega_b^+)} - \tilde{a}_4(\cdot, \cdot),
\]

for some constant $C > 0$, i.e., the sesquilinear forms $-a_3$ and $a_4$ are strongly elliptic over $Y_0$ and $Y_1$, respectively.

**Proof.** Recall that the functions $U_{n,l}$ defined in (3.11) are basis functions of the space $Y_1$, $l = 1, 2$. It is easy to check that

\[
\int_{\Omega_b^+} U_{n,l} \cdot \overline{U}_{n',l'} \, dx = \delta_{n,n'} \delta_{l,l'} \Lambda_1 \Lambda_2 \int_{\beta}^{b+1} \exp(i\beta_n x_3) \exp(-i\overline{\beta}_n x_3) dx_3
\]

\[
= \begin{cases} 
\delta_{n,n'} \delta_{l,l'} \Lambda_1 \Lambda_2 & \text{if } \beta_n \in \mathbb{R}, \\
\delta_{n,n'} \delta_{l,l'} e^{-2|\beta_n|}(1 - e^{-2|\beta_n|})(2|\beta_n|)^{-1} \Lambda_1 \Lambda_2 & \text{if } \beta_n \notin \mathbb{R},
\end{cases}
\]

and that, by using Lemma 3.1,

\[
\int_{\Omega_b^+} \text{curl } U_{n,l} \cdot \text{curl } \overline{U}_{n',l'} \, dx = \delta_{n,n'} \delta_{l,l'} k^d \sqrt{\alpha_n^2 + |\beta_n|^2} - d \int_{\Omega_b^+} U_{n,1-l} \cdot \overline{U}_{n,1-l} \, dx.
\]

Therefore, we can represent the $H(\nabla, \Omega_b^+)$-norm of $U_{n,l}$ as $0, 1$ as

\[
\|U_{n,0}\|_{H(\nabla, \Omega_b^+)}^2 = \begin{cases} 
(1+k^2)\Lambda_1 \Lambda_2 & \text{if } \beta_n \in \mathbb{R}, \\
e^{-2|\beta_n|}(1-e^{-2|\beta_n|})(1+2|\beta_n|^2+k^2)(2|\beta_n|)^{-1} \Lambda_1 \Lambda_2 & \text{if } \beta_n \notin \mathbb{R},
\end{cases}
\]

\[
\|U_{n,1}\|_{H(\nabla, \Omega_b^+)}^2 = \begin{cases} 
(1+k^2)\Lambda_1 \Lambda_2 & \text{if } \beta_n \in \mathbb{R}, \\
e^{-2|\beta_n|}(1-e^{-2|\beta_n|})(1+\frac{k^4}{2|\beta_n|^2+k^2})(2|\beta_n|)^{-1} \Lambda_1 \Lambda_2 & \text{if } \beta_n \notin \mathbb{R}.
\end{cases}
\]
On the other hand, using simple calculations, we have, for \( n \notin Y \),
\[
\int_{\Gamma_b} e_3 \times U_{n,l} \cdot \text{curl} \nabla_{n,l} \, ds = (-i)(\cos \theta_n)^{2l-1} k^{2l} \sqrt{\alpha_n^2 + |\beta_n|^2}^{2l-2l} \int_{\Gamma_b} |e_3 \times U_{n,1-l}|^2 \, ds,
\]
by Lemma 3.1 (i) and (ii). Furthermore,
\[
\int_{\Gamma_b} |e_3 \times U_{n,1-l}|^2 \, ds = \begin{cases} 
|e^{i\beta_n b}|^2 \Lambda_1 \Lambda_2 & \text{if } l = 1, \\
|e^{i\beta_n b}|^2 \Lambda_1 \Lambda_2 | \cos \theta_n |^2 & \text{if } l = 0,
\end{cases}
\]
by the definitions of \( U_{n,l} \) given in (3.11). Combining the previous two equalities yields
\[
\text{Re} \left\{ \int_{\Gamma_b} e_3 \times U_{n,1} \cdot \text{curl} \nabla_{n,1} ds \right\} = \begin{cases} 
0 & \text{if } \beta_n \in \mathbb{R}, \\
\frac{|\beta_n| k^2}{2|\beta_n|^2 + k^2} e^{-2|\beta_n|b} \Lambda_1 \Lambda_2 & \text{if } \beta_n \notin \mathbb{R},
\end{cases}
\]
\[
\text{Re} \left\{ \int_{\Gamma_b} e_3 \times U_{n,0} \cdot \text{curl} \nabla_{n,0} ds \right\} = \begin{cases} 
0 & \text{if } \beta_n \in \mathbb{R}, \\
-|\beta_n| e^{-2|\beta_n|b} \Lambda_1 \Lambda_2 & \text{if } \beta_n \notin \mathbb{R}.
\end{cases}
\]
Since \(|\beta_n| \sim \sqrt{1 + |n|^2} \) as \(|n| \to +\infty\), there holds
\[
-\text{Re} \left\{ \int_{\Gamma_b} e_3 \times U_{n,0} \cdot \text{curl} \nabla_{n,0} ds \right\} \geq C \|U_{n,0}\|^2_{H(\text{curl}, \Omega_h)},
\] (4.34)
\[
\text{Re} \left\{ \int_{\Gamma_b} e_3 \times U_{n,1} \cdot \text{curl} \nabla_{n,1} ds \right\} \geq C \|U_{n,1}\|^2_{H(\text{curl}, \Omega_h)},
\] (4.35)
whenever \( \beta_n \notin \mathbb{R} \), with \( C > 0 \) being a constant independent of \( l \) and \( n \). Therefore, given \( E_0^+ = \sum_{n \in \mathbb{Z}^2} C_{n,0} U_{n,0} \in Y_0 \), we deduce from (4.34) that
\[
-\text{Re} \left( a_3(E_0^+, E_0^+) \right) = - \sum_{n \in \mathbb{Z}^2} |C_n,0|^2 \text{Re} \left[ \int_{\Gamma_b} e_3 \times U_{n,0} \cdot \text{curl} \nabla_{n,0} ds \right]
\geq C \sum_{\beta_n \notin \mathbb{R}} |C_n,0|^2 \|U_{n,0}\|^2_{H(\text{curl}, \Omega_h^+)}
= C||E_0^+||^2_{H(\text{curl}, \Omega_h^+)} - C \sum_{\beta_n \in \mathbb{R}} |C_n,0|^2 \|U_{n,0}\|^2_{H(\text{curl}, \Omega_h^+)}
= C||E_0^+||^2_{H(\text{curl}, \Omega_h^+)} - \tilde{a}_3(E_0^+, E_0^+).
\]
Since the set \( \{ n \in \mathbb{Z}^2 : \beta_n \in \mathbb{R} \} \) consists of a finite number of indices, the form \( \tilde{a}_3(\cdot, \cdot) : Y_0 \times Y_0 \to \mathbb{R} \) is compact. Thus the sesquilinear form \(-a_3\) is strongly elliptic over \( Y_0 \). The proof for \( a_4 \) can be carried out analogously by employing (4.35). \( \Box \)

**Remark 4.6.** By the definition of \( U_{n,l} \), one can check that, for the component-wise gradient \( \nabla U_{n,l} \),
\[
\int_{\Omega_h^+} |\nabla U_{n,l}|^2 \, dx = (|\alpha_n|^2 + |\beta_n|^2) \int_{\Omega_h^+} |U_{n,l}|^2 \, dx.
\] (4.36)
Thus, comparing (4.36) with (4.33) leads to
\[ ||U_{n,0}||_{H^1(\Omega_b^+)^3}^2 = ||U_{n,0}||_{H(\text{curl},\Omega_b^+)}^2.\]

This implies that the \(H^1\)-norm and \(H(\text{curl})\)-norm of the elements from \(Y_0\) are identical, i.e.,
\[ ||E_0^+||_{H^1(\Omega_b^+)^3} = ||E_0^+||_{H(\text{curl},\Omega_b^+)}, \quad \text{if} \quad E_0^+ \in Y_0.\]

However, this is not true for the space \(Y_1\).

We turn to investigating the properties of \(a_5\) defined in (4.28).

**Lemma 4.7.** The sesquilinear form \(a_5\) is compact over \(X_0 \times Y_1\).

**Proof.** For \(V_1^+ \in Y_1 \subset Y\), define the operator \(J(V_1^+) := \text{curl } V_1^+\). Obviously, we get
\[
||J(V_1^+)||_{L^2(\Omega_b^+)^3} \leq ||V_1^+||_{H(\text{curl},\Omega_b^+)^3},
\]
\[
||\text{curl } J(V_1^+)||_{L^2(\Omega_b^+)^3} = ||k^2 V_1^+||_{L^2(\Omega_b^+)^3} \leq k^2 ||V_1^+||_{H(\text{curl},\Omega_b^+)^3}.
\]

Hence, by Lemma 3.1 (i), \(J\) is a bounded linear map from \(Y_1\) into \(Y_0\). In view of the equivalence of the norms \(||JV_1^+||_{H(\text{curl},\Omega_b^+)^3}\) and \(||JV_1^+||_{H(\Omega_b^+)^3}\) (see Remark 4.6), we see that \(J\) is also bounded from the subspace \(Y_1\) of \(H(\text{curl},\Omega_b^+)^3\) into the subspace \(Y_0\) of \(H^1(\Omega_b^+)^3\), with the trace \(J(V_1^+)|_{\Gamma_b} \in H^{1/2}(\Gamma_b)^3\). Thus, there exists an extension \(W\) of \((\text{curl } V_1^+)|_{\Gamma_b}\) from \(H^{1/2}(\Gamma_b)^3\) into \(H^1(\Omega_b)^3\) such that \(W = \text{curl } V_1^+\) on \(\Gamma_b\) and \(\nu \times W = 0\) on \(\Gamma\). Using integration by parts,
\[
a_5(E_0, V_1^+) = \int_{\Gamma_b} e_3 \times E_0 \cdot J(V_1^+) ds = \int_{\Gamma_b} e_3 \times E_0 \cdot W ds = \int_{\Omega_b} \{\text{curl } E_0 \cdot W - E_0 \cdot \text{curl } W\} dx.
\]

From the compact imbedding of \(W \in H^1(\Omega_b)^3\) into \(L^2(\Omega_b)^3\) and that of \(E_0 \in X_0\) into \(L^2(\Omega_b)^3\), it follows that the sesquilinear form \(a_5(E_0, V_1^+)\) is compact over \(X_0 \times Y_1^+\).

Combining Lemmas 4.2, 4.3, 4.5 and 4.7, we are now in a position to prove the Fredholm property of the variational formulation (3.14).

**Proof of Theorem 4.1.** It suffices to verify that the sesquilinear form \(a - a_6\) is Fredholm over \(\mathbb{H}\) with index zero. To do this, we define the spaces \(\mathbb{H}_j = X_j \oplus Y_j\) for \(j = 0, 1\), so that we can rewrite \(\mathbb{H} = X \times Y = \mathbb{H}_1 \times \mathbb{H}_2\). Define the sesquilinear forms
\[
b_0((E_0, E_0^+), (V_0, V_0^+)) := a_2(E_0, V_0) - a_3(E_0^+, V_0^+) + a_5(E_0, V_0^+) - a_5(V_0, E_0^+),
\]
for \((E_0, E_0^+), (V_0, V_0^+) \in \mathbb{H}_0\) and
\[
b_1((\nabla p, E_1^+), (\nabla \xi, V_1^+)) := -a_1(\nabla p, \nabla \xi) - a_4(E_1^+, V_1^+) + a_5(\nabla p, V_1^+) - a_5(\nabla \xi, E_1^+),
\]
for \((\nabla p, E_1^+), (\nabla \xi, V_1^+) \in \mathbb{H}_1\).
for all \((\nabla p, E^+_1)\) and \((\nabla \xi, V^+_1)\) in \(H_1\). Now split the form in Table 1 in blocks corresponding to the splitting \(H = H_1 \times H_2\). Then the restriction to \(H_1\) is the form \(b_0\) with the strongly elliptic quadratic form

\[
\text{Re} \ b_0 \left( (E_0, E^+_0), (E_0, E^+_0) \right) = a_2(E_0, E_0) - a_3(E^+_0, E^+_0).
\]

The restriction to \(H_1\) is the form \(b_1\), and \(-b_1\) has the strongly elliptic quadratic form

\[
\text{Re} \ b_1 \left( (\nabla p, E^+_1), (\nabla p, E^+_1) \right) = a_1(\nabla p, \nabla p) + a_4(E^+_1, E^+_1).
\]

Consequently, the diagonal blocks of the \(2 \times 2\) splitting correspond to Fredholm operators with index zero. On the other hand, the full form in Table 1 differs from the diagonal block matrix only by compact terms. Hence the form \(a\) generates a Fredholm operator with index zero.

5 Proof of Theorem 2.1

Since the problem (BVP') and (3.14) are equivalent (see Lemma 3.3), to prove Theorem 2.1 we only need to prove the existence of solutions to (3.14) with \(E^\text{in}\) replaced by \(E^\text{gen}\) given in (2.6). Consider the homogenous adjoint problem of the variational formulation (3.14): find \((V, V^+) \in H\) such that

\[
a\left( (W, W^+), (V, V^+) \right) = 0 \quad (5.37)
\]

for all \((W, W^+) \in H\). By the Fredholm alternative, it suffices to verify that

\[
a\left( (0, E^\text{gen}), (V, V^+) \right) = 0
\]

for any solution \((V, V^+)\) to (5.37). The following lemma describes properties of the solution \((V, V^+)\), which will be used later for proving Theorem 2.1.

**Lemma 5.1.** Assume that the pair \((V, V^+) \in H\) is a solution to the homogeneous adjoint problem (5.37). Then

\[
V_T|_{\Gamma_b}, \ (\text{curl} \ V^+)|_{\Gamma_b} \in \text{Span}\left\{ (U_{n,l})|_{\Gamma_b} : \beta_n \notin \mathbb{R}, l = 1, 2 \right\} \cup \left\{ U_{n,0}|_{\Gamma_b} : \beta_n = 0 \right\} \quad (5.38)
\]

**Proof.** Analogously to the proof of (3.20), one can prove that \(\text{curl} \ \text{curl} \ V - k^2 V = 0\) holds in \(\Omega_b\), leading to the identity (3.18) with \((V, E)\) replaced by \((W, V)\). By the definition of \(a(\cdot, \cdot)\), there holds

\[
0 = a\left( (W, W^+), (V, V^+) \right) = \int_{\Gamma_b} \left\{ e_3 \times W \cdot \text{curl} \, \nabla - \text{curl} \, W^+ \cdot e_3 \times \nabla \right\} \, ds + \int_{\Gamma_b} e_3 \times (W - W^+) \cdot \text{curl} \, \nabla^+ \, ds
\]

\[
+ \eta \sum_{n \in Y} \int_{\Gamma_b} e_3 \times (W - W^+) \cdot e_3 \times U_{n,0} \, ds \int_{\Gamma_b} e_3 \times (V - V^+) \cdot e_3 \times U_{n,0} \, ds \quad (5.39)
\]
for all \((W, W^+) \in \mathbb{H}\). In the following we will prove (5.38) by choosing different test functions \((W, W^+) \in \mathbb{H}\).

(i) Choose \(W \equiv 0, W^+ = U_{n,0}\) for some \(n \in \Upsilon\) in (5.39). Since \((\text{curl } U_{n,0})_T = 0\) on \(\Gamma_b\) (see Lemma 3.1), simple calculations leads to

\[
\int_{\Gamma_b} \left[ \text{curl } V^+ + \eta \Lambda_1 \Lambda_2 e_3 \times (V - V^+) \right] \cdot e_3 \times \overline{U}_{n,0} \, ds = 0.
\]

However, one can verify using Lemma 3.1 (i) and (ii) that

\[
\int_{\Gamma_b} \{\text{curl } V^+ \cdot e_3 \times \overline{U}_{n,0}\} \, ds = 0 \quad \text{for } V^+ \in Y.
\]

Hence,

\[
\int_{\Gamma_b} e_3 \times (V - V^+) \cdot e_3 \times \overline{U}_{n,0} \, ds = 0, \quad \text{if } n \in \Upsilon. \tag{5.40}
\]

(ii) Choose \(W \equiv 0\) and \(W^+ = U_{n,1}\) for some \(n \in \Upsilon\) in (5.39). Making use of \(e_3 \times U_{n,1} = 0\) for \(n \in \Upsilon\), we derive from (5.39) that

\[
\int_{\Gamma_b} \{\text{curl } U_{n,1} \cdot e_3 \times \overline{V}\} \, ds = 0.
\]

This together with Lemma 3.1 (i) gives the relation

\[
\int_{\Gamma_b} \{e_3 \times \overline{V} \cdot U_{n,0}\} \, ds = 0 \quad \text{if } n \in \Upsilon. \tag{5.41}
\]

(iii) Inserting (5.40) and (5.41) into (5.39) with \(W \equiv 0\) and taking into account Lemma 3.2, we obtain

\[
\int_{\Gamma_b} \{\text{curl } [V]_{mo} + \text{curl } V^+\} \cdot e_3 \times W^+ \, ds = 0
\]

for all \(W^+ \in Y\). By Lemma 3.1 (iii), the above identity implies that

\[
\left\{ (\text{curl } [V]_{mo})_T + (\text{curl } V^+)_T \right\} |_{\Gamma_b} \in \text{Span}\{U_{n,0} : n \in \Upsilon\}. \tag{5.42}
\]

Since \(V^+ \in Y\), we have \(\text{curl } V^+ \in H(\text{curl}, \Omega^+_b)\) and thus the trace \((\text{curl } V^+_T)\) on \(\Gamma_b\) belongs to \(H^{1/2}_{-} (\text{Curl}, \Gamma_b)\). Using Lemma 3.1 (iii), we may assume that on \(\Gamma_b\)

\[
(\text{curl } V^+_T) = \sum_{n: \beta_n = 0} B_{n,0} U_{n,0} + \sum_{n: \beta_n \neq 0} B_{n,1} e_3 \times U_{n,1} + \sum_{l,n: \beta_n \neq 0} B_{n,l} e_3 \times U_{n,l}
\]

with \(B_{n,l} \in \mathbb{C}\). Combining the previous two formulas, we deduce from the definition of the modification operator \([\cdot]_{mo}\) in Lemma 3.2 that \(B_{n,1} = 0\) for \(\beta_n = 0\), and that

\[
\left(\text{curl } [V]_{mo}\right)_T + \sum_{l,n: \beta_n \neq 0} B_{n,l} e_3 \times U_{n,l} = 0 \quad \text{on } \Gamma_b.
\]
Therefore,
\[
(curl \, V^+)_{\Gamma} = \sum_{n: \beta_n = 0} B_{n,0} U_{n,0} - \left( curl \, [V]_{mo} \right)_{\Gamma} \quad \text{on } \Gamma_b.
\]  
(5.43)

(iv) Inserting (5.40) and (5.41) into (5.37) with \(W^+ = 0\) and \(W = V\), we find (see (3.13) with \(E^+ \equiv 0\) and \(E = V\))
\[
0 = \text{Im} \, a(\mathbf{V}, 0), (\mathbf{V}, V^+) = \text{Im} \int_{\Gamma_b} e_3 \times \mathbf{V} \cdot \left( curl \, V^+ \right) \, ds,
\]  
(5.44)

where the function \( (curl \, V^+)_{\Gamma} \) is given in (5.43). According to Lemma 3.1 (iii), we may represent \( e_3 \times V \mid_{\Gamma_b} \) as
\[
e_3 \times V = \sum_{l,n: \beta_n \neq 0} C_{n,l} e_3 \times U_{n,l} + \sum_{n: \beta_n = 0} \{ C_{n,0} e_3 \times U_{n,0} + C_{n,1} U_{n,0} \}, \quad C_{n,l} \in \mathbb{C},
\]
on \(\Gamma_b\). However, by (5.41) there holds \( C_{n,1} = 0 \) for \( n \in \mathbb{T}\). Thus, applying Lemma 3.1 gives
\[
e_3 \times V = \sum_{l,n: \beta_n \neq 0} C_{n,l} (-1)^l (U_{n,1-l}) T (\cos \theta_n)^{2l-1} + \sum_{n: \beta_n = 0} C_{n,0} e_3 \times U_{n,0} \quad \text{on } \Gamma_b,
\]  
(5.45)

and
\[
\left( curl \, [V]_{mo} \right)_{T} \quad \text{on } \Gamma_b.
\]
\[
= - \sum_{n: \beta_n \neq 0} C_{n,l} (curl U_{n,l})_T + \sum_{l,n: \beta_n \neq \mathbb{R}} C_{n,l} (curl U_{n,l})_T
\]
\[
= - \sum_{l,n: \beta_n > 0} i(-1)^l k C_{n,l} (U_{n,1-l})_T + \sum_{l,n: \beta_n \notin \mathbb{R}} i(-1)^l \sqrt{|\alpha_n|^2 + |\beta_n|^{2l-2}} k^{2l} C_{n,l} (U_{n,1-l})_T
\]
(5.46)
on \(\Gamma_b\). Inserting the above identity (5.46) into (5.43) and using (5.45), we derive from (5.44) that
\[
0 = \text{Im} \left\{ -ik \sum_{l,n: \beta_n > 0} |C_{n,l}|^2 \| (U_{n,1-l})_T \|_{L^2(\Gamma_b)}^2 (\cos \theta_n)^{2l-1} \right\}
\]
\[
+ \text{Im} \left\{ - \sum_{l,n: \beta_n \notin \mathbb{R}} |C_{n,l}|^2 (ik)^{2l} \| (U_{n,1-l})_T \|_{L^2(\Gamma_b)}^2 \beta_n^{1-2l} (|\alpha_n|^2 + |\beta_n|^{2l})^{1-2l} \right\}
\]
\[
= -k \sum_{l,n: \beta_n > 0} |C_{n,l}|^2 \| (U_{n,1-l})_T \|_{L^2(\Gamma_b)}^2 (\cos \theta_n)^{2l-1},
\]
which, together with the definition of \( \cos \theta_n \) defined in Lemma 3.1, leads to
\[
C_{n,l} = 0 \quad \text{for all } \beta_n > 0, \ l = 1, 2.
\]  
(5.47)

Finally, combining (5.47) and (5.45) we have proved (5.38) for \( V_T \mid_{\Gamma_b} \), and combining (5.47), (5.46) and (5.43) leads to the desired result for \( (curl \, V^+)_{\Gamma} \).
\[\square\]
We proceed with the proof of Theorem 2.1, i.e., to show the existence of a solution \((E, E^+) \in \mathbb{H}\) to the variational problem (3.14) for the incident wave \(E^\text{gen}_{\text{inc}}\).

Assume that \((V, V^+) \in \mathbb{H}\) satisfies \(a((W, W^+), (V, V^+)) = 0\) for all \((W, W^+) \in \mathbb{H}\). Using Lemmas 5.1 and 3.1, it is easy to check that

\[
a((0, E^\text{gen}_{\text{inc}}), (V, V^+)) = -\int_{\Gamma_b} \left\{ \text{curl} \ E^\text{inc}_{\text{gen}} \cdot e_3 \times V + e_3 \times E^\text{inc}_{\text{gen}} \cdot \text{curl} V^+ \right\} \, ds = 0.
\]

This means that each solution to the homogenous adjoint problem (5.37) is orthogonal to the right-hand side of the variational problem (3.14) in the sense of (5.48). According to Theorem 4.1, the Fredholm alternative applied to the variational problem (3.14) yields the existence of the solution \((E, E^+) \in \mathbb{H}\) to problem (3.14) for the incident plane waves \(E^\text{inc}_{\text{gen}}\) defined in (2.6).

The claim (5.38) implies that the solution \(V^+\) takes the form

\[
V^+(x) = \sum_{\beta_n \in \mathbb{R}} C_{n,l}(x) + \sum_{\beta_n = 0} C_{n,l}(x) \in Y, \quad x \in \Omega^+_b,
\]

and particularly, the coefficients of the propagating modes for \(\beta_n > 0\) vanish. By analogous arguments, this assertion even remains valid for the solution \((V, V^+)\) to the homogenous variational problem

\[
a((V, V^+), (W, W^+)) = 0, \forall (W, W^+) \in \mathbb{H}.
\]

In other words, the coefficient \(C_{n,l}\) of the difference of two solutions of (BVP) are zero if \(\beta_n > 0\).

The proof of Theorem 2.1 is thus completed.

6 Appendix

For the reader’s convenience, we prove that the variational formulation (2.7) is uniquely solvable for small wavenumbers \(k > 0\). Since Rayleigh frequencies can be excluded for small wavenumbers, by Remark 3.4 we see that such a unique solvability also applies to our variational formulation (3.14) provided \(k\) is sufficiently small.

Lemma 6.1. There exists a sufficiently small wavenumber \(k_0 > 0\) such that the variational formulation (2.7) admits a unique solution \(E \in X\) for all \(k \in (0, k_0]\).

Proof. To prove Lemma 6.1, we need to replace equation (2.7) on the \(k\)-dependent \(\alpha\)-quasi-periodic space \(H(\text{curl}, \Omega_b)\) by an equivalent variational problem acting on the \((\Lambda_1, \Lambda_2)\)-periodic Sobolev space. Introduce the spaces \(H^1_p(\Omega_b), H_p(\text{curl}, \Omega_b), H^s_{t,p}(\Gamma_b), H^s_{t,p}(\text{Div}, \Gamma_b)\) and \(H^s_{t,p}(\text{Curl}, \Gamma_b)\) in the same way as \(H^1(\Omega_b), H(\text{curl}, \Omega_b), H^s_{t}(\Gamma_b), H^s_{t}(\text{Div}, \Gamma_b)\) and \(H^s_{t}(\text{Curl}, \Gamma_b)\) in Section 2, but with \(\alpha = (0, 0)^\top\). Define the operator \(\nabla_\alpha := \nabla + i(\alpha, 0)^\top\) and, analogously to \(X\) in Section 2, the space

\[
D := \left\{ F : \Omega_b \to \mathbb{C}^3, F \in H_p(\text{curl}, \Omega_b), \nu \times F = 0 \text{ on } \Gamma \right\}.
\]
Let \( \tau_n := (2\pi n_1/\Lambda_1, 2\pi n_2/\Lambda_2)^T = \alpha_n - \alpha \) for \( n = (n_1, n_2)^T \in \mathbb{Z}^2 \). Given

\[
\tilde{F}(x') = \sum_{n \in \mathbb{Z}^2} \tilde{F}_n e^{i\tau_n \cdot x'} \in H^{-1/2}_{t,p}(\text{Div}, \Gamma_b), \tag{6.49}
\]

it follows from the definition of the operator \( \mathcal{R} \) (see (2.8)) that

\[
\mathcal{R}(\tilde{E}) = \mathcal{T}(\tilde{F}) \exp(\alpha \cdot x') \quad \text{for} \quad \tilde{E}(x') = e^{i\alpha \cdot x'} \tilde{F}(x') \in H^{-1/2}_{t,p}(\text{Div}, \Gamma_b),
\]

where the operator \( \mathcal{T} : H^{-1/2}_{t,p}(\text{Div}, \Gamma_b) \to H^{-1/2}_{t,p}(\text{Curl}, \Gamma_b) \) is the Dirichlet-to-Neumann map over the space \( H^{-1/2}_{t,p}(\text{Div}, \Gamma_b) \) defined by

\[
(\mathcal{T} \tilde{F})(x') = -\sum_{n \in \mathbb{Z}^2} \frac{1}{\beta_n} \left[ k^2 \tilde{F}_n - (\alpha_n \cdot \tilde{F}_n) \alpha_n \right] \exp(i\tau_n \cdot x'), \quad n = (n_1, n_2)^T \in \mathbb{Z}^2. \tag{6.50}
\]

Note that \( \mathcal{T} \) is well defined for small wavenumbers \( k \in (0, k_0] \), since \( \beta_n \neq 0 \) if \( k_0 \) is sufficiently small. The spaces \( H^{-1/2}_{t,p}(\Gamma_b), H^{-1/2}_{t,p}(\text{Div}, \Gamma_b) \) and \( H^{-1/2}_{t,p}(\text{Curl}, \Gamma_b) \) will be equipped with the norms analogous to the quasi-periodic ones in Section 2, only with the coefficient \( E_n \) replaced by \( \tilde{F}_n \) given in (6.49) and \( \alpha_n \) replaced by \( \tau_n \).

Set \( F^{in}(x) := \exp(-i\alpha \cdot x') E^{in}(x) \) and

\[
F(x) = \exp(-i\alpha \cdot x') E(x), \quad \psi(x) = \exp(-i\alpha \cdot x') \varphi(x) \in H_p(\text{curl}, \Omega_b),
\]

for \( E, \varphi \in H(\text{curl}, \Omega_b) \). We now consider the variational formulation

\[
a_p(F, \psi) := \int_{\Omega_b} \left[ \nabla_x \times F \cdot \nabla_x \times \psi - k^2 F \cdot \overline{\psi} \right] \, dx - \int_{\Gamma_b} \mathcal{T}(e_3 \times F) \cdot (e_3 \times \overline{\psi}) \, ds
\]

\[
= \int_{\Gamma_b} \left[ (\nabla_x \times F^{in})_T - \mathcal{T}(e_3 \times F^{in}) \right] \cdot (e_3 \times \overline{\psi}) \, ds \tag{6.51}
\]

for all \( \psi \in D \), which is the counterpart of problem (2.7) in the periodic space \( H_p(\text{curl}, \Omega_b) \).

The problem (6.51) can be rewritten as the operator equation \( B(F) = f \) in the dual space \( D' \) of \( D \), where for \( \psi \in D \) the dualities \( \langle B(F), \psi \rangle \) and \( \langle f, \psi \rangle \) between \( D' \) and \( D \) are defined by the the sesquilinear form \( a_p(F, \psi) \) and the right hand of (6.51), respectively. By Lemma 4.2, we have the Hodge-decomposition \( D = D_0 \oplus D_1 \), with

\[
D_1 := \left\{ \nabla_q q : q \in H^1_p(\Omega_b), q = 0 \text{ on } \Gamma \right\},
\]

\[
D_0 := \left\{ F_0 \in D : \int_{\Omega_b} \nabla_q q \cdot F_0 \, dx = 0 \text{ for all } \nabla_q q \in X_1 \right\}.
\]

This allows the decomposition

\[
F = F_0 + \nabla_q q, \quad \psi = G_0 + \nabla_q g, \quad F_0, G_0 \in D_0, \nabla_q q, \nabla_q g \in D_1.
\]

Now, the sesquilinear form \( a_p \) in (6.51) can be rewritten as

\[
a_p(F, \psi) = a_p(F_0, G_0) + a_p(\nabla_q q, G_0) + a_p(\nabla_q q, \nabla_q g) + a_p(F_0, \nabla_q g),
\]

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and the operator $B$ takes the form

$$B = \begin{pmatrix} B_1 & B_3 \\ B_2 & B_4 \end{pmatrix},$$

$B_1 : D_0 \to D_0', \quad \langle B_1(F_0), G_0 \rangle = a_p(F_0, G_0), \quad \forall G_0 \in D_0,$

$B_2 : D_0 \to D_1', \quad \langle B_2(F_0), \nabla_a g \rangle = a_p(F_0, \nabla_a g), \quad \forall \nabla_a g \in D_1,$

$B_3 : D_1 \to D_0', \quad \langle B_3(\nabla_a q), G_0 \rangle = a_p(\nabla_a q, G_0), \quad \forall G_0 \in D_0,$

$B_4 : D_1 \to D_1', \quad \langle B_4(\nabla_a g), \nabla_a g \rangle = a_p(\nabla_a q, \nabla_a g), \quad \forall \nabla_a g \in D_1.$

We first prove that the form $a_p$ is coercive over $D_0$ for a small wavenumber $k$. Using the explicit representation of $T$, one can prove that, for $\tilde{F}$ given in (6.49),

$$\text{Re} \int_{\Gamma_b} T(\bar{F}) \cdot \bar{F} \, ds = \Lambda_1 \Lambda_2 \sum_{n:|\alpha_n|>k} \frac{1}{|\alpha_n|^2 - k^2} \left[ k^2 |\tilde{F}_n|^2 - |\alpha_n \cdot \tilde{F}_n|^2 \right].$$

Hence,

$$-\text{Re} \int_{\Gamma_b} T(\bar{F}) \cdot \bar{F} \, ds \geq -\Lambda_1 \Lambda_2 \sum_{n:|\alpha_n|>k} \frac{1}{|\alpha_n|^2 - k^2} k^2 |\tilde{F}_n|^2$$

$$\geq -C_1 \Lambda_1 \Lambda_2 \sum_{n \in \mathbb{Z}^2} \frac{1}{|r_n|^2 + 1} k^2 |\tilde{F}_n|^2$$

$$= -C_1 \Lambda_1 \Lambda_2 k^2 |\tilde{F}|^2_{H_{-1/2}^{1/2}(\Gamma_b)}$$

$$\geq -C_1 \Lambda_1 \Lambda_2 k^2 |\tilde{F}|^2_{\mathcal{H}_{-1/2}^{1/2}(\text{Div }, \Gamma_b)}.$$  

Applying the previous estimate to the trace $e_3 \times F_0$ for $F_0 \in D_0$ and using the continuity of the trace mapping from $H_p(\text{curl }, \Omega_b)$ to $H_{-1/2}^{1/2}(\text{Div }, \Gamma_b)$, we arrive at

$$-\text{Re} \left\{ \int_{\Gamma_b} T(e_3 \times F_0) \cdot (e_3 \times \bar{F}_0) \, ds \right\} \geq -k^2 C_1 \Lambda_1 \Lambda_2 |e_3 \times F_0|^2_{H_{-1/2}^{1/2}(\text{Div }, \Gamma_b)}$$

$$\geq -k^2 C_2 |F_0|^2_{H(\text{curl }, \Omega_b)}.$$  

Therefore,

$$\text{Re} \ a_p(F_0, F_0) \geq ||\nabla_a \times F_0||^2_{L^2(\Omega_b)^3} - k^2 ||F_0||^2_{L^2(\Omega_b)^3} - k^2 C_2 ||F_0||^2_{H(\text{curl }, \Omega_b)^3}. \quad (6.52)$$

Recalling that the function $E_0 := \exp(i \alpha \cdot x') F_0$ belongs to the space $X_0 \subset X$ which is divergence free, we have the Friedrichs-type estimate (see [22, Corollary 4.8])

$$||E_0||^2_{L^2(\Omega_b)^3} \leq C_3 ||\nabla \times E_0||^2_{L^2(\Omega_b)^3},$$

for some constant $C_3 > 0$ independent of $E_0 \in X_0$, which is equivalent to

$$||F_0||^2_{L^2(\Omega_b)^3} \leq C_3 ||\nabla_a \times F_0||^2_{L^2(\Omega_b)^3}. \quad (6.53)$$
Combining (6.52) and (6.53) leads to the coercivity of the form $a_p$ over $D_1$ for small wavenumbers $k < k_0$. This implies that the operator $B^{-1}$ exists with the bounded norm $\|B^{-1}\|_{D_0 \to D_0} \leq C$ for some constant $C > 0$ independent of $k \in (0, k_0]$.

Next we claim that the form $-a_p$ is also coercive over $D_1$. In fact, the function $(Q(x'), 0)^\top := e_3 \times \nabla_a g|_{\Gamma_b}$ can be expanded into

$$Q(x') = \sum_{n \in \mathbb{Z}^2} (-\alpha_n^{(2)}, \alpha_n^{(1)})^\top Q_n \exp(i\tau_n \cdot x'), \quad Q_n \in \mathbb{C}. \quad (6.54)$$

Thus, using the representation of $T$ given in (6.50), we find

$$-\text{Re} a_p(\nabla_a q, \nabla_a q) = k^2 \|\nabla_a q\|^2_{L^2(\Omega_b)^3} + \sum_{n:|\alpha_n|>k} |\beta_n|^{-1}k^2 \|\alpha_n\|^2 |Q_n|^2 \geq C_0 k^2 \|\nabla_a q\|^2_{H(\text{curl, } \Omega_b)^3}.$$

As a consequence, we have $\|B^{-1}_4\|_{D_1 \to D_1} \leq k^{-2}C_0^{-1}$, where the constant $C_0$ does not depend on $k$.

The operator $B$ can be written as the matrix operator

$$B = \begin{pmatrix} B_1 & 0 \\ B_2 & B_4 \end{pmatrix} + \begin{pmatrix} 0 & B_3 \\ 0 & 0 \end{pmatrix}, \quad \begin{pmatrix} B_1 & 0 \\ B_2 & B_4 \end{pmatrix}^{-1} = \begin{pmatrix} B_1^{-1} & 0 \\ -B_2^{-1}B_1^{-1} & B_4^{-1} \end{pmatrix} =: \mathcal{M}.$$

Thus the operator equation $B(F) = f$ is equivalent to

$$\begin{pmatrix} I & 0 \\ 0 & I \end{pmatrix} + \begin{pmatrix} 0 & B_1^{-1}B_3 \\ 0 & -B_2^{-1}B_2B_1^{-1}B_3 \end{pmatrix} \begin{pmatrix} F_0 \\ \nabla_a g \end{pmatrix} = \mathcal{M}f, \quad (6.55)$$

where $I$ denotes the identity operator. To prove the invertibility of $B$, it suffices to show

$$\|B_3\|_{D_1 \to D_0} + \|B_2\|_{D_0 \to D_0} \leq C_4 k^2, \quad (6.56)$$

with some $C_4 > 0$ independent of $k \in (0, k_0]$. Consider the sesquilinear form corresponding to $B_2$:

$$a_p(F_0, \nabla_a g) = -\int_{\Gamma_b} \mathcal{T}(e_3 \times F_0) \cdot (e_3 \times \nabla_a g) \, ds.$$

Expand the first two components of $e_3 \times F_0$, $e_3 \times \nabla_a g$ into the series in (6.49) and (6.54), respectively. Then, by (6.50) we get

$$|a_p(F_0, \nabla_a g)| = k^2 \sum_{n \in \mathbb{Z}^2} \frac{1}{i\beta_n} \widetilde{F}_n \cdot (-\alpha_n^{(2)}, \alpha_n^{(1)})^\top Q_n \leq C_5 k^2 \left( \sum_{n \in \mathbb{Z}^2} (1 + |\tau_n|^2)^{1/2} |Q_n|^2 \right)^{1/2} \left( \sum_{n \in \mathbb{Z}^2} (1 + |\tau_n|^2)^{-1/2} |\widetilde{F}_n|^2 \right)^{1/2} \leq C_6 k^2 \|e_3 \times \nabla_a g\|_{H^{-1/2}(\text{Div, } \Gamma_b)} \|e_3 \times F_0\|_{H^{-1/2}(\text{Div, } \Gamma_b)}.$$
This combined with the continuity of the trace mapping from $H_p(\text{curl}, \Omega_b)$ to $H^{-1/2}_p(\text{Div}, \Gamma_b)$ leads to the estimate in (6.56) for $B_2$. For the proof concerning $B_3$, we can proceed analogously.

We now conclude that the operator on the left hand side of (6.55) is a small perturbation of the identity for all $k < k_0$ if $k_0$ is sufficiently small. Hence, problem (3.14) always admits a unique solution $E$ of the form $E = \exp(i\alpha \cdot x')F$ with $F = F_0 + \nabla_\alpha q$, $F_0 \in D_0$, $\nabla_\alpha q \in D_1$. □

References


