HOMOGENIZATION FOR NON-SELF-ADJOINT PERIODIC ELLIPTIC OPERATORS ON AN INFINITE CYLINDER*

NIKITA N. SENIK[†]

Abstract. We consider the problem of homogenization for non-self-adjoint second-order elliptic differential operators $\mathcal{A}^{\varepsilon}$ of divergence form on $L_2(\mathbb{R}^{d_1} \times \mathbb{T}^{d_2})$, where d_1 is positive and d_2 is non-negative. The coefficients of the operator $\mathcal{A}^{\varepsilon}$ are periodic in the first variable with period ε and smooth in a certain sense in the second. We show that, as ε gets small, $(\mathcal{A}^{\varepsilon} - \mu)^{-1}$ and $\nabla_{x_2}(\mathcal{A}^{\varepsilon} - \mu)^{-1}$ for an appropriate μ converge in the operator norm to, respectively, $(\mathcal{A}^0 - \mu)^{-1}$ and $\nabla_{x_2}(\mathcal{A}^0 - \mu)^{-1}$, where \mathcal{A}^0 is an operator whose coefficients depend only on x_2 . We also obtain an approximation for $\nabla_{x_1}(\mathcal{A}^{\varepsilon} - \mu)^{-1}$ and find the next term in the approximation for $(\mathcal{A}^{\varepsilon} - \mu)^{-1}$. Estimates for the rates of convergence and the rates of approximation are provided and are sharp with respect to the order.

Key words. homogenization, operator error estimates, periodic differential operators, effective operator, corrector

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1. Introduction. The periodic homogenization problem consists in studying asymptotic behavior of solutions of differential equations with rapidly oscillating coefficients. Let Ξ be a domain in \mathbb{R}^d , and let A be a one-periodic matrix-valued function on \mathbb{R}^d that satisfies, for almost every $x \in \mathbb{R}^d$,

$$|c_A|\xi|^2 \le \langle A(x)\xi,\xi\rangle_{\mathbb{C}^d} \le C_A|\xi|^2, \qquad \xi \in \mathbb{C}^d$$

where $0 < c_A \leq C_A < \infty$. Consider the traditional model of homogenization: given a weakly convergent sequence $\{f_{\varepsilon}\}_{\varepsilon>0}$ in $L_2(\Xi)$, find $\{u_{\varepsilon}\}_{\varepsilon>0}$ in the Sobolev space $H^1(\Xi)$ such that

$$-\operatorname{div} A(\varepsilon^{-1}x)\nabla u_{\varepsilon} + \mu u_{\varepsilon} = f_{\varepsilon} \quad \text{in } \Xi,$$
$$u_{\varepsilon} = 0 \quad \text{on } \partial \Xi$$

(with a fixed $\mu \in \mathbb{C}$). The coefficients of this equation are ε -periodic and so rapidly oscillate if ε is small. The key fact about homogenization is that the solutions of such problems converge, as the period tends to 0, to solutions of problems whose coefficients have no oscillations. Classical arguments (as in [BLP], [BP], or [ZhKO]) show that u_{ε} has a weak limit in $H^1(\Xi)$, which is the solution of the problem

$$-\operatorname{div} A^0 \nabla u_0 + \mu u_0 = f_0 \quad \text{in } \Xi,$$
$$u_0 = 0 \quad \text{on } \partial \Xi.$$

where A^0 is a constant matrix and f_0 is the weak limit of f_{ε} . In applications, this means that we can approximate highly heterogeneous media by a homogeneous one.

In the case that Ξ is bounded, the embedding $\mathring{H}^1(\Xi) \to L_2(\Xi)$ is compact. Hence, passing to a subsequence if necessary, one finds that u_{ε} converges, in fact,

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[†]St. Petersburg State University, St. Petersburg 199034, Russia (N.N.Senik@gmail.com).

in the $L_2(\Xi)$ -norm. An elementary argument [AlC], again using the compactness of the embedding, then yields the norm-resolvent convergence of the corresponding operators. (Notice, however, that it says nothing about the rate of convergence.) This result has several important implications. It tells us, in particular, that the spectrum of the original operator converges to the spectrum of the limit operator in the sense of Hausdorff.

Matters are different if Ξ is not bounded, or, more precisely, quasi-bounded (say, an infinite strip on the plane with constant width or the entire plane). In that case, the embedding $\mathring{H}^1(\Xi) \to L_2(\Xi)$ fails to be compact [AF], and therefore one needs a different technique in order to obtain the norm-resolvent convergence. Such a technique was suggested by Birman and Suslina [BSu1] (see also [BSu2]), and it settled the case of $\Xi = \mathbb{R}^d$ for a broad class of purely periodic elliptic problems. What is more, they established a sharp-order bound on the rate of convergence. Many related results appeared in subsequent years; see, for example, [Gr1], [Gr2], [Zh], [ZhP], [Bo], [KLS], and references therein. In the recent paper [ChC], a result of this kind was proved for some of the high-contrast problems.

The present paper is motivated by the study of homogenization problems on unbounded domains when the coefficients of the corresponding operators are periodic only in certain directions. These arise naturally in many applications—for instance, in the theory of waveguides and in elasticity—and were investigated in [S-HT], [OShY], [Su1], [BCSu], and [Se1].

Let Ξ be the cylinder $\mathbb{R}^{d_1} \times \mathbb{T}^{d_2}$ (\mathbb{T} is the flat torus). In [Su1], Suslina studied the homogenization problem for elliptic self-adjoint operators $\mathcal{A}^{\varepsilon}$ on Ξ of the form

$$\mathcal{A}^{\varepsilon} = -\operatorname{div}_{x_1} A_{11}(\varepsilon^{-1}x_1, x_2) \nabla_{x_1} - \operatorname{div}_{x_2} A_{22}(\varepsilon^{-1}x_1, x_2) \nabla_{x_2}$$

Here, A_{11} and A_{22} are periodic in the first variable and Lipschitz in the second. She proved that $\mathcal{A}^{\varepsilon}$ converges in the norm-resolvent sense to an operator \mathcal{A}^{0} , whose coefficients depend only on the nonperiodic variable x_2 , and furthermore

$$\|(\mathcal{A}^{\varepsilon}+1)^{-1}-(\mathcal{A}^{0}+1)^{-1}\|_{\mathbf{B}(L_{2}(\Xi))}\leq C\varepsilon.$$

Such problems were further analyzed in [Se1], where we extended that result to selfadjoint operators with lower-order terms and also obtained an approximation for the resolvent in $\mathbf{B}(L_2(\Xi), H^1(\Xi))$. (Strictly speaking, the work [Se1] deals with the case $d_1 = d_2 = 1$, although it is possible to use the techniques of that article to treat the other cases.) But operators with nondiagonal terms were left beyond the scope of these papers, and it was our original intention to fill this gap.

In this article, we will be concerned with an elliptic non-self-adjoint operator $\mathcal{A}^{\varepsilon}$ on Ξ of the form

$$\mathcal{A}^{\varepsilon} = -\operatorname{div} A(\varepsilon^{-1}x_1, x_2)\nabla + \operatorname{div} a_1(\varepsilon^{-1}x_1, x_2) + a_2^*(\varepsilon^{-1}x_1, x_2)\nabla + q(\varepsilon^{-1}x_1, x_2),$$

where A, a_1 , a_2 , and q are periodic in the first variable with respect to a lattice in \mathbb{R}^{d_1} and have weak derivatives with respect to the second variable. We further assume that the coefficients together with the derivatives belong to certain spaces of Sobolev multipliers. We find approximations for $(\mathcal{A}^{\varepsilon} - \mu)^{-1}$ and $\nabla (\mathcal{A}^{\varepsilon} - \mu)^{-1}$ in the operator norm and prove the following estimates:

(1.1)
$$\| (\mathcal{A}^{\varepsilon} - \mu)^{-1} - (\mathcal{A}^{0} - \mu)^{-1} \|_{\mathbf{B}(L_{2}(\Xi))} \leq C\varepsilon,$$

(1.2)
$$\|\nabla_{x_2}(\mathcal{A}^{\varepsilon} - \mu)^{-1} - \nabla_{x_2}(\mathcal{A}^0 - \mu)^{-1}\|_{\mathbf{B}(L_2(\Xi))^{d_2}} \le C\varepsilon$$

and

(1.3)
$$\|\nabla_{x_1}(\mathcal{A}^{\varepsilon} - \mu)^{-1} - \nabla_{x_1}(\mathcal{A}^0 - \mu)^{-1} - \varepsilon \nabla_{x_1} \mathcal{K}^{\varepsilon}_{\mu} \|_{\mathbf{B}(L_2(\Xi))^{d_1}} \le C\varepsilon,$$

(1.4)
$$\| (\mathcal{A}^{\varepsilon} - \mu)^{-1} - (\mathcal{A}^{0} - \mu)^{-1} - \varepsilon \mathcal{C}^{\varepsilon}_{\mu} \|_{\mathbf{B}(L_{2}(\Xi))} \leq C \varepsilon^{2}$$

(see the statements of Theorems 3.1–3.3 in section 3). Here, \mathcal{A}^0 is the effective operator and $\mathcal{K}^{\varepsilon}_{\mu}$ and $\mathcal{C}^{\varepsilon}_{\mu}$ are correctors. The effective operator has a form similar to that of $\mathcal{A}^{\varepsilon}$, with coefficients depending only on the nonperiodic variable x_2 , while the correctors involve rapidly oscillating functions. The estimates are sharp with respect to the order and the constants on the right may be expressed explicitly in terms of the problem parameters. Some of our results were announced in [Se2].

The estimates (1.2) and (1.4) have no analogue in [Su1] and [Se1]. The former is new and the latter has appeared for the first time in [BSu3] for certain self-adjoint operators on the entire space. A more recent development [Su4] has extended that result to self-adjoint operators with lower-order terms. We also mention here the paper [P], where an estimate similar to (1.4) was obtained for operators on \mathbb{R}^d with smooth coefficients. However, all those results apply only to purely periodic operators.

On the other hand, we may regard the problem we address here as a special case of more general locally periodic homogenization problems (where the coefficients may depend on both x and $\varepsilon^{-1}x$). From this point of view, estimates of type (1.1) and (1.3) are known; see, for instance, [PT], where symmetric operators with no lower-order terms were treated. In contrast, the estimate (1.2) is a feature of our problem. As for (1.4), the arguments provided here can be used to prove a similar result for locally periodic operators as well. Our results on locally periodic operators will appear elsewhere [Se4]. For a detailed treatment, see [Se3].

The operator-theoretic method of Birman and Suslina deals only with purely periodic operators and cannot be extended to locally periodic ones. Nevertheless, the abstract results they obtained may be adapted, by ad hoc means, to get the approximations for operators with A having block-diagonal structure, as shown in [Su1] and [Se1]. However, operators with more general A do not fit into this framework. So, if we are to handle these cases, we must develop a different approach.

Our program is as follows. We first reduce the problem to a problem on a fundamental domain for the lattice. This is done by applying the scaling transformation and the Gelfand transform, both with respect to the periodic variable. The Gelfand transform is a standard tool in the spectral approach to homogenization (see, for example, [BLP, Chapter 4], [ZhKO, Chapter 2, section 6], [Sev], [BSu1], and [COV]). But we would like to emphasize that our method is not spectral in nature, and the use of the Floquet–Bloch theory is a mere convenience. The crucial step is to obtain suitable versions of the resolvent identity (see (4.26) and (4.30) below) that enable us to verify the desired inequalities by rather elementary means.

We note that the torus \mathbb{T}^{d_2} can be replaced by any flat manifold without boundary $(\mathbb{R}^{d_2}, \text{ for instance})$. We hope that the techniques presented in this article will prove useful in studying homogenization problems on domains of type $\mathbb{R}^{d_1} \times (0, 1)^{d_2}$ with Dirichlet or Neumann boundary conditions as well.

The paper is organized as follows. In section 2, we give the necessary background information, introduce the operator $\mathcal{A}^{\varepsilon}$, as well as the effective operator and the correctors, formulate the problem under consideration, and provide an example of $\mathcal{A}^{\varepsilon}$. Section 3 contains presentation of the main results. In section 4, we deal with the problem on the fundamental domain and prove the results.

2. Basic definitions and problem formulation. We begin with some notation.

2.1. Preliminaries. The symbol $\|\cdot\|_U$ denotes the norm on a normed space U. Let U and V be Banach spaces. We use the notation $\mathbf{B}(U, V)$ to denote the Banach space of bounded linear operators from U to V. When U = V, the space $\mathbf{B}(U) =$ $\mathbf{B}(U,U)$ becomes a Banach algebra with identity \mathcal{I} . The inner product on a pre-Hilbert space U is denoted by $(\cdot, \cdot)_U$. In the finite-dimensional case $U = \mathbb{C}^n$, the norm and the inner product are denoted by $|\cdot|$ and $\langle \cdot, \cdot \rangle$, respectively. We shall identify the spaces $\mathbf{B}(\mathbb{C}^n, \mathbb{C}^m)$ with $\mathbb{C}^{m \times n}$.

Let Σ be a domain in \mathbb{R}^d and U a Banach space. Then $L_p(\Sigma; U)$, with $1 \leq p \leq \infty$, is the Banach space of strongly measurable functions $u: \Sigma \to U$ satisfying

$$||u||_{L_p(\Sigma;U)} = \left(\int_{\Sigma} ||u(x)||_U^p dx\right)^{1/p} < \infty$$

if $p < \infty$ and

$$|u||_{L_{\infty}(\Sigma;U)} = \underset{x \in \Sigma}{\operatorname{ess\,sup}} ||u(x)||_{U} < \infty$$

if $p = \infty$. In case $U = \mathbb{C}^n$, we shall write $\|\cdot\|_{p,\Sigma}$ for the norm on $L_p(\Sigma)^n = L_p(\Sigma; U)$ and $(\cdot, \cdot)_{2,\Sigma}$ for the inner product on $L_2(\Sigma)^n = L_2(\Sigma; U)$. We denote by $W_p^m(\Sigma)$, with $m \in \mathbb{N}$ and $1 \leq p \leq \infty$, the Banach space of those measurable functions $u: \Sigma \to \mathbb{C}$ that possess all weak derivatives up to and including order m and such that

$$||u||_{m,p,\Sigma} = ||u||_{W_p^m(\Sigma)} = \left(\sum_{|\alpha| \le m} ||D^{\alpha}u||_{L_p(\Sigma)}^p\right)^{1/p} < \infty$$

if $p < \infty$ and

$$\|u\|_{m,\infty,\Sigma} = \|u\|_{W^m_{\infty}(\Sigma)} = \max_{|\alpha| \le m} \|D^{\alpha}u\|_{L_{\infty}(\Sigma)} < \infty$$

if $p = \infty$. Here we use the notation $D = -i\nabla$. The Hilbert space $W_2^m(\Sigma)$ is denoted by $H^m(\Sigma)$, and $H^m(\Sigma)^*$ is its dual space under the pairing $(\cdot, \cdot)_{2,\Sigma}$. If Σ is not open, then $W_n^m(\Sigma)$ will be understood to mean the Sobolev space on the interior of Σ .

Multipliers between Sobolev spaces are (generalized) functions such that the corresponding multiplication operators are bounded. Here we shall be brief; a thorough treatment of Sobolev multipliers may be found in [MSh]. Let Σ be a Lipschitz domain in \mathbb{R}^d , and let m and n be nonnegative integers satisfying $m \ge n$. Then γ is a Sobolev multiplier between $H^m(\Sigma)$ and $H^n(\Sigma)$ (written $\gamma \in \mathbf{M}(H^m(\Sigma), H^n(\Sigma))$) provided that the operator of multiplication $\gamma: H^m(\Sigma) \to H^n(\Sigma)$ is continuous. The space $\mathbf{M}(H^m(\Sigma), H^n(\Sigma)^*)$ of Sobolev multipliers between $H^m(\Sigma)$ and the dual of $H^n(\Sigma)$ is defined in the same way. Notice that an element of $\mathbf{M}(H^m(\Sigma), H^n(\Sigma)^*)$ is a complex-valued distribution.

We shall normally write $\|\gamma\|_{\mathbf{M}}$ for the norm of a Sobolev multiplier γ . This should lead to no confusion, since, once we discover that $\gamma \in \mathbf{M}(H^m(\Sigma), H^n(\Sigma))$ (or $\gamma \in \mathbf{M}(H^m(\Sigma), H^n(\Sigma)^*)$), we fix the spaces $H^m(\Sigma)$ and $H^n(\Sigma)$ (or the spaces $H^m(\Sigma)$ and $H^n(\Sigma)^*$).

Given a positive δ , the scaling transformation S^{δ} is defined to be the map that assigns to each measurable function u on Σ the measurable function v on $\delta^{-1}\Sigma$ given by $v(y) = \delta^{d/2} u(\delta y)$. Then S^{δ} is an isomorphism of $H^m(\Sigma)$ onto $H^m(\delta^{-1}\Sigma)$, with

$$\|\mathcal{S}^{\delta}\|_{\mathbf{B}(H^m(\Sigma), H^m(\delta^{-1}\Sigma))} \le \max\{1, \delta^m\},\$$

and an isometry provided that m = 0. By duality S^{δ} extends to $H^m(\Sigma)^*$, so that $S^{\delta} \colon H^m(\Sigma)^* \to H^m(\delta^{-1}\Sigma)^*$ is also an isomorphism and

$$\|\mathcal{S}^{\delta}\|_{\mathbf{B}(H^m(\Sigma)^*, H^m(\delta^{-1}\Sigma)^*)} \le \max\{1, \delta^{-m}\}.$$

Let $\{\lambda_m\}_{m \in [d]}$ be a basis of \mathbb{R}^d . Here, [d] denotes the set of integers $\{1, 2, \ldots, d\}$. Then the basis generates the lattice

$$\Lambda = \left\{ \lambda \in \mathbb{R}^d \colon \lambda = \sum_{m \in [d]} n_m \lambda_m, n_m \in \mathbb{Z} \right\}$$

with the basic cell

$$\Omega = \Big\{ x \in \mathbb{R}^d \colon x = \sum_{m \in [d]} x_m \lambda_m, x_m \in [-2^{-1}, 2^{-1}) \Big\}.$$

The dual lattice Λ^* is generated by the basis $\{\lambda_m^*\}_{m \in [d]}$ that is defined by the equations $\langle \lambda_m, \lambda_n^* \rangle = 2\pi \delta_{mn}$. We denote the Brillouin zone by Ω^* :

$$\Omega^* = \left\{ k \in \mathbb{R}^d \colon |k| < |k - \lambda^*|, \lambda^* \in \Lambda^* \setminus \{0\} \right\}.$$

Notice that the closure of Ω^* is a convex polyhedron containing the ball of radius $r_{\Lambda} = 2^{-1} \min_{\lambda^* \in \Lambda^* \setminus \{0\}} |\lambda^*|$ centered at the origin.

Lattices are intimately related to Fourier series. If u is any function in $L_2(\Omega)$, then there is a unique representation

$$u(x) = |\Omega|^{-1/2} \sum_{\lambda^* \in \Lambda^*} \hat{u}_{\lambda^*} e^{-i\langle x, \lambda^* \rangle},$$

where the series converges in $L_2(\Omega)$. The corresponding mapping $u \mapsto \{\hat{u}_{\lambda^*}\}_{\lambda^* \in \Lambda^*}$ is an isometric isomorphism of $L_2(\Omega)$ onto $l_2(\Lambda^*)$.

Let $\widetilde{W}_{p}^{m}(\Omega)$ denote the subspace of $W_{p}^{m}(\Omega)$ consisting of functions whose periodic extensions are in $W_{p,\text{loc}}^{m}(\mathbb{R}^{d})$. The symbol $\widetilde{W}_{p,0}^{m}(\Omega)$ stands for the space of functions in $\widetilde{W}_{p}^{m}(\Omega)$ with zero mean value. We shall write $\widetilde{H}^{m}(\Omega)$ and $\widetilde{H}_{0}^{m}(\Omega)$ for $\widetilde{W}_{2}^{m}(\Omega)$ and $\widetilde{W}_{2,0}^{m}(\Omega)$. Observe that, for each $k \in \mathbb{R}^{d}$ and any $u \in \widetilde{H}^{1}(\Omega)$, we have

$$\|(D+k)u\|_{2,\Omega}^{2} = \sum_{\lambda^{*} \in \Lambda^{*}} |\lambda^{*} + k|^{2} |\hat{u}_{\lambda^{*}}|^{2},$$

which yields a variant of Poincaré's inequality:

(2.1)
$$||u - u_{\Omega}||_{2,\Omega} \le C_{\Omega} ||(D + k)u||_{2,\Omega}$$

for all $k \in \Omega^*$. Here, $u_{\Omega} = |\Omega|^{-1} \int_{\Omega} u(x) dx$ and $C_{\Omega} = r_{\Lambda}^{-1}$.

Another operator that is closely related to lattices is the Gelfand transform $\mathcal{G}: L_2(\mathbb{R}^d) \to L_2(\Omega^* \times \Omega)$ given by

$$(\mathcal{G}u)(k,x) = |\Omega^*|^{-1/2} \sum_{\lambda \in \Lambda} u(x+\lambda) e^{-i\langle x+\lambda,k \rangle}$$

the series converging in $L_2(\Omega^* \times \Omega)$. It is well known that \mathcal{G} is an isometric isomorphism of $L_2(\mathbb{R}^d)$ onto $L_2(\Omega^* \times \Omega)$ and an isomorphism of $H^1(\mathbb{R}^d)$ onto $L_2(\Omega^*; \tilde{H}^1(\Omega))$. By duality, the Gelfand transform extends to $H^1(\mathbb{R}^d)^*$, so $\mathcal{G}: H^1(\mathbb{R}^d)^* \to L_2(\Omega^*; \tilde{H}^1(\Omega))^*$.

2.2. Problem formulation. We fix a positive integer d_1 and a nonnegative integer d_2 ; the first will be the number of the periodic directions, and the second will be the number of the nonperiodic directions. We suppose for specificity that d_2 is positive; the case $d_2 = 0$ is similar, with obvious changes. Let $d = d_1 + d_2$. Set $\Xi = \mathbb{R}^{d_1} \times \mathbb{T}_L^{d_2}$, L > 0, where $\mathbb{T}_L^{d_2}$ stands for the d_2 -dimensional flat torus $\mathbb{R}^{d_2}/(L\mathbb{Z}^{d_2})$, that is, a cube in \mathbb{R}^{d_2} with opposite sides identified. Now for each $x \in \Xi$, we have $x = (x_1, x_2)$, where $x_1 \in \mathbb{R}^{d_1}$ and $x_2 \in \mathbb{T}_L^{d_2}$. The *m*th coordinate of x_1 and the *n*th coordinate of x_2 are denoted by $x_{1,m}$ and $x_{2,n}$, respectively. The symbols D_1 and $D_{1,m}$ will stand for D_{x_1} and $D_{x_{1,m}}$, and similarly $D_2 = D_{x_2}$ and $D_{2,n} = D_{x_{2,n}}$. Let Λ be a lattice in \mathbb{R}^{d_1} acting on Ξ . If we denote a basic cell of Λ by Ω_1 and

Let Λ be a lattice in \mathbb{R}^{a_1} acting on Ξ . If we denote a basic cell of Λ by Ω_1 and the torus $\mathbb{T}_L^{d_2}$ by Ω_2 , then $\Omega = \Omega_1 \times \Omega_2$ is a fundamental domain for Λ , and $\{\Omega_\lambda\}_{\lambda \in \Lambda}$, where $\Omega_\lambda = \lambda + \Omega$, is a tiling of Ξ .

We now introduce a class of allowed coefficients. Let U and V be complex Sobolev spaces over the interior of Ω or subspaces of such spaces. We define $\mathbf{S}(U, V)$ to be the set of all complex-valued generalized functions $\gamma \in C_0^{\infty}(\Xi)^*$ such that (1) γ is periodic with respect to Λ , (2) $\gamma \in \mathbf{M}(U, V)$, and (3) $D_2\gamma \in \mathbf{M}(U, V)^{d_2}$. We shall write $\mathbf{S}(U)$ in place of $\mathbf{S}(U, U)$.

Let A be a matrix-valued function in $\mathbf{S}(L_2(\Omega))^{d \times d}$ with Re A uniformly positive definite, a_1 and a_2 be vector-valued functions in $\mathbf{S}(H^1(\Omega), L_2(\Omega))^d$, and q be a complex-valued distribution in $\mathbf{S}(H^1(\Omega), H^1(\Omega)^*)$. Assume also that

(2.2)
$$\|a_1\|_{\mathbf{M}} + \|a_2\|_{\mathbf{M}} + \|q\|_{\mathbf{M}} < \|(\operatorname{Re} A)^{-1}\|_{\mathbf{M}}^{-1}$$

This last requirement is not as restrictive as it might seem to be. It will turn out that the hypothesis (2.2) is, in a sense, a weaker property than the relative $-\Delta$ -form-boundedness of $a_1^*D + D^*a_2 + q$ with relative bound zero, so that (2.2) is satisfied in most cases—see Remark 2.2 below.

Remark 2.1. Our intention is to replace the scale of Lebesgue spaces with that of multiplier spaces, which prove to be perfectly suited to the problem in question. In particular, for this reason we use $\mathbf{M}(L_2(\Omega))$ for the space $L_{\infty}(\Omega)$.

Let $\mathcal{D}_1 = \begin{pmatrix} D_1 \\ 0 \end{pmatrix}$ and $\mathcal{D}_2 = \begin{pmatrix} 0 \\ D_2 \end{pmatrix}$. We also write \mathcal{D} for D. Given $\varepsilon \in \mathscr{C} = (0, 1]$, we introduce the notation $\gamma^{\varepsilon} = (\mathcal{S}^{1/\varepsilon} \otimes \mathcal{I})\gamma(\mathcal{S}^{\varepsilon} \otimes \mathcal{I})$ for any Sobolev multiplier γ (if γ is a function, then $\gamma^{\varepsilon}(x) = \gamma(\varepsilon^{-1}x_1, x_2)$ for $x \in \Xi$) and define the form a^{ε} on $H^1(\Xi)$ by

(2.3)
$$\boldsymbol{a}^{\varepsilon}[\boldsymbol{u}] = (A^{\varepsilon}\mathcal{D}\boldsymbol{u}, \mathcal{D}\boldsymbol{u})_{2,\Xi} + (\mathcal{D}\boldsymbol{u}, a_1^{\varepsilon}\boldsymbol{u})_{2,\Xi} + (a_2^{\varepsilon}\boldsymbol{u}, \mathcal{D}\boldsymbol{u})_{2,\Xi} + (q^{\varepsilon}\boldsymbol{u}, \boldsymbol{u})_{2,\Xi}.$$

Notice that $\gamma \mapsto \gamma^{\varepsilon}$ is a bounded map of $\mathbf{M}(H^m(\Omega), L_2(\Omega))$ onto $\mathbf{M}(H^m(\Omega^{\varepsilon}), L_2(\Omega^{\varepsilon}))$ and of $\mathbf{M}(H^m(\Omega), H^1(\Omega)^*)$ onto $\mathbf{M}(H^m(\Omega^{\varepsilon}), H^1(\Omega^{\varepsilon})^*)$, with norms not exceeding 1. Here $\Omega^{\varepsilon} = \varepsilon \Omega_1 \times \Omega_2$. Then, since $A^{\varepsilon}, a_n^{\varepsilon}, n \in [2]$, and q^{ε} are periodic with respect to $\varepsilon \Lambda$ and since $\{(\Omega_{\lambda})^{\varepsilon}\}_{\lambda \in \Lambda}$ is a tiling of Ξ , we see that $A^{\varepsilon} \in \mathbf{M}(L_2(\Xi))^{d \times d}, a_n^{\varepsilon} \in$ $\mathbf{M}(H^1(\Xi), L_2(\Xi))^d$ and $q^{\varepsilon} \in \mathbf{M}(H^1(\Xi), H^1(\Xi)^*)$. Furthermore, the corresponding norms are majorized by the multiplier norms of A, a_n , and q, respectively. Now it is clear that a^{ε} is bounded,

(2.4)
$$|a^{\varepsilon}[u,v]| \le C_{\flat} ||u||_{1,2,\Xi} ||v||_{1,2,\Xi}, \quad u,v \in H^{1}(\Xi),$$

where

$$C_{\flat} = \|A\|_{\mathbf{M}} + \|a_1\|_{\mathbf{M}} + \|a_2\|_{\mathbf{M}} + \|q\|_{\mathbf{M}}.$$

Observe also that

(2.5)
$$\operatorname{Re} \mathfrak{a}^{\varepsilon}[u] \ge c_* \|\mathcal{D}u\|_{2,\Xi}^2 - c_{\natural} \|u\|_{2,\Xi}^2, \qquad u \in H^1(\Xi),$$

where

(2.6)
$$c_* = \|(\operatorname{Re} A)^{-1}\|_{\mathbf{M}}^{-1} - \|a_1\|_{\mathbf{M}} - \|a_2\|_{\mathbf{M}} - \|q\|_{\mathbf{M}},$$

(2.7)
$$c_{\natural} = 2^{-1} (\|a_1\|_{\mathbf{M}} + \|a_2\|_{\mathbf{M}}) + \|q\|_{\mathbf{M}}.$$

Since c_* is positive, it follows that a^{ε} is coercive.

Thus, a^{ε} is strictly *m*-sectorial, with sector

$$\mathcal{S}_1 = \left\{ z \in \mathbb{C} \colon |\operatorname{Im} z| \le c_*^{-1} C_\flat(\operatorname{Re} z + c_* + c_\flat) \right\}.$$

Let $\mathcal{A}^{\varepsilon}_{\mu} \colon H^1(\Xi) \to H^1(\Xi)^*$ be the operator associated with the form $a^{\varepsilon}_{\mu} = a^{\varepsilon} - \mu$. Then $\mathcal{A}^{\varepsilon}_{\mu}$ is an isomorphism whenever $\mu \notin S_1$.

Remark 2.2. The hypothesis (2.2) is needed in order for the form a^{ε} to be coercive. In fact, it can be weakened to allow those $a_n, n \in [2]$, and q that satisfy, for any $u \in H^1(\Omega)$,

(2.8)
$$\|a_n u\|_{2,\Omega}^2 \le c_{a_n} \|Du\|_{2,\Omega}^2 + C_{a_n} \|u\|_{2,\Omega}^2,$$

(2.9)
$$|(qu, u)_{2,\Omega}| \le c_q ||Du||_{2,\Omega}^2 + C_q ||u||_{2,\Omega}^2$$

with

(2.10)
$$c_{a_1}^{1/2} + c_{a_2}^{1/2} + c_q < \|(\operatorname{Re} A)^{-1}\|_{\mathbf{M}}^{-1}.$$

Indeed, since we are interested in estimating operator norms (see Theorems 3.1–3.3) and since S^{δ} is an isomorphism, we may replace $\mathcal{A}^{\varepsilon}$ by $\hat{\mathcal{A}}^{\varepsilon} = S^{\delta} \mathcal{A}^{\varepsilon} (S^{\delta})^{-1}$. (Here, we realize the torus $\Omega_2 = \mathbb{T}_L^{d_2}$ as the cube $[0, L]^{d_2}$ with opposite sides identified, and, in this sense, $\delta^{-1} \mathbb{T}_L^{d_2}$ is well defined and equals $\mathbb{T}_{\delta^{-1}L}^{d_2}$.) It is easy to see that the coefficients of $\hat{\mathcal{A}}^1$ are given by $\hat{\mathcal{A}} = \delta^{-2} S^{\delta} \mathcal{A} S^{1/\delta}$, $\hat{a}_n = \delta^{-1} S^{\delta} a_n S^{1/\delta}$, and $\hat{q} = S^{\delta} q S^{1/\delta}$. Therefore, if we take δ so that $\delta^2 \leq \min\{c_{a_1}C_{a_1}^{-1}, c_{a_2}C_{a_2}^{-1}, c_qC_q^{-1}\}$, then

$$\begin{aligned} \|\hat{a}_1\|_{\mathbf{M}} + \|\hat{a}_2\|_{\mathbf{M}} + \|\hat{q}\|_{\mathbf{M}} &\leq \delta^{-2} \left(c_{a_1}^{1/2} + c_{a_2}^{1/2} + c_q \right) \\ &< \delta^{-2} \| (\operatorname{Re} A)^{-1} \|_{\mathbf{M}}^{-1} = \| (\operatorname{Re} \hat{A})^{-1} \|_{\mathbf{M}}^{-1}; \end{aligned}$$

that is, the hypothesis (2.2) holds for $\hat{\mathcal{A}}^{\varepsilon}$. We note that the class of operators such as $\mathcal{A}^{\varepsilon}$ here is broad enough to cover most cases that arise in applications—see an example below.

We are interested in approximations for $(\mathcal{A}_{\mu}^{\varepsilon})^{-1}$ and $\mathcal{D}(\mathcal{A}_{\mu}^{\varepsilon})^{-1}$ in the operator norm on $L_2(\Xi)$. In order to describe these approximations, we define the effective operator and two different correctors.

2.3. Effective operator. Let N be the weak solution of

(2.11)
$$\mathcal{D}_1^* A(\mathcal{D}_1 N + I) = 0$$

in $L_2(\Omega_2; \widetilde{H}_0^1(\Omega_1))^{1 \times d}$, and let M be the weak solution of

$$\mathcal{D}_1^*(A\mathcal{D}_1M + a_2) = 0$$

in $L_2(\Omega_2; \widetilde{H}_0^1(\Omega_1))$. We know that N and M exist and are unique, since we may rewrite these problems as

$$(2.13) D_1^* A_{11} D_1 u = D_1^* f$$

with an f in $L_2(\Omega)^{d_1}$ and $u \in L_2(\Omega_2; \widetilde{H}_0^1(\Omega_1))$ to be found. Notice in passing that such a u satisfies

(2.14)
$$(A_{11}(\cdot, x_2) D_1 u(\cdot, x_2) - f(\cdot, x_2), D_1 v)_{2,\Omega_1} = 0$$

for almost every $x_2 \in \Omega_2$ and all $v \in \widetilde{H}^1(\Omega_1)$.

We now provide some elementary properties of N and M (cf. [Su2, Proposition 8.2]).

LEMMA 2.3. Let u be the weak solution of the problem (2.13), where the function f is in $\mathbf{M}(H^m(\Omega); L_2(\Omega))^{d_1}$, m a nonnegative integer. Then $D_1 u \in \mathbf{M}(H^m(\Omega_2), L_2(\Omega))^{d_1}$ and

$$||D_1 u||_{\mathbf{M}} \le |\Omega_1|^{1/2} ||(\operatorname{Re} A)^{-1}||_{\mathbf{M}} ||f||_{\mathbf{M}}$$

If, in addition, $D_{2,n}f \in \mathbf{M}(H^m(\Omega); L_2(\Omega))^{d_1}$ for some $n \in [d_2]$, then $D_{2,n}D_1u \in \mathbf{M}(H^m(\Omega_2), L_2(\Omega))^{d_1}$ and

$$\|D_{2,n}D_1u\|_{\mathbf{M}} \le |\Omega_1|^{1/2} \|(\operatorname{Re} A)^{-1}\|_{\mathbf{M}} (\|D_{2,n}A\|_{\mathbf{M}}\|(\operatorname{Re} A)^{-1}\|_{\mathbf{M}}\|f\|_{\mathbf{M}} + \|D_{2,n}f\|_{\mathbf{M}}).$$

Proof. Let $v = u|w|^2$ with $w \in C^m(\Omega_2)$. Then $v \in L_2(\Omega_2; \widetilde{H}^1_0(\Omega_1))$, and we can apply both sides of (2.13) to v, obtaining

$$||(D_1u)w||_{2,\Omega} \le ||(\operatorname{Re} A_{11})^{-1}||_{\mathbf{M}} ||fw||_{2,\Omega}.$$

This proves the first assertion.

Suppose now $D_{2,n}f \in \mathbf{M}(H^m(\Omega); L_2(\Omega))^{d_1}$. We know that $D_{2,n}f \in L_2(\Omega)^{d_1}$, so $D_{2,n}u$ exists and belongs to $L_2(\Omega_2; \widetilde{H}_0^1(\Omega_1))$, which may be verified by using the difference quotient technique of Nirenberg. Therefore, we can write

$$D_1^* A_{11} D_1 D_{2,n} u = D_1^* (D_{2,n} f - (D_{2,n} A_{11}) D_1 u).$$

Applying both sides of the last equality to $v = (D_{2,n}u)|w|^2$ with $w \in C^m(\Omega_2)$ yields

$$\|(D_1D_{2,n}u)w\|_{2,\Omega} \le \|(\operatorname{Re} A_{11})^{-1}\|_{\mathbf{M}} (\|(D_{2,n}A_{11})(D_1u)w\|_{2,\Omega} + \|(D_{2,n}f)w\|_{2,\Omega}),$$

and the second assertion follows.

From the above lemma and the Poincaré inequality (2.1), we conclude that $N \in \mathbf{S}(L_2(\Omega_2), L_2(\Omega))^{1 \times d}$ and $\mathcal{D}_1 N \in \mathbf{S}(L_2(\Omega_2), L_2(\Omega))^{d \times d}$, while $M \in \mathbf{S}(H^1(\Omega_2), L_2(\Omega))$ and $\mathcal{D}_1 M \in \mathbf{S}(H^1(\Omega_2), L_2(\Omega))^d$.

We now turn to the effective coefficients.

(2.15)
$$A^{0} = |\Omega_{1}|^{-1} \int_{\Omega_{1}} A(\mathcal{D}_{1}N + I) \, dy_{1}$$

Then, from the properties of A and N, we have $A^0 \in \mathbf{S}(L_2(\Omega_2))^{d \times d}$. It is a standard fact (see [ZhKO, section 1.6]) that if Re A is positive definite, then

(2.16)
$$\operatorname{Re} A^0 \ge \left(|\Omega_1|^{-1} \int_{\Omega_1} (\operatorname{Re} A)^{-1} dy_1 \right)^{-1}$$

This implies that $\operatorname{Re} A^0$ is also positive definite and furthermore $(\operatorname{Re} A^0)^{-1}$ is in $\mathbf{M}(L_2(\Omega_2))^{d \times d}$ and $\|(\operatorname{Re} A^0)^{-1}\|_{\mathbf{M}} \leq \|(\operatorname{Re} A)^{-1}\|_{\mathbf{M}}$.

Next, we define the functions

(2.17)
$$a_1^0 = |\Omega_1|^{-1} \int_{\Omega_1} (\mathcal{D}_1 N + I)^* a_1 \, dy_1,$$

(2.18)
$$a_2^0 = |\Omega_1|^{-1} \int_{\Omega_1} (A\mathcal{D}_1 M + a_2) \, dy_1.$$

Both of these are in $\mathbf{S}(H^1(\Omega_2), L_2(\Omega_2))^d$, as can be seen from the properties of A, a_1 , a_2 and N, M.

Finally, let q^0 correspond to the form

(2.19)
$$(q^0 u, u)_{2,\Omega_2} = |\Omega_1|^{-1} (qu, u)_{2,\Omega} + |\Omega_1|^{-1} (a_1^* \mathcal{D}_1 M u, u)_{2,\Omega}$$

on $H^1(\Omega_2)$. By the properties of a_1 , q, and M, we obtain $q^0 \in \mathbf{S}(H^1(\Omega_2), H^1(\Omega_2)^*)$. Notice that, in the case when q is a function, we have, as usual,

$$q^{0} = |\Omega_{1}|^{-1} \int_{\Omega_{1}} q \, dy_{1} + |\Omega_{1}|^{-1} \int_{\Omega_{1}} a_{1}^{*} \mathcal{D}_{1} M \, dy_{1}.$$

We are almost ready to define the effective operator. Consider the form a^0 on $H^1(\Xi)$ given by

(2.20)
$$a^{0}[u] = (A^{0}\mathcal{D}u, \mathcal{D}u)_{2,\Xi} + (\mathcal{D}u, a_{1}^{0}u)_{2,\Xi} + (a_{2}^{0}u, \mathcal{D}u)_{2,\Xi} + (q^{0}u, u)_{2,\Xi},$$

where

$$(q^0 u, u)_{2,\Xi} = \int_{\mathbb{R}^{d_1}} (q^0 u(x_1, \cdot), u(x_1, \cdot))_{2,\Omega_2} dx_1.$$

Then a^0 is plainly bounded,

(2.21)
$$|a^0[u,v]| \le C^0_{\flat} ||u||_{1,2,\Xi} ||v||_{1,2,\Xi}, \qquad u,v \in H^1(\Xi),$$

with

$$C_{\flat}^{0} = \|A^{0}\|_{\mathbf{M}} + \|a_{1}^{0}\|_{\mathbf{M}} + \|a_{2}^{0}\|_{\mathbf{M}} + \|q^{0}\|_{\mathbf{M}}$$

In a moment, we shall see that it is coercive.

LEMMA 2.4. Let \check{a}^0 be the form on $H^1(\Xi) \oplus L_2(\Xi; \widetilde{H}^1_0(\Omega_1))$ given by

$$\begin{split} \check{a}^{0}[\check{u}] &= |\Omega_{1}|^{-1} \int_{\Xi} \int_{\Omega_{1}} \left(\left\langle A(y_{1}, x_{2}) \, \check{\mathcal{D}}\check{u}(x, y_{1}), \check{\mathcal{D}}\check{u}(x, y_{1}) \right\rangle \\ &+ \left\langle \check{\mathcal{D}}\check{u}(x, y_{1}), a_{1}(y_{1}, x_{2}) \, \check{u}_{1}(x) \right\rangle \\ &+ \left\langle a_{2}(y_{1}, x_{2}) \, \check{u}_{1}(x), \check{\mathcal{D}}\check{u}(x, y_{1}) \right\rangle \right) dx \, dy_{1} \\ &+ |\Omega_{1}|^{-1} \int_{\mathbb{R}^{d_{1}}} (q\check{u}_{1}(x_{1}, \cdot), \check{u}_{1}(x_{1}, \cdot))_{2,\Omega} dx_{1}, \end{split}$$

where $\check{u} = (\check{u}_1, \check{u}_2)$ and $\check{\mathcal{D}}\check{u}(x, y_1) = \mathcal{D}_x\check{u}_1(x) + \mathcal{D}_{y_1}\check{u}_2(x, y_1)$. Then \check{a}^0 is coercive and

for all $\check{u} \in H^1(\Xi) \oplus L_2(\Xi; \widetilde{H}^1_0(\Omega_1)).$

Proof. While the proof is quite similar to that of (2.5), there is a difference: the variables x_1 and y_1 in the definition of \check{a}^0 are "mixed," so that we cannot treat the lower-order terms as before.

We begin with a first-order term. By Cauchy's inequality, we have

$$\int_{\Xi} \int_{\Omega_1} \left\langle \check{\mathcal{D}}\check{u}(x,y_1), a_1(y_1,x_2) \check{u}_1(x) \right\rangle dx \, dy_1 \bigg| \\ \leq \|\check{\mathcal{D}}\check{u}\|_{2,\Xi \times \Omega_1} \left(\int_{\Xi} \int_{\Omega_1} |a_1(y_1,x_2) \check{u}_1(x)|^2 dx \, dy_1 \right)^{1/2}.$$

Let v denote the mapping $y \mapsto (\int_{\mathbb{R}^{d_1}} |\check{u}_1(x_1, y_2)|^2 dx_1)^{1/2}$. Then $v \in H^1(\Omega)$, with $||v||_{1,2,\Omega} \leq |\Omega_1|^{1/2} ||\check{u}_1||_{1,2,\Xi}$, and

$$\int_{\Xi} \int_{\Omega_1} |a_1(y_1, x_2) \,\check{u}_1(x)|^2 dx \, dy_1 = \|a_1 v\|_{2,\Omega}^2.$$

As a result,

$$\begin{aligned} \left| \int_{\Xi} \int_{\Omega_1} \left\langle \check{\mathcal{D}}\check{u}(x, y_1), a_1(y_1, x_2) \,\check{u}_1(x) \right\rangle dx \, dy_1 \right| \\ &\leq |\Omega_1|^{1/2} \|a_1\|_{\mathbf{M}} \|\check{\mathcal{D}}\check{u}\|_{2,\Xi \times \Omega_1} \|\check{u}_1\|_{1,2,\Xi} \\ &\leq \|a_1\|_{\mathbf{M}} \left(\|\check{\mathcal{D}}\check{u}\|_{2,\Xi \times \Omega_1}^{2,\Xi \times \Omega_1} + 2^{-1} |\Omega_1| \|\check{u}_1\|_{2,\Xi}^{2} \right) \end{aligned}$$

We have used here the fact that, by Stokes' theorem,

$$\|\check{\mathcal{D}}\check{u}\|_{2,\Xi\times\Omega_{1}}^{2} = |\Omega_{1}|\|\mathcal{D}_{x}\check{u}_{1}\|_{2,\Xi}^{2} + \|\mathcal{D}_{y_{1}}\check{u}_{2}\|_{2,\Xi\times\Omega_{1}}^{2}.$$

We may likewise prove that

$$\left| \int_{\Xi} \int_{\Omega_1} \left\langle a_2(y_1, x_2) \, \check{u}_1(x), \check{\mathcal{D}} \check{u}(x, y_1) \right\rangle dx \, dy_1 \right|$$

$$\leq \|a_2\|_{\mathbf{M}} \left(\|\check{\mathcal{D}} \check{u}\|_{2, \Xi \times \Omega_1}^2 + 2^{-1} |\Omega_1| \|\check{u}_1\|_{2, \Xi}^2 \right)$$

and

$$\begin{aligned} \left| \int_{\mathbb{R}^{d_1}} (q\check{u}_1(x_1, \cdot), \check{u}_1(x_1, \cdot))_{2,\Omega} dx_1 \right| \\ &= \left| (qv, v)_{2,\Omega} \right| \le \|q\|_{\mathbf{M}} \left(\|\check{\mathcal{D}}\check{u}\|_{2,\Xi \times \Omega_1}^2 + |\Omega_1| \|\check{u}_1\|_{2,\Xi}^2 \right). \end{aligned}$$

Combining these inequalities with

$$\operatorname{Re} \int_{\Xi} \int_{\Omega_1} \left\langle A(y_1, x_2) \, \check{\mathcal{D}} \check{u}(x, y_1), \check{\mathcal{D}} \check{u}(x, y_1) \right\rangle dx \, dy_1$$

$$\geq \| (\operatorname{Re} A)^{-1} \|_{\mathbf{M}}^{-1} \| \check{\mathcal{D}} \check{u} \|_{2, \Xi \times \Omega_1}^2,$$

which is obvious, gives (2.22).

Remark 2.5. The form \check{a}^0 is associated with the two-scale homogenized system, first proposed by Allaire [Al] in the context of two-scale convergence. See also [LNW] for a self-contained approach to this matter.

Now we wish to relate the form \check{a}^0 to a^0 . Fix a $u \in H^1(\Xi)$. Let $\check{u}_1(x) = u(x)$ and $\check{u}_2(x, y_1) = N(y_1, x_2) \mathcal{D}_x u(x) + M(y_1, x_2) u(x)$. We claim that $\check{u} = (\check{u}_1, \check{u}_2)$ belongs to $H^1(\Xi) \oplus L_2(\Xi; \widetilde{H}_0^1(\Omega_1))$. Indeed, $\check{u}_2(x, y_1)$ has a derivative with respect to y_1 , and,

by reasoning explained in the proof of Lemma 2.4, it lies in $L_2(\Xi \times \Omega_1)^{d_1}$ (notice here that $\mathcal{D}_1 N$ and $\mathcal{D}_1 M$ are multipliers). Applying identities for N and M in the form (2.14), we find that $\check{a}^0[\check{u},\check{v}] = 0$ for all $\check{v} \in \{0\} \oplus L_2(\Xi; \widetilde{H}_0^1(\Omega_1))$, and therefore $\check{a}^0[\check{u}] = \check{a}^0[\check{u}, u \oplus 0]$. Now it follows from the definitions of the effective coefficients that

$$\check{a}^0[\check{u}] = a^0[u]$$

for every $u \in H^1(\Xi)$, which is the desired relation.

Since \check{a}^0 is coercive, the above identity tells us that so is a^0 , with

Hence, the form a^0 is strictly *m*-sectorial, with sector

$$\mathscr{S}_0 = \left\{ z \in \mathbb{C} \colon |\operatorname{Im} z| \le c_*^{-1} C_\flat^0(\operatorname{Re} z + c_* + c_\flat) \right\}.$$

Corresponding to $a_{\mu}^{0} = a^{0} - \mu$ there is an operator $\mathcal{A}_{\mu}^{0} = \mathcal{A}^{0} - \mu \colon H^{1}(\Xi) \to H^{1}(\Xi)^{*}$, which is an isomorphism provided that $\mu \notin \mathscr{S}_{0}$. For such a μ , $(\mathcal{A}_{\mu}^{0})^{-1}$ maps $L_{2}(\Xi)$ onto $H^{2}(\Xi)$. (This can be shown by using the difference quotient technique of Nirenberg; see the proof of Lemma 4.2 for further details on this matter.) We denote the largest of the sectors \mathscr{S}_{0} and \mathscr{S}_{1} by \mathscr{S} .

2.4. Correctors. We introduce two types of correctors. The first, denoted $\mathcal{K}^{\varepsilon}_{\mu}$, will be needed to obtain the approximation for $\mathcal{D}_1(\mathcal{A}^{\varepsilon}_{\mu})^{-1}$ and is defined as follows. Let $\mathcal{P}^{\varepsilon}$ be the pseudodifferential operator in the x_1 -variable with symbol $\chi_{\varepsilon^{-1}\Omega_1^*}$, where $\chi_{\varepsilon^{-1}\Omega_1^*}$ is the characteristic function of the set $\varepsilon^{-1}\Omega_1^*$, or, to put it differently,

$$\mathcal{P}^{\varepsilon} = (\mathcal{F} \otimes \mathcal{I})^* \chi_{\varepsilon^{-1}\Omega_1^*} (\mathcal{F} \otimes \mathcal{I}).$$

Here \mathcal{F} is the Fourier transform in $L_2(\mathbb{R}^{d_1})$. Then the corrector $\mathcal{K}^{\varepsilon}_{\mu} \colon L_2(\Xi) \to H^1(\Xi)$ for $\mathcal{A}^{\varepsilon}$ is given by

(2.24)
$$\mathcal{K}^{\varepsilon}_{\mu} = (N^{\varepsilon}\mathcal{D} + M^{\varepsilon})(\mathcal{A}^{0}_{\mu})^{-1}\mathcal{P}^{\varepsilon}.$$

We remark that, while $(N^{\varepsilon}\mathcal{D} + M^{\varepsilon})(\mathcal{A}^{0}_{\mu})^{-1}f$, with $f \in L_{2}(\Xi)$, is not generally in $H^{1}(\Xi)$ (not even in $L_{2}(\Xi)$), the function $\mathcal{K}^{\varepsilon}_{\mu}f$ always is, which may be proved by applying the scaling transformation and the Gelfand transform (see (4.18)) and then using the properties of N and M (see Lemma 4.4). What is more, these calculations show that $\mathcal{K}^{\varepsilon}_{\mu}$ is a bounded operator.

The second corrector, denoted C^{ε}_{μ} , will be needed for a more subtle result. If k is a vector in \mathbb{R}^{d_1} and \mathcal{K} is the corresponding element of $\mathbb{R}^{d_1} \oplus \{0\}$, then we define differential expressions

$$\begin{aligned} \mathcal{S}(k;y_1) &= \left((\pounds + \mathcal{D}_2)^* A(y_1, \cdot) + a_1^*(y_1, \cdot) \right) (\pounds + \mathcal{D}_2) + (\pounds + \mathcal{D}_2)^* a_2(y_1, \cdot) + q(y_1, \cdot), \\ \mathcal{T}(k;y_1) &= \left((\pounds + \mathcal{D}_2)^* A(y_1, \cdot) + a_1^*(y_1, \cdot) \right) \mathcal{D}_{y_1} \end{aligned}$$

and families of operators

$$\mathcal{A}^{0}_{\mu}(k) = (\pounds + \mathcal{D}_{2})^{*} A^{0}(\pounds + \mathcal{D}_{2}) + (a^{0}_{1})^{*}(\pounds + \mathcal{D}_{2}) + (\pounds + \mathcal{D}_{2})^{*} a^{0}_{2} + q^{0} - \mu,$$

$$\mathcal{K}_{\mu}(k; y_{1}) = \left(N(y_{1}, \cdot) \left(\pounds + \mathcal{D}_{2}\right) + M(y_{1}, \cdot)\right) (\mathcal{A}^{0}_{\mu}(k))^{-1}.$$

Let $(\mathcal{A}_{\mu}^{\varepsilon})^+$ be the adjoint of $\mathcal{A}_{\mu}^{\varepsilon}$. For the operator $(\mathcal{A}_{\mu}^{\varepsilon})^+$, we construct the effective operator $(\mathcal{A}_{\mu}^{0})^+$ and the corrector $(\mathcal{K}_{\mu}^{\varepsilon})^+$, as well as the families $\mathcal{A}_{\mu}^{0}(k)^+$

and $\mathcal{K}_{\mu}(k; y_1)^+$. (It may be noted in passing that $(\mathcal{A}^0_{\mu})^+$ is the adjoint of \mathcal{A}^0_{μ} .) Finally, let \mathcal{L}_{μ} be the pseudodifferential operator in the x_1 -variable with operator-valued symbol $k \mapsto \mathcal{L}_{\mu}(k) \colon L_2(\Omega_2) \to L_2(\Omega_2)$, where

$$\mathcal{L}_{\mu}(k) = |\Omega_{1}|^{-1} \int_{\Omega_{1}} \left(\mathcal{K}_{\mu}(k;y_{1})^{+} \right)^{*} \left(\mathcal{S}(k;y_{1}) \left(\mathcal{A}_{\mu}^{0}(k) \right)^{-1} + \mathcal{T}(k;y_{1}) \mathcal{K}_{\mu}(k;y_{1}) \right) dy_{1};$$

that is,

$$\mathcal{L}_{\mu} = (\mathcal{F} \otimes \mathcal{I})^* \mathcal{L}_{\mu}(\cdot) (\mathcal{F} \otimes \mathcal{I}).$$

The operator \mathcal{L}^+_{μ} is constructed similarly. The corrector $\mathcal{C}^{\varepsilon}_{\mu} \colon L_2(\Xi) \to L_2(\Xi)$ is then defined by the formula

(2.25)
$$\mathcal{C}^{\varepsilon}_{\mu} = \left(\mathcal{K}^{\varepsilon}_{\mu} - \mathcal{L}_{\mu}\right) + \left(\left(\mathcal{K}^{\varepsilon}_{\mu}\right)^{+} - \mathcal{L}^{+}_{\mu}\right)^{*}.$$

We will see in what follows that $\mathcal{C}^{\varepsilon}_{\mu}$ is continuous.

Remark 2.6. Notice that, since $\mathcal{A}^0_{\mu}(\cdot)$ is the symbol of \mathcal{A}^0_{μ} in the above indicated sense, \mathcal{L}_{μ} can be written as

$$\mathcal{L}_{\mu} = (\mathcal{A}^0_{\mu})^{-1} \mathcal{M} (\mathcal{A}^0_{\mu})^{-1},$$

where $\mathcal{M}: H^2(\Xi) \to H^1(\Xi)^*$ is a third-order differential operator with coefficients depending only on x_2 .

We conclude this section with an example of the operator $\mathcal{A}^{\varepsilon}$.

2.5. An example. Let d > 1 and p > d. From the Ehrling lemma, we know that if $\gamma \in L_p(\Omega)$, then $\gamma \in \mathbf{M}(H^1(\Omega), L_2(\Omega))$ and for all $\epsilon > 0$ there is a $C_{\gamma}(\epsilon) > 0$, depending on d, p, Ω , and $\|\gamma\|_{p,\Omega}$, such that

(2.26)
$$\|\gamma u\|_{2,\Omega}^2 \le \epsilon \|Du\|_{2,\Omega}^2 + C_{\gamma}(\epsilon)\|u\|_{2,\Omega}^2, \quad u \in H^1(\Omega).$$

As an example of a multiplier between $H^1(\Omega)$ and $H^1(\Omega)^*$, let δ_{Σ} be the Dirac distribution on a d-1 dimensional Lipschitz surface Σ in Ω and let σ be a function in $L_{p-1}(\Sigma)$. Again, for each $\epsilon > 0$ there is a $C_{\sigma}(\epsilon)$, depending on d, p, Ω, Σ , and $\|\sigma\|_{p-1,\Sigma}$, such that

(2.27)
$$\left| (\sigma \delta_{\Sigma} u, u)_{2,\Omega} \right| \le \epsilon \|Du\|_{2,\Omega}^2 + C_{\sigma}(\epsilon) \|u\|_{2,\Omega}^2, \qquad u \in H^1(\Omega).$$

Equipped with this information, we consider a periodic operator on $L_2(\Xi)$ of the form

$$\mathcal{H}^{\varepsilon} = (\mathcal{D} - A_1^{\varepsilon})^* g^{\varepsilon} (\mathcal{D} - A_2^{\varepsilon}) + V^{\varepsilon}$$

We may think of $\mathcal{H}^{\varepsilon}$ as a (possibly non-self-adjoint) periodic Schrödinger operator with magnetic and electric potentials that is associated with metric g^{ε} . Suppose that g is a periodic function in $\operatorname{Lip}(\Omega_2; L_{\infty}(\Omega_1))^{d \times d}$ and $\operatorname{Re} g$ is uniformly positive definite. Let A_1 and A_2 be periodic functions in $W_p^1(\Omega_2; L_p(\Omega_1))^d$. Finally, let Σ be a d-1 dimensional periodic Lipschitz surface in Ξ . Then we assume that V is the sum of a periodic function $\hat{V} \in W_{p/2}^1(\Omega_2; L_{p/2}(\Omega_1))$ and a distribution $\sigma \delta_{\Sigma}$ with periodic $\sigma \in W_{p-1}^1(\Sigma \cap \Omega)$. Clearly, $\mathcal{H}^{\varepsilon}$ thus defined can be expressed in the form

$$\mathcal{H}^{\varepsilon} = \mathcal{D}^* A^{\varepsilon} \mathcal{D} + (a_1^{\varepsilon})^* \mathcal{D} + \mathcal{D}^* a_2^{\varepsilon} + q^{\varepsilon},$$

where the coefficients satisfy the properties (2.8)–(2.10) in Remark 2.2. So our result applies to $\mathcal{H}^{\varepsilon}$.

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It is straightforward to construct an analogous example for the case d = 1. Now we take Σ to be a discrete periodic set of points in \mathbb{R} and assume that $g \in L_{\infty}(\Omega)$ with Re g uniformly positive definite, $A_1, A_2 \in L_2(\Omega)$, and $V = \hat{V} + \sigma \delta_{\Sigma}$, where \hat{V} lies in $L_1(\Omega)$ and σ is a periodic function on Σ .

We note that the potential V^{ε} may also involve a singular term $\varepsilon^{-1}W^{\varepsilon}$ with a suitable function W. We refer the reader to [Su3, section 11] for the details.

3. Main results. We now state the principal results of the present paper.

THEOREM 3.1. If $\mu \notin S$, then for any $\varepsilon \in \mathscr{C}$ we have

(3.1)
$$\|(\mathcal{A}^{\varepsilon}_{\mu})^{-1} - (\mathcal{A}^{0}_{\mu})^{-1}\|_{\mathbf{B}(L_{2}(\Xi))} \lesssim \varepsilon,$$

(3.2)
$$\|\mathcal{D}_2(\mathcal{A}^{\varepsilon}_{\mu})^{-1} - \mathcal{D}_2(\mathcal{A}^0_{\mu})^{-1}\|_{\mathbf{B}(L_2(\Xi))^d} \lesssim \varepsilon.$$

The estimates are sharp with respect to the order, and the constants depend only on d, r_{Λ}, μ , and the multiplier norms of the coefficients.

THEOREM 3.2. If $\mu \notin S$, then for any $\varepsilon \in \mathscr{C}$ we have

(3.3)
$$\|\mathcal{D}_1(\mathcal{A}^{\varepsilon}_{\mu})^{-1} - \mathcal{D}_1(\mathcal{A}^0_{\mu})^{-1} - \varepsilon \mathcal{D}_1 \mathcal{K}^{\varepsilon}_{\mu}\|_{\mathbf{B}(L_2(\Xi))^d} \lesssim \varepsilon.$$

The estimate is sharp with respect to the order, and the constant depends only on d, r_{Λ} , μ , and the multiplier norms of the coefficients.

THEOREM 3.3. If $\mu \notin S$, then for any $\varepsilon \in \mathscr{C}$ we have

(3.4)
$$\| (\mathcal{A}^{\varepsilon}_{\mu})^{-1} - (\mathcal{A}^{0}_{\mu})^{-1} - \varepsilon \mathcal{C}^{\varepsilon}_{\mu} \|_{\mathbf{B}(L_{2}(\Xi))} \lesssim \varepsilon^{2}.$$

The estimate is sharp with respect to the order, and the constant depends only on d, r_{Λ} , μ , and the multiplier norms of the coefficients.

Remark 3.4. Although it is possible to write down all the constants explicitly, we do not do so here. In particular, we write $\alpha \leq \beta$ to mean $\alpha \leq C\beta$, where C is a positive constant depending only on d, r_{Λ} , μ , and the multiplier norms of the coefficients.

Remark 3.5. Using the resolvent identity, we can transfer the estimates in Theorems 3.1–3.3 to those $\mu \in \mathcal{S}$ for which $\mathcal{A}^{\varepsilon}_{\mu}$ (at least for each ε in an interval $(0, \varepsilon_{\mu}]$) and \mathcal{A}^{0}_{μ} , when viewed as operators on $L_{2}(\Xi)$, have bounded inverses with norms majorized by constants independent of ε . For instance, the estimates hold if μ is in the resolvent set of the effective operator. In this case, we may choose ε_{μ} to be equal to 1 if $\mu \notin \mathcal{S}_{1}$ and to be so small that, for some $\mu_{0} \notin \mathcal{S}$,

$$\varepsilon_{\mu} < \frac{1}{C_{\mu_0} |\mu - \mu_0| (1 + |\mu - \mu_0| \| (\mathcal{A}^0_{\mu})^{-1} \|_{\mathbf{B}(L_2(\Xi))})}$$

if $\mu \in S_1$. Here C_{μ_0} is the implied constant in the estimate (3.1) with μ_0 in place of μ . The condition on ε_{μ} for $\mu \in S_1$ comes from suitable resolvent identities and the fact that $\mathcal{A}^{\varepsilon}_{\mu_0}$ converges to $\mathcal{A}^{0}_{\mu_0}$ in the norm-resolvent sense.

Remark 3.6. The hypothesis that the coefficients have weak derivatives with respect to the nonperiodic variable is crucial to our analysis and reflects the fact that the roles of the two variables are quite different. Roughly speaking, only the first variable is involved in the homogenization procedure, while the second plays the role of a parameter (see, for example, the definitions of N and M, where this is literally the case). In particular, the hypothesis ensures that N and M belong to $H^1(\Xi)$ and that the preimage of $L_2(\Xi)$ under \mathcal{A}^0_{μ} is $H^2(\Xi)$; as a consequence, the range of $\mathcal{K}^{\varepsilon}_{\mu}$ is contained in $H^1(\Xi)$.

Remark 3.7. While $(\mathcal{A}_{\mu}^{\varepsilon})^{-1}$ and $\mathcal{D}_{2}(\mathcal{A}_{\mu}^{\varepsilon})^{-1}$ have limits, the operator $\mathcal{D}_{1}(\mathcal{A}_{\mu}^{\varepsilon})^{-1}$ can fail to converge, because, though the norm of $\varepsilon \mathcal{D}_{1} \mathcal{K}_{\mu}^{\varepsilon}$ is bounded uniformly in ε , it need not go to zero. However, if, for example, $\mathcal{D}_{1}^{*}A = 0$ and $\mathcal{D}_{1}^{*}a_{2} = 0$ (in the weak sense), then $\mathcal{K}_{\mu}^{\varepsilon} = 0$, and $\mathcal{D}_{1}(\mathcal{A}_{\mu}^{\varepsilon})^{-1}$ is therefore convergent as well. Notice that, in this case, the effective coefficients are obtained by simply taking the mean over Ω_{1} .

Remark 3.8. We may replace $\mathcal{P}^{\varepsilon}$ with another smoothing. For instance, the Steklov averaging operator (see [Zh]) or the scale-splitting operator (see [Gr1]) can be used instead. This follows from the inequalities

(3.5)
$$\| (\mathcal{D}_1 N) w \|_{2,\Xi}^2 \lesssim \| N \otimes \mathcal{D}_1 w \|_{2,\Xi}^2 + \| w \|_{2,\Xi}^2,$$

(3.6)
$$\|(\mathcal{D}_1 M)w\|_{2,\Xi}^2 \lesssim \|M\mathcal{D}_1 w\|_{2,\Xi}^2 + \|w\|_{1,2,\Xi}^2,$$

which hold for any $w \in C_0^{\infty}(\Xi)$ (the proof of the inequalities is parallel to that of Lemma 2.3), and properties of the smoothing operators (cf. similar techniques in [PSu, Lemma 3.5]). The reason why we chose $\mathcal{P}^{\varepsilon}$ is merely one of convenience: as we shall see below, $\mathcal{P}^{\varepsilon}$ takes a rather simple form after passing to the fundamental domain.

Remark 3.9. As already indicated, the operator $\mathcal{P}^{\varepsilon}$ guarantees that the range of $\mathcal{K}^{\varepsilon}_{\mu}$ is contained in $H^1(\Xi)$. This means that, in general, it is not possible to remove $\mathcal{P}^{\varepsilon}$. However, this can be done in certain cases. For example, if $N \in \mathbf{M}(L_2(\Omega))^{1 \times d}$ and $M \in \mathbf{M}(H^1(\Omega), L_2(\Omega))$, then the classical corrector

$$\hat{\mathcal{K}}^{\varepsilon}_{\mu} = (N^{\varepsilon}\mathcal{D} + M^{\varepsilon})(\mathcal{A}^{0}_{\mu})^{-1}$$

as well as the composition $\mathcal{D}_1 \hat{\mathcal{K}}^{\varepsilon}_{\mu}$ are bounded on $L_2(\Xi)$ (by (3.5) and (3.6)), and the estimates (3.3) and (3.4) remain true with $\hat{\mathcal{K}}^{\varepsilon}_{\mu}$ in place of $\mathcal{K}^{\varepsilon}_{\mu}$ and

$$\hat{\mathcal{C}}^{\varepsilon}_{\mu} = \left(\hat{\mathcal{K}}^{\varepsilon}_{\mu} - \mathcal{L}_{\mu}\right) + \left((\hat{\mathcal{K}}^{\varepsilon}_{\mu})^{+} - \mathcal{L}^{+}_{\mu}\right)^{*}$$

in place of $\mathcal{C}^{\varepsilon}_{\mu}$.

4. Problem on the fundamental domain. Our strategy is to reduce $\mathcal{A}^{\varepsilon}_{\mu}$ to an operator on the fundamental domain Ω and then formulate Theorems 3.1–3.3 in terms of this latter operator.

Let $\tau = (k, \varepsilon) \in \mathcal{T} = \Omega_1^* \times \mathcal{E}$. We introduce the notation $\mathcal{D}_1(\tau) = \mathcal{D}_1 + \mathcal{K}$, $\mathcal{D}_2(\tau) = \varepsilon \mathcal{D}_2$ and $\mathcal{D}(\tau) = \mathcal{D}_1(\tau) + \mathcal{D}_2(\tau)$ and set

$$||u||_{1,2,\Omega;\tau} = \left(||\mathcal{D}(\tau)u||_{2,\Omega}^2 + |\tau|^2 ||u||_{2,\Omega}^2\right)^{1/2}$$

and

$$\begin{aligned} \|u\|_{1_{1},2,\Omega;\tau} &= \left(\|\mathcal{D}_{1}(\tau)u\|_{2,\Omega}^{2} + |\tau|^{2}\|u\|_{2,\Omega}^{2}\right)^{1/2},\\ \|u\|_{1_{2},2,\Omega;\tau} &= \left(\|\mathcal{D}_{2}(\tau)u\|_{2,\Omega}^{2} + |\tau|^{2}\|u\|_{2,\Omega}^{2}\right)^{1/2} \end{aligned}$$

for any u for which the right-hand sides make sense.

Now let us define periodic Sobolev spaces over the interior of Ω . Recall that we view Λ as acting on Ξ and $\Omega \subset \mathbb{R}^d$ is a fundamental domain for Λ . The space $\widetilde{W}_p^m(\Omega)$ consists of all functions that are in $W_p^m(\Omega)$ and that have periodic extensions in $W_p^m(K)$ for each compact set $K \subset \Xi$. Let $\widetilde{W}_{p,0}^m(\Omega)$ be the subspace of functions in $\widetilde{W}_p^m(\Omega)$ with zero mean over Ω_1 . As usual, $\widetilde{H}^m(\Omega) = \widetilde{W}_2^m(\Omega)$ and $\widetilde{H}_0^m(\Omega) = \widetilde{W}_{2,0}^m(\Omega)$.

We define the form $a(\tau)$ on $\widetilde{H}^1(\Omega)$ by

(4.1)
$$a(\tau)[u] = (A\mathcal{D}(\tau)u, \mathcal{D}(\tau)u)_{2,\Omega} + \varepsilon(\mathcal{D}(\tau)u, a_1u)_{2,\Omega} + \varepsilon(a_2u, \mathcal{D}(\tau)u)_{2,\Omega} + \varepsilon^2(qu, u)_{2,\Omega}.$$

Note that, when estimating $||a_n u||_{2,\Omega}$, $n \in [2]$, and $|(qu, u)_{2,\Omega}|$, we can replace u by uv with $v(x) = e^{i\langle x_1, k \rangle}$. Hence,

(4.2)
$$\varepsilon \|a_n u\|_{2,\Omega} \le \|a_n\|_{\mathbf{M}} (\varepsilon^2 \|\mathcal{D}_1(\tau) u\|_{2,\Omega}^2 + \|\mathcal{D}_2(\tau) u\|_{2,\Omega}^2 + \varepsilon^2 \|u\|_{2,\Omega}^2)^{1/2},$$

(4.3) $\varepsilon^{2}|(qu, u)_{2,\Omega}| \leq ||q||_{\mathbf{M}} (\varepsilon^{2} ||\mathcal{D}_{1}(\tau)u||_{2,\Omega}^{2} + ||\mathcal{D}_{2}(\tau)u||_{2,\Omega}^{2} + \varepsilon^{2} ||u||_{2,\Omega}^{2});$

in particular, this means that

(4.4)
$$\varepsilon \|a_n u\|_{2,\Omega} \le \|a_n\|_{\mathbf{M}} \left(\|\mathcal{D}(\tau)u\|_{2,\Omega}^2 + \varepsilon^2 \|u\|_{2,\Omega}^2 \right)^{1/2},$$

(4.5)
$$\varepsilon^2 |(qu, u)_{2,\Omega}| \le ||q||_{\mathbf{M}} \left(||\mathcal{D}(\tau)u||_{2,\Omega}^2 + \varepsilon^2 ||u||_{2,\Omega}^2 \right)$$

for all $u \in H^1(\Omega)$. Therefore, the same reasoning as for a^{ε} gives

(4.6)
$$|a(\tau)[u,v]| \le C_{\flat} ||u||_{1,2,\Omega;\tau} ||v||_{1,2,\Omega;\tau}, \quad u,v \in \widetilde{H}^{1}(\Omega),$$

and

(4.7)
$$\operatorname{Re} \boldsymbol{a}(\tau)[u] \ge c_* \|\mathcal{D}(\tau)u\|_{2,\Omega}^2 - \varepsilon^2 c_{\natural} \|u\|_{2,\Omega}^2, \qquad u \in \widetilde{H}^1(\Omega).$$

Define $\mathcal{A}_{\mu}(\tau) = \mathcal{A}(\tau) - \varepsilon^2 \mu$: $\widetilde{H}^1(\Omega) \to \widetilde{H}^1(\Omega)^*$ to be the operator associated with the form $a_{\mu}(\tau) = a(\tau) - \varepsilon^2 \mu$. It follows that $\mathcal{A}_{\mu}(\tau)$ is an isomorphism if $\mu \notin S_1$.

LEMMA 4.1. For any $\mu \notin S$ and $\tau \in \mathcal{T}$ we have

$$\begin{aligned} \|(\mathcal{A}_{\mu}(\tau))^{-1}\|_{\mathbf{B}(L_{2}(\Omega))} \lesssim |\tau|^{-2}, \\ \|\mathcal{D}(\tau)(\mathcal{A}_{\mu}(\tau))^{-1}\|_{\mathbf{B}(L_{2}(\Omega))^{d}} \lesssim |\tau|^{-1}, \\ \|\mathcal{D}(\tau)(\mathcal{A}_{\mu}(\tau))^{-1}\mathcal{D}(\tau)\|_{\mathbf{B}(L_{2}(\Omega))^{d\times d}} \lesssim 1, \\ \|\mathcal{D}(\tau)\mathcal{D}_{2}(\tau)(\mathcal{A}_{\mu}(\tau))^{-1}\|_{\mathbf{B}(L_{2}(\Omega))^{d\times d}} \lesssim 1, \end{aligned}$$

where the constants depend on d, μ , and the multiplier norms of the coefficients.

Proof. We do the case $\mu \in \mathcal{R}$, where

$$\mathscr{R} = \{ z \in \mathbb{C} \colon \operatorname{Re} z < -c_{\natural} \}.$$

The general case then follows by the resolvent identity.

Expanding $u \in \tilde{H}^1(\Omega)$ in a Fourier series

$$u(x) = |\Omega_1|^{-1/2} \sum_{\lambda^* \in \Lambda^*} \hat{u}_{\lambda^*}(x_2) e^{-i\langle x_1, \lambda^* \rangle}$$

we find that, for all $k \in \Omega_1^*$,

(4.8)
$$\|\mathcal{D}_1(\tau)u\|_{2,\Omega}^2 = \sum_{\lambda^* \in \Lambda^*} |\lambda^* + k|^2 \|\hat{u}_{\lambda^*}\|_{2,\Omega_2}^2 \ge |k|^2 \|u\|_{2,\Omega}^2,$$

which means that

$$|\tau|^2 ||u||^2_{2,\Omega} \le ||\mathcal{D}(\tau)u||^2_{2,\Omega} + \varepsilon^2 ||u||^2_{2,\Omega}.$$

Combining this with (4.7) gives the first estimate. The second is immediate from the first and (4.7), and the third follows at once from (4.7). It remains to prove the last.

We shall use the classical technique of difference quotients. To that end, we introduce a little notation. Let $e_{2,m}$, $m \in [d_2]$, be the unit vector along the $x_{2,m}$ -axis. If uis any function on Ω , then we define the difference quotient $D_{2,m}^h$ in the variable $x_{2,m}$ of size $h \in \mathbb{R} \setminus \{0\}$ by setting $D_{2,m}^h u = -ih^{-1}(\mathcal{T}_{2,m}^h u - u)$, where $(\mathcal{T}_{2,m}^h u)(x) = u(x + he_{2,m})$. Observe that

$$(D_{2,m}^h)^* = D_{2,m}^{-h}$$

and

$$D^{h}_{2,m}(uv) = (D^{h}_{2,m}u)\mathcal{T}^{h}_{2,m}v + u(D^{h}_{2,m}v)$$

For $f \in L_2(\Omega)$ fixed, we set $w = (\mathcal{A}_{\mu}(\tau))^{-1} f$. Then

(4.9)
$$\mathbf{a}_{\mu}(\tau)[D_{2,m}^{h}w] = (f, D_{2,m}^{-h}D_{2,m}^{h}w)_{2,\Omega} - \left([D_{2,m}^{h}, \mathcal{A}_{\mu}(\tau)]w, D_{2,m}^{h}w\right)_{2,\Omega}.$$

If we show that

(4.10)
$$\left| \left([D_{2,m}^h, \mathcal{A}_\mu(\tau)] u, v \right)_{2,\Omega} \right| \lesssim \| u \|_{1,2,\Omega;\tau} \| v \|_{1,2,\Omega;\tau}, \qquad u, v \in \widetilde{H}^1(\Omega),$$

where the constant is independent of h, then, by the estimates that we just proved, the right-hand side of (4.9) will not exceed

$$\begin{split} \varepsilon^{-1} \| \mathcal{D}(\tau) D_{2,m}^{h} w \|_{2,\Omega} \| f \|_{2,\Omega} + |\tau|^{-1} C \| D_{2,m}^{h} w \|_{1,2,\Omega;\tau} \| f \|_{2,\Omega} \\ &\leq 2^{-1} c_{*} \| \mathcal{D}(\tau) D_{2,m}^{h} w \|_{2,\Omega}^{2} + \varepsilon^{2} |c_{\natural} + \operatorname{Re} \mu| \| D_{2,m}^{h} w \|_{2,\Omega}^{2} + \varepsilon^{-2} C \| f \|_{2,\Omega}^{2} \end{split}$$

with some constant C. We have used here the fact that, for any h,

$$\|D_{2,m}^{-h}D_{2,m}^{h}w\|_{2,\Omega} \le \|D_{2,m}D_{2,m}^{h}w\|_{2,\Omega}.$$

On the other hand, since $-(c_{\natural} + \operatorname{Re} \mu) > 0$, it follows from (4.7) that

$$\operatorname{Re} a_{\mu}(\tau)[D_{2,m}^{h}w] \ge c_{*} \|\mathcal{D}(\tau)D_{2,m}^{h}w\|_{2,\Omega}^{2} + \varepsilon^{2}|c_{\natural} + \operatorname{Re} \mu| \|D_{2,m}^{h}w\|_{2,\Omega}^{2}.$$

As a result,

$$\|\mathcal{D}(\tau)\varepsilon D_{2,m}^h w\|_{2,\Omega} \lesssim \|f\|_{2,\Omega}$$

uniformly in h. Thus, there exists $\mathcal{D}(\tau)\mathcal{D}_{2,m}(\tau)w$, with

$$\|\mathcal{D}(\tau)\mathcal{D}_{2,m}(\tau)w\|_{2,\Omega} \lesssim \|f\|_{2,\Omega},$$

as desired.

We conclude by proving (4.10). Since

$$([D_{2,m}^h, \mathcal{A}_{\mu}(\tau)]\mathcal{T}_{2,m}^{-h}u, v)_{2,\Omega} = ((D_{2,m}^h A)\mathcal{D}(\tau)u, \mathcal{D}(\tau)v)_{2,\Omega} - \varepsilon (\mathcal{D}(\tau)u, (D_{2,m}^h a_1)v)_{2,\Omega} + \varepsilon ((D_{2,m}^h a_2)u, \mathcal{D}(\tau)v)_{2,\Omega} + \varepsilon^2 ((D_{2,m}^h q)u, v)_{2,\Omega},$$

we see that it suffices to show that each coefficient of this form is still a multiplier with norm bounded by a constant independent of h, because then an argument similar to the one we used for proving (4.6) will lead to (4.10). Obviously,

$$(D_{2,m}^{h}\gamma)(x) = \int_{(0,1)} (\mathcal{T}_{2,m}^{th} D_{2,m}\gamma)(x) \, dt$$

for every function γ that has a derivative $D_{2,m}\gamma \in L_2(\Omega)$. A duality argument shows that if $D_{2,m}\gamma \in H^1(\Omega)^*$, then

$$\left((D_{2,m}^{h}\gamma)u, u \right)_{2,\Omega} = \int_{(0,1)} \left((D_{2,m}\gamma)\mathcal{T}_{2,m}^{-th}u, \mathcal{T}_{2,m}^{-th}u \right)_{2,\Omega} dt$$

for any $u \in C^1(\bar{\Omega})$. Hence, $\|D_{2,m}^h A\|_{\mathbf{M}} \leq \|D_{2,m} A\|_{\mathbf{M}}, \|D_{2,m}^h a_n\|_{\mathbf{M}} \leq \|D_{2,m} a_n\|_{\mathbf{M}}, n \in [2]$, and $\|D_{2,m}^h q\|_{\mathbf{M}} \leq \|D_{2,m} q\|_{\mathbf{M}}$. This completes the proof.

Let $\mathcal{A}_{\mu}(\tau)^{+}$: $\widetilde{H}^{1}(\Omega) \to \widetilde{H}^{1}(\Omega)^{*}$ be the formal adjoint of $\mathcal{A}_{\mu}(\tau)$. From (4.6) and (4.7), we see that $\mathcal{A}_{\mu}(\tau)^{+}$ is also an isomorphism whenever $\mu \notin \mathcal{S}_{1}$. Moreover, the conclusion of Lemma 4.1 holds for $\mathcal{A}_{\mu}(\tau)^{+}$. It is easy to note the relationship between $\mathcal{A}_{\mu}(\tau)$ and $\mathcal{A}_{\mu}(\tau)^{+}$. Indeed, a suitable restriction of $\mathcal{A}_{\mu}(\tau)^{+}$ is the adjoint of the restriction of $\mathcal{A}_{\mu}(\tau)$, so that, if $\mu \notin \mathcal{S}_{1}$,

(4.11)
$$((\mathcal{A}_{\mu}(\tau))^{-1}f, f)_{2,\Omega} = (f, (\mathcal{A}_{\mu}(\tau)^{+})^{-1}f)_{2,\Omega}, \quad f \in L_{2}(\Omega).$$

We shall think of $L_2(\Omega)$ as the tensor product $L_2(\Omega_1) \otimes L_2(\Omega_2)$. Recall that \mathcal{G} is the Gelfand transform and $\mathcal{S}^{\varepsilon}$ is the scaling transformation. Clearly, $\mathcal{GS}^{\varepsilon} \otimes \mathcal{I}$ maps $H^1(\Xi)$ onto $L_2(\Omega_1^*; \widetilde{H}^1(\Omega))$ and, for any $u \in H^1(\Xi)$,

$$a^{\varepsilon}_{\mu}[u] = \varepsilon^{-2} \int_{\Omega^*_1} a_{\mu}(\tau) [\tilde{u}(k, \cdot)] dk,$$

where $\tilde{u} = (\mathcal{GS}^{\varepsilon} \otimes \mathcal{I})u$; that is,

(4.12)
$$(\mathcal{GS}^{\varepsilon} \otimes \mathcal{I})(\mathcal{A}^{\varepsilon}_{\mu})^{-1}(\mathcal{GS}^{\varepsilon} \otimes \mathcal{I})^{-1} = \int_{\Omega_{1}^{\varepsilon}}^{\oplus} \varepsilon^{2}(\mathcal{A}_{\mu}(\tau))^{-1} dk.$$

Now, we would like to do the same for the operator \mathcal{A}^0_{μ} . To this end, let $a^0(\tau)$ be the form on $\widetilde{H}^1(\Omega)$ defined by

(4.13)
$$a^{0}(\tau)[u] = (A^{0}\mathcal{D}(\tau)u, \mathcal{D}(\tau)u)_{2,\Omega} + \varepsilon(\mathcal{D}(\tau)u, a_{1}^{0}u)_{2,\Omega} + \varepsilon(a_{2}^{0}u, \mathcal{D}(\tau)u)_{2,\Omega} + \varepsilon^{2}(q^{0}u, u)_{2,\Omega},$$

where

$$(q^0 u, u)_{2,\Omega} = \int_{\Omega_1} (q^0 u(x_1, \cdot), u(x_1, \cdot))_{2,\Omega_2} dx_1.$$

The same arguments used to obtain (2.21) and (2.23) now show that

(4.14)
$$|\boldsymbol{a}^{0}(\tau)[u,v]| \leq C_{\flat}^{0} \|\boldsymbol{u}\|_{1,2,\Omega;\tau} \|v\|_{1,2,\Omega;\tau}, \quad u,v \in \widetilde{H}^{1}(\Omega)$$

and

(4.15)
$$\operatorname{Re} \mathfrak{a}^{0}(\tau)[u] \geq c_{*} \|\mathcal{D}(\tau)u\|_{2,\Omega}^{2} - \varepsilon^{2} c_{\natural} \|u\|_{2,\Omega}^{2}, \qquad u \in \widetilde{H}^{1}(\Omega).$$

Let $\mathcal{A}^{0}_{\mu}(\tau) = \mathcal{A}^{0}(\tau) - \varepsilon^{2} \mu$: $\widetilde{H}^{1}(\Omega) \to \widetilde{H}^{1}(\Omega)^{*}$ be the operator corresponding to $a^{0}_{\mu}(\tau) = a^{0}(\tau) - \varepsilon^{2} \mu$. Then $\mathcal{A}^{0}_{\mu}(\tau)$ is an isomorphism if $\mu \notin \mathscr{S}_{0}$.

LEMMA 4.2. For any $\mu \notin S$ and $\tau \in \mathcal{T}$ we have

$$\begin{aligned} \|(\mathcal{A}^{0}_{\mu}(\tau))^{-1}\|_{\mathbf{B}(L_{2}(\Omega))} \lesssim |\tau|^{-2} \\ \|\mathcal{D}(\tau)(\mathcal{A}^{0}_{\mu}(\tau))^{-1}\|_{\mathbf{B}(L_{2}(\Omega))^{d}} \lesssim |\tau|^{-1} \\ \|\mathcal{D}(\tau)(\mathcal{A}^{0}_{\mu}(\tau))^{-1}\mathcal{D}(\tau)\|_{\mathbf{B}(L_{2}(\Omega))^{d\times d}} \lesssim 1, \\ \|\mathcal{D}(\tau)\mathcal{D}(\tau)(\mathcal{A}^{0}_{\mu}(\tau))^{-1}\|_{\mathbf{B}(L_{2}(\Omega))^{d\times d}} \lesssim 1, \end{aligned}$$

where the constants depend on d, μ , and the multiplier norms of the coefficients.

Proof. The proof is similar to that of Lemma 4.1.

Since, for every $u \in H^1(\Xi)$,

$$\mathbf{a}^0_{\mu}[u] = \varepsilon^{-2} \int_{\Omega_1^*} \mathbf{a}^0_{\mu}(\tau) [\tilde{u}(k, \cdot)] \, dk,$$

with $\tilde{u} = (\mathcal{GS}^{\varepsilon} \otimes \mathcal{I})u$, it follows that

(4.16)
$$(\mathcal{GS}^{\varepsilon} \otimes \mathcal{I})(\mathcal{A}^{0}_{\mu})^{-1} (\mathcal{GS}^{\varepsilon} \otimes \mathcal{I})^{-1} = \int_{\Omega^{*}_{1}}^{\oplus} \varepsilon^{2} (\mathcal{A}^{0}_{\mu}(\tau))^{-1} dk.$$

As a result, we may restate Theorem 3.1 in terms of the fibers $\mathcal{A}_{\mu}(\tau)$ and $\mathcal{A}_{\mu}^{0}(\tau)$.

THEOREM 4.3. Let $\mu \notin S$. Then for all $\tau \in \mathcal{T}$ it holds that

$$\|(\mathcal{A}_{\mu}(\tau))^{-1} - (\mathcal{A}_{\mu}^{0}(\tau))^{-1}\|_{\mathbf{B}(L_{2}(\Omega))} \lesssim |\tau|^{-1} \\ \|\mathcal{D}_{2}(\tau)(\mathcal{A}_{\mu}(\tau))^{-1} - \mathcal{D}_{2}(\tau)(\mathcal{A}_{\mu}^{0}(\tau))^{-1}\|_{\mathbf{B}(L_{2}(\Omega))^{d}} \lesssim 1,$$

where the constants depend only on d, r_{Λ} , μ , and the multiplier norms of the coefficients.

Let \mathcal{P} denote the orthogonal projection in $L_2(\Omega)$ onto $\mathbb{C} \otimes L_2(\Omega_2)$. Notice that

$$(\mathcal{GS}^{\varepsilon}\otimes\mathcal{I})\mathcal{P}^{\varepsilon}(\mathcal{GS}^{\varepsilon}\otimes\mathcal{I})^{-1}=\int_{\Omega_{1}^{*}}^{\oplus}\mathcal{P}\,dk.$$

Define $\mathcal{K}_{\mu}(\tau) \colon L_2(\Omega) \to \widetilde{H}^1(\Omega)$ by

(4.17)
$$\mathcal{K}_{\mu}(\tau) = (N\mathcal{D}(\tau) + \varepsilon M)(\mathcal{A}^{0}_{\mu}(\tau))^{-1}\mathcal{P}.$$

LEMMA 4.4. For any $\mu \notin \mathcal{S}$ and $\tau \in \mathcal{T}$ we have

$$\begin{aligned} \|\mathcal{D}_{1}\mathcal{K}_{\mu}(\tau)\|_{\mathbf{B}(L_{2}(\Omega))^{d}} \lesssim |\tau|^{-1}, \\ \|\mathcal{D}_{1}\mathcal{D}_{2}(\tau)\mathcal{K}_{\mu}(\tau)\|_{\mathbf{B}(L_{2}(\Omega))^{d\times d}} \lesssim 1, \end{aligned}$$

where the constants depend on d, μ , and the multiplier norms of the coefficients.

Proof. Let $f \in L_2(\Omega)$, and let $u = (\mathcal{A}^0_\mu(\tau))^{-1} \mathcal{P} f$ and $U = \mathcal{K}_\mu(\tau) f$. Then, by Lemma 4.2, it follows that

$$\|\mathcal{D}_{1}U\|_{2,\Omega} \le |\Omega_{1}|^{-1/2} (\|\mathcal{D}_{1}N\|_{\mathbf{M}} + \|\mathcal{D}_{1}M\|_{\mathbf{M}}) \|u\|_{1_{2},2,\Omega;\tau} \lesssim |\tau|^{-1} \|f\|_{2,\Omega}$$

and

$$\begin{aligned} \|\mathcal{D}_{1}\mathcal{D}_{2}(\tau)U\|_{2,\Omega} &\leq |\Omega_{1}|^{-1/2} \left(\|\mathcal{D}_{1}\mathcal{D}_{2}N\|_{\mathbf{M}} + \|\mathcal{D}_{1}\mathcal{D}_{2}M\|_{\mathbf{M}}\right)|\tau|\|u\|_{1_{2},2,\Omega;\tau} \\ &+ |\Omega_{1}|^{-1/2} \left(\|\mathcal{D}_{1}N\|_{\mathbf{M}} + \|\mathcal{D}_{1}M\|_{\mathbf{M}}\right)\|\mathcal{D}_{2}(\tau)u\|_{1_{2},2,\Omega;\tau} \lesssim \|f\|_{2,\Omega}, \end{aligned}$$

as required.

We remark that, since $\mathcal{PK}_{\mu}(\tau) = 0$ (by the definitions of N and M), we may use Poincaré's inequality to see that $\mathcal{K}_{\mu}(\tau)$ and $\mathcal{D}_{2}(\tau)\mathcal{K}_{\mu}(\tau)$ satisfy estimates similar to those for $\mathcal{D}_{1}\mathcal{K}_{\mu}(\tau)$ and $\mathcal{D}_{1}\mathcal{D}_{2}(\tau)\mathcal{K}_{\mu}(\tau)$, respectively. This means that, unlike the case of $\mathcal{A}_{\mu}(\tau)$ and $\mathcal{A}_{\mu}^{0}(\tau)$, where both $\mathcal{D}_{1}(\tau)$ and $\mathcal{D}_{2}(\tau)$ make the norms of the corresponding compositions smaller, roughly speaking, by multiplying each of these norms by

 $|\tau|$, the differentiation \mathcal{D}_1 will not change the order of the norm of $\mathcal{K}_{\mu}(\tau)$. The reason, of course, is that the corrector $\mathcal{K}^{\varepsilon}_{\mu}$ involves functions that rapidly oscillate in the first variable.

In the same fashion as above, we may prove that

(4.18)
$$(\mathcal{GS}^{\varepsilon} \otimes \mathcal{I}) \mathcal{K}^{\varepsilon}_{\mu} (\mathcal{GS}^{\varepsilon} \otimes \mathcal{I})^{-1} = \int_{\Omega^{*}_{1}}^{\oplus} \varepsilon \mathcal{K}_{\mu}(\tau) \, dk.$$

Theorem 3.2 now takes the following form.

THEOREM 4.5. Let $\mu \notin S$. Then for all $\tau \in \mathcal{T}$ it holds that

$$\|\mathcal{D}_1(\tau)(\mathcal{A}_{\mu}(\tau))^{-1} - \mathcal{D}_1(\tau)(\mathcal{A}^0_{\mu}(\tau))^{-1} - \mathcal{D}_1(\tau)\mathcal{K}_{\mu}(\tau)\|_{\mathbf{B}(L_2(\Omega))^d} \lesssim 1,$$

where the constant depends only on d, r_{Λ} , μ , and the multiplier norms of the coefficients.

Let $\mathcal{S}(\tau) \colon \widetilde{H}^1(\Omega) \to \widetilde{H}^1(\Omega)^*$ and $\mathcal{T}(\tau) \colon \widetilde{H}^1(\Omega) \to \widetilde{H}^1(\Omega)^*$ be given by

(4.19)
$$\mathcal{S}(\tau) = \left((\mathbf{k} + \mathcal{D}_2(\tau))^* A + \varepsilon a_1^* \right) (\mathbf{k} + \mathcal{D}_2(\tau)) + \varepsilon (\mathbf{k} + \mathcal{D}_2(\tau))^* a_2 + \varepsilon^2 q,$$

(4.20) $\mathcal{T}(\tau) = \left((\mathcal{K} + \mathcal{D}_2(\tau))^* A + \varepsilon a_1^* \right) \mathcal{D}_1.$

Clearly, $S(\tau)$ and $T(\tau)$ are bounded operators, satisfying estimates like that for $A_{\mu}(\tau)$. It is in fact possible to improve these estimates by using (4.2) and (4.3) instead of (4.4) and (4.5):

$$(4.21) \quad |(\mathcal{S}(\tau)u, v)_{2,\Omega}| \lesssim \left(\varepsilon \|\mathcal{D}_1(\tau)u\|_{2,\Omega} + \|u\|_{1_2,2,\Omega;\tau}\right) \left(\varepsilon \|\mathcal{D}_1(\tau)v\|_{2,\Omega} + \|v\|_{1_2,2,\Omega;\tau}\right),$$

$$(4.22) \quad |(\mathcal{T}(\tau)u, v)_{2,\Omega}| \lesssim \|\mathcal{D}_1 u\|_{2,\Omega} \big(\varepsilon \|\mathcal{D}_1(\tau)v\|_{2,\Omega} + \|v\|_{1_2,2,\Omega;\tau}\big)$$

if $u, v \in \widetilde{H}^1(\Omega)$. The operators $\mathcal{S}(\tau)^+$ and $\mathcal{T}(\tau)^+$ are defined likewise. Of course, estimates similar to (4.21) and (4.22) hold for $\mathcal{S}(\tau)^+$ and $\mathcal{T}(\tau)^+$ as well. Notice that

(4.23)
$$\mathcal{A}_{\mu}(\tau) = \mathcal{D}_{1}^{*}A\mathcal{D}_{1} + \mathcal{S}(\tau) + \mathcal{T}(\tau) + (\mathcal{T}(\tau)^{+})^{*} - \varepsilon^{2}\mu$$

We break Theorem 3.3 into two parts. The first is formulated as follows.

LEMMA 4.6. For any $\mu \notin S$ and $\varepsilon \in \mathscr{C}$ we have

$$\|\mathcal{L}_{\mu}(\mathcal{I}-\mathcal{P}^{\varepsilon})\|_{\mathbf{B}(L_{2}(\Xi))} \lesssim \varepsilon,$$

where the constant depends on d, r_{Λ} , μ , and the multiplier norms of the coefficients.

Proof. We estimate the operator norm of the symbol $\mathcal{L}_{\mu}(\cdot)$. Fix $k \in \mathbb{R}^{d_1} \setminus \{0\}$. Let $f \in L_2(\Omega_2)$, and let $u = (\mathcal{A}^0_{\mu}(\tau))^{-1}f$ and $U = \mathcal{K}_{\mu}(\tau)f$, $U^+ = \mathcal{K}_{\mu}(\tau)^+f$, where $\tau = (k, 1)$. Then

$$\left(\mathcal{L}_{\mu}(k)f,f\right)_{2,\Omega_{2}} = |\Omega_{1}|^{-1} \left(\mathcal{S}(\tau)u + \mathcal{T}(\tau)U,U^{+}\right)_{2,\Omega}$$

It should be noted that while Lemmas 4.2 and 4.4 are only asserted to be valid for $\mathcal{A}^{0}_{\mu}(\tau)$ and $\mathcal{K}_{\mu}(\tau)$ with $\tau \in \Omega_{1}^{*} \times \mathscr{C}$, they may be extended to $\mathcal{A}^{0}_{\mu}(\tau)\mathcal{P}$ and $\mathcal{K}_{\mu}(\tau)$ with $\tau \in \mathbb{R}^{d_{1}} \times \mathscr{C}$. Indeed, the condition $k \in \Omega_{1}^{*}$ is used only to ensure the inequality (4.8). But when u does not depend on x_{1} , we have equality for each $k \in \mathbb{R}^{d_{1}}$. Thus, the estimates (4.21) and (4.22) together with these extended versions of Lemmas 4.2 and 4.4, as well as Poincaré's inequality, give

$$|(\mathcal{L}_{\mu}(k)f, f)_{2,\Omega_{2}}| \lesssim (||u||_{1_{2},2,\Omega;\tau} + ||\mathcal{D}_{1}U||_{2,\Omega}) ||\mathcal{D}_{1}U^{+}||_{1_{2},2,\Omega;\tau} \lesssim |k|^{-1} ||f||_{2,\Omega_{2}}^{2}$$

Now if $g \in C_0^{\infty}(\Xi)$ and $\hat{g} = (\mathcal{F} \otimes \mathcal{I})g$, then

$$\begin{split} |(\mathcal{L}_{\mu}(\mathcal{I}-\mathcal{P}^{\varepsilon})g,g)_{2,\Xi}| &\leq \int_{\mathbb{R}^{d_{1}}\setminus\varepsilon^{-1}\Omega_{1}^{*}} \left| \left(\mathcal{L}_{\mu}(k)\hat{g}(k,\cdot),\hat{g}(k,\cdot)\right)_{2,\Omega_{2}} \right| dk \\ &\lesssim \int_{\mathbb{R}^{d_{1}}\setminus\varepsilon^{-1}\Omega_{1}^{*}} |k|^{-1} \|\hat{g}(k,\cdot)\|_{2,\Omega}^{2} dk \leq \varepsilon r_{\Lambda}^{-1} \|g\|_{2,\Xi}^{2}. \end{split}$$

This is the result that we wished to prove.

The lemma takes care of $\mathcal{L}_{\mu}(\mathcal{I} - \mathcal{P}^{\varepsilon})$ and $\mathcal{L}^{+}_{\mu}(\mathcal{I} - \mathcal{P}^{\varepsilon})$ in the estimate (3.4), so we may concentrate our attention on $\mathcal{L}_{\mu}\mathcal{P}^{\varepsilon}$ and $\mathcal{L}^{+}_{\mu}\mathcal{P}^{\varepsilon}$. Let $\mathcal{A}^{0}_{\mu}(\tau)^{+}$ and $\mathcal{K}_{\mu}(\tau)^{+}$ play the roles of $\mathcal{A}^{0}_{\mu}(\tau)$ and $\mathcal{K}_{\mu}(\tau)$ for $\mathcal{A}_{\mu}(\tau)^{+}$. Then, since

$$((\mathcal{GS}^{\varepsilon}\otimes\mathcal{I})\mathcal{P}^{\varepsilon}u)(k,x) = |\Omega_1|^{-1/2}((\mathcal{S}^{1/\varepsilon}\mathcal{F}\otimes\mathcal{I})u)(k,x_2)$$

for every $u \in L_2(\Xi)$, we find that

(4.24)
$$(\mathcal{GS}^{\varepsilon} \otimes \mathcal{I}) \mathcal{L}_{\mu} \mathcal{P}^{\varepsilon} (\mathcal{GS}^{\varepsilon} \otimes \mathcal{I})^{-1} = \int_{\Omega_{1}^{*}}^{\oplus} \varepsilon \mathcal{L}_{\mu}(\tau) \mathcal{P} dk,$$

where $\mathcal{L}_{\mu}(\tau) \colon L_2(\Omega) \to L_2(\Omega)$ is given by

(4.25)
$$\mathcal{L}_{\mu}(\tau) = \left(\mathcal{K}_{\mu}(\tau)^{+}\right)^{*} \left(\mathcal{S}(\tau)(\mathcal{A}_{\mu}^{0}(\tau))^{-1} + \mathcal{T}(\tau)\mathcal{K}_{\mu}(\tau)\right).$$

Define $\mathcal{L}_{\mu}(\tau)^{+}$ similarly. Obviously, in order to prove Theorem 3.3, we need to establish the following result.

THEOREM 4.7. Let $\mu \notin S$. Then for all $\tau \in \mathcal{T}$ it holds that

$$\left\| (\mathcal{A}_{\mu}(\tau))^{-1} - (\mathcal{A}_{\mu}^{0}(\tau))^{-1} - (\mathcal{K}_{\mu}(\tau) - \mathcal{L}_{\mu}(\tau)) \mathcal{P} - \mathcal{P} (\mathcal{K}_{\mu}(\tau)^{+} - \mathcal{L}_{\mu}(\tau)^{+})^{*} \right\|_{\mathbf{B}(L_{2}(\Omega))} \lesssim 1,$$

where the constant depends only on d, r_{Λ} , μ , and the multiplier norms of the coefficients.

We now turn to the proofs. Our first goal is to verify the identity

(4.26)
$$(\mathcal{A}_{\mu}(\tau))^{-1}\mathcal{P} - (\mathcal{A}_{\mu}^{0}(\tau))^{-1}\mathcal{P} - \mathcal{K}_{\mu}(\tau)$$
$$= -(\mathcal{A}_{\mu}(\tau))^{-1}\mathcal{P}^{\perp} \big(\mathcal{S}(\tau)(\mathcal{A}_{\mu}^{0}(\tau))^{-1}\mathcal{P} + \mathcal{T}(\tau)\mathcal{K}_{\mu}(\tau)\big)$$
$$- (\mathcal{A}_{\mu}(\tau))^{-1} \big(\mathcal{S}(\tau)^{+} + \mathcal{T}(\tau)^{+} - \varepsilon^{2}\bar{\mu}\big)^{*}\mathcal{K}_{\mu}(\tau).$$

Denote the operator on the left by $\mathcal{U}_{\mu}(\tau)$. If $f, g \in L_2(\Omega)$, then we set $u = (\mathcal{A}^0_{\mu}(\tau))^{-1}f$, $U = \mathcal{K}_{\mu}(\tau)f$, and $v^+ = (\mathcal{A}_{\mu}(\tau)^+)^{-1}g$. By (4.11), we have

$$(\mathcal{U}_{\mu}(\tau)f,g)_{2,\Omega} = a^0_{\mu}(\tau)[\mathcal{P}u,v^+] - a_{\mu}(\tau)[\mathcal{P}u+U,v^+]$$

(here we are using the fact that $\mathcal{D}(\tau)$ commutes with \mathcal{P} on periodic functions; hence, so does $\mathcal{A}^{0}_{\mu}(\tau)$). Looking at the definitions of the effective coefficients, we see that

(4.27)
$$\mathcal{A}^{0}_{\mu}(\tau)\mathcal{P} = \mathcal{P}\big(\mathcal{A}_{\mu}(\tau) + \mathcal{T}(\tau)(N\mathcal{D}(\tau) + \varepsilon M)\big)\mathcal{P},$$

from which we obtain

(4.28)
$$a^0_\mu(\tau)[\mathcal{P}u,v^+] - a_\mu(\tau)[\mathcal{P}u,v^+] = (\mathcal{T}(\tau)U,\mathcal{P}v^+)_{2,\Omega} - (\mathcal{A}_\mu(\tau)\mathcal{P}u,\mathcal{P}^\perp v^+)_{2,\Omega}.$$

On the other hand, it follows from (4.23) that

(4.29)
$$a_{\mu}(\tau)[U,v^{+}] = (\mathcal{T}(\tau)U,\mathcal{P}v^{+})_{2,\Omega} + (\mathcal{S}(\tau)U,\mathcal{P}v^{+})_{2,\Omega} + (\mathcal{A}_{\mu}(\tau)U,\mathcal{P}^{\perp}v^{+})_{2,\Omega}.$$

The first term on the right-hand side of this last equality cancels with the first term on the right-hand side of (4.28), so

$$(\mathcal{U}_{\mu}(\tau)f,g)_{2,\Omega} = -(\mathcal{A}_{\mu}(\tau)(\mathcal{P}u+U),\mathcal{P}^{\perp}v^{+})_{2,\Omega} - (\mathcal{S}(\tau)U,\mathcal{P}v^{+})_{2,\Omega})$$

Notice that, by the definitions of N and M,

$$\mathcal{D}_1^* A \mathcal{D}_1 U + (\mathcal{T}(\tau)^+)^* \mathcal{P} u = 0.$$

Since $\mathcal{T}(\tau)\mathcal{P} = 0$, $\mathcal{T}(\tau)^+\mathcal{P} = 0$, and $\mathcal{PK}_{\mu}(\tau) = 0$, this and the identity (4.23) imply that

$$(\mathcal{U}_{\mu}(\tau)f,g)_{2,\Omega} = -\left(\mathcal{S}(\tau)\mathcal{P}u + \mathcal{T}(\tau)U, \mathcal{P}^{\perp}v^{+}\right)_{2,\Omega} - \left(\mathcal{S}(\tau)U + (\mathcal{T}(\tau)^{+})^{*}U - \varepsilon^{2}\mu U, v^{+}\right)_{2,\Omega}.$$

Then, using (4.11) and the fact that $S(\tau)^+$ is the formal adjoint of $S(\tau)$, we get (4.26). Another important identity is

(4.30)
$$\begin{aligned} (\mathcal{A}_{\mu}(\tau))^{-1}\mathcal{P} - (\mathcal{A}_{\mu}^{0}(\tau))^{-1}\mathcal{P} - \mathcal{K}_{\mu}(\tau) + \mathcal{L}_{\mu}(\tau)\mathcal{P} + \mathcal{P}(\mathcal{L}_{\mu}(\tau)^{+})^{*} \\ &= -((\mathcal{A}_{\mu}(\tau)^{+})^{-1} - \mathcal{K}_{\mu}(\tau)^{+})^{*}\mathcal{P}^{\perp} \left(\mathcal{S}(\tau)(\mathcal{A}_{\mu}^{0}(\tau))^{-1}\mathcal{P} + \mathcal{T}(\tau)\mathcal{K}_{\mu}(\tau)\right) \\ &- \left((\mathcal{A}_{\mu}(\tau)^{+})^{-1} - (\mathcal{A}_{\mu}^{0}(\tau)^{+})^{-1}\mathcal{P} - \mathcal{K}_{\mu}(\tau)^{+}\right)^{*} \left(\mathcal{S}(\tau)^{+} + \mathcal{T}(\tau)^{+}\right)^{*} \mathcal{K}_{\mu}(\tau) \\ &- \left(\mathcal{S}(\tau)^{+}\mathcal{K}_{\mu}(\tau)^{+} - \varepsilon^{2}\bar{\mu}(\mathcal{A}_{\mu}(\tau)^{+})^{-1}\right)^{*} \mathcal{K}_{\mu}(\tau). \end{aligned}$$

To prove this, we just note that $\mathcal{T}(\tau)^+ \mathcal{P} = 0$ and $\mathcal{PK}_{\mu}(\tau)^+ = 0$ and then apply (4.26). We denote the operator on the left by $\mathcal{V}_{\mu}(\tau)$.

With these results in hand, it is easy to complete the proofs of the theorems.

Proof of Theorem 4.3. We write

$$(\mathcal{A}_{\mu}(\tau))^{-1} - (\mathcal{A}_{\mu}^{0}(\tau))^{-1} = \mathcal{U}_{\mu}(\tau) + (\mathcal{A}_{\mu}(\tau))^{-1}\mathcal{P}^{\perp} - (\mathcal{A}_{\mu}^{0}(\tau))^{-1}\mathcal{P}^{\perp} + \mathcal{K}_{\mu}(\tau).$$

Poincaré's inequality (2.1) and Lemmas 4.1 and 4.2 (for $\mathcal{A}_{\mu}(\tau)^{+}$ and $\mathcal{A}_{\mu}^{0}(\tau)^{+}$, respectively) show that the norms of $|\tau|(\mathcal{A}_{\mu}(\tau))^{-1}\mathcal{P}^{\perp}$, $\mathcal{D}_{2}(\tau)(\mathcal{A}_{\mu}(\tau))^{-1}\mathcal{P}^{\perp}$ as well as $|\tau|(\mathcal{A}_{\mu}^{0}(\tau))^{-1}\mathcal{P}^{\perp}$, $\mathcal{D}_{2}(\tau)(\mathcal{A}_{\mu}^{0}(\tau))^{-1}\mathcal{P}^{\perp}$ are uniformly bounded. In Lemma 4.4, we proved that so are the norms of $|\tau|\mathcal{K}_{\mu}(\tau)$ and $\mathcal{D}_{2}(\tau)\mathcal{K}_{\mu}(\tau)$. Thus, it is enough to check that

(4.31)
$$\|\mathcal{U}_{\mu}(\tau)\|_{\mathbf{B}(L_{2}(\Omega))} \lesssim |\tau|^{-1},$$

(4.32)
$$\|\mathcal{D}_2(\tau)\mathcal{U}_\mu(\tau)\|_{\mathbf{B}(L_2(\Omega))} \lesssim 1$$

Let the notation be as above. We use (4.21) and (4.22) together with the Poincaré inequality (2.1) to estimate each term in (4.26). The result is that

(4.33)
$$|(\mathcal{U}_{\mu}(\tau)f,g)_{2,\Omega}| \lesssim |\tau|^{-1} (|\tau|||u||_{1_{2},2,\Omega;\tau} + ||\mathcal{D}_{1}U||_{1_{2},2,\Omega;\tau}) \times (||\mathcal{D}_{1}(\tau)v^{+}||_{1_{2},2,\Omega;\tau} + |\tau|||v^{+}||_{1_{2},2,\Omega;\tau}).$$

Combining this with Lemmas 4.1 (for $\mathcal{A}_{\mu}(\tau)^+$), 4.2, and 4.4 gives (4.31).

The inequality (4.32) is proved in a like manner. We set $w^+ = (\mathcal{A}_{\mu}(\tau)^+)^{-1} \mathcal{D}_2(\tau)^* g$ with $g \in L_2(\Omega)^d$ such that $\mathcal{D}_2(\tau)^* g \in L_2(\Omega)$ and then estimate the form

(4.34)
$$(\mathcal{U}_{\mu}(\tau)f, \mathcal{D}_{2}(\tau)^{*}g)_{2,\Omega} = -\left(\mathcal{S}(\tau)\mathcal{P}u + \mathcal{T}(\tau)U, \mathcal{P}^{\perp}w^{+}\right)_{2,\Omega} \\ -\left(U, \left(\mathcal{S}(\tau)^{+} + \mathcal{T}(\tau)^{+} - \varepsilon^{2}\bar{\mu}\right)w^{+}\right)_{2,\Omega}.$$

However, a modification is required to eliminate the mixed second derivatives of w^+ which arise when we estimate the right-hand side (cf. (4.33), where a similar term, namely, $\mathcal{D}_2(\tau)\mathcal{D}_1(\tau)v^+$, causes no difficulty).

We do so as follows. Let $\varphi \in L_2(\Omega_2; \widetilde{H}^1(\Omega_1))$, and let ψ be the solution of

$$\mathcal{D}_2(\tau)^* \mathcal{D}_2(\tau) \psi + |\tau|^2 \psi = \varphi$$

in $L_2(\Omega_1; H^1(\Omega_2))$. Obviously, ψ has first derivatives and mixed second derivatives, as well as pure second derivatives in x_2 . Fix $\ell \in \mathbb{R}^{d_1} \oplus \{0\}$ with $|\ell| = |\tau|$ and define the operator $\mathcal{E}(\tau) \colon L_2(\Omega_2; \tilde{H}^1(\Omega_1)) \to \tilde{H}^1(\Omega)^d$ that assigns to each $\varphi \in L_2(\Omega_2; \tilde{H}^1(\Omega_1))$ the function $(\ell + \mathcal{D}_2(\tau))\psi$. It follows that $(\ell + \mathcal{D}_2(\tau))^* \mathcal{E}(\tau)$ is the identity mapping. A straightforward calculation (using the fact that $\|\mathcal{E}(\tau)\varphi\|_{2,\Omega} = \|\psi\|_{1_2,2,\Omega;\tau}$) shows that $\mathcal{E}(\tau)$ is bounded and

(4.35)
$$|\tau| \|\mathcal{D}_1(\tau)\mathcal{E}(\tau)\varphi\|_{2,\Omega} + \|\mathcal{E}(\tau)\varphi\|_{1_2,2,\Omega;\tau} \le \|\mathcal{D}_1(\tau)\varphi\|_{2,\Omega} + 3\|\varphi\|_{2,\Omega}.$$

Now, we may rewrite the first expression on the right-hand side of (4.34) as

$$\begin{split} \left(\mathcal{S}(\tau)\mathcal{P}u + \mathcal{T}(\tau)U, \mathcal{P}^{\perp}w^{+} \right)_{2,\Omega} \\ &= \varepsilon \left([\mathcal{D}_{2}, \mathcal{S}(\tau)]\mathcal{P}u + [\mathcal{D}_{2}, \mathcal{T}(\tau)]U, \mathcal{E}(\tau)\mathcal{P}^{\perp}w^{+} \right)_{2,\Omega} \\ &+ \left(\mathcal{S}(\tau)(\ell + \mathcal{D}_{2}(\tau))\mathcal{P}u + \mathcal{T}(\tau)(\ell + \mathcal{D}_{2}(\tau))U, \mathcal{E}(\tau)\mathcal{P}^{\perp}w^{+} \right)_{2,\Omega}. \end{split}$$

Applying (4.21) and (4.22) and similar results for the commutators of \mathcal{D}_2 with $\mathcal{S}(\tau)$ and $\mathcal{T}(\tau)$ (notice that these commutators have the same forms as $\mathcal{S}(\tau)$ and $\mathcal{T}(\tau)$), we conclude that

$$\begin{split} \left(\mathcal{S}(\tau)\mathcal{P}u + \mathcal{T}(\tau)U, \mathcal{P}^{\perp}w^{+} \right)_{2,\Omega} \right| \\ \lesssim \left(\|\mathcal{D}_{2}(\tau)u\|_{1_{2},2,\Omega;\tau} + |\tau|\|u\|_{1_{2},2,\Omega;\tau} + \|\mathcal{D}_{1}U\|_{1_{2},2,\Omega;\tau} \right) \\ \times \left(|\tau|\|\mathcal{D}_{1}(\tau)\mathcal{E}(\tau)\mathcal{P}^{\perp}w^{+}\|_{2,\Omega} + \|\mathcal{E}(\tau)\mathcal{P}^{\perp}w^{+}\|_{1_{2},2,\Omega;\tau} \right). \end{split}$$

Therefore, by the estimate (4.35) together with the Poincaré inequality,

(4.36)
$$\frac{\left| \left(\mathcal{S}(\tau) \mathcal{P} u + \mathcal{T}(\tau) U, \mathcal{P}^{\perp} w^{+} \right)_{2,\Omega} \right|}{\lesssim \left(\| \mathcal{D}_{2}(\tau) u \|_{1_{2},2,\Omega;\tau} + |\tau| \| u \|_{1_{2},2,\Omega;\tau} + \| \mathcal{D}_{1} U \|_{1_{2},2,\Omega;\tau} \right) \| \mathcal{D}_{1}(\tau) w^{+} \|_{2,\Omega}.$$

The second expression on the right of (4.34) is handled in the same way as before. The upshot is that

(4.37)
$$\begin{aligned} |(\mathcal{U}_{\mu}(\tau)f, \mathcal{D}_{2}(\tau)^{*}g)_{2,\Omega}| \\ \lesssim \left(\|\mathcal{D}_{2}(\tau)u\|_{1_{2},2,\Omega;\tau} + |\tau|\|u\|_{1_{2},2,\Omega;\tau} + \|\mathcal{D}_{1}U\|_{1_{2},2,\Omega;\tau} \right) \|w^{+}\|_{1,2,\Omega;\tau}, \end{aligned}$$

whence (4.32) follows by Lemmas 4.1 (for $\mathcal{A}_{\mu}(\tau)^+$), 4.2, and 4.4.

Proof of Theorem 4.5. The proof follows the same pattern as the previous one. Again, the assertion is reduced, by Poincaré's inequality and Lemmas 4.1 and 4.2, to the estimation of $\mathcal{D}_1(\tau)\mathcal{U}_\mu(\tau)$. Then the arguments that we used to obtain (4.37) go through without change to yield the desired conclusion.

Proof of Theorem 4.7. We write

$$(\mathcal{A}_{\mu}(\tau))^{-1} - (\mathcal{A}_{\mu}^{0}(\tau))^{-1} - (\mathcal{K}_{\mu}(\tau) - \mathcal{L}_{\mu}(\tau))\mathcal{P} - \mathcal{P}(\mathcal{K}_{\mu}(\tau)^{+} - \mathcal{L}_{\mu}(\tau)^{+})^{*} = \mathcal{V}_{\mu}(\tau) + (\mathcal{P}^{\perp}(\mathcal{A}_{\mu}(\tau)^{+})^{-1} - \mathcal{P}^{\perp}(\mathcal{A}_{\mu}^{0}(\tau)^{+})^{-1} - \mathcal{K}_{\mu}(\tau)^{+})^{*}$$

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(recall that $\mathcal{V}_{\mu}(\tau)$ is the operator on the left side of (4.30)). Since Theorem 4.5 holds true for $\mathcal{A}_{\mu}(\tau)^+$, it follows that

$$\|\mathcal{P}^{\perp}(\mathcal{A}_{\mu}(\tau)^{+})^{-1} - \mathcal{P}^{\perp}(\mathcal{A}_{\mu}^{0}(\tau)^{+})^{-1} - \mathcal{K}_{\mu}(\tau)^{+}\|_{\mathbf{B}(L_{2}(\Omega))} \lesssim 1,$$

where we have used Poincaré's inequality and the fact that $\mathcal{PK}_{\mu}(\tau)^{+} = 0$. Thus, we are left with estimating the operator $\mathcal{V}_{\mu}(\tau)$.

Fix $f, g \in L_2(\Omega)$. Let u, U, and v^+ be as above, and let $u^+ = (\mathcal{A}^0_{\mu}(\tau)^+)^{-1}g$ and $U^+ = \mathcal{K}_{\mu}(\tau)^+g$. Then, by (4.30),

$$\begin{aligned} \left(\mathcal{V}_{\mu}(\tau)f,g\right)_{2,\Omega} &= -\left(\mathcal{S}(\tau)\mathcal{P}u + \mathcal{T}(\tau)U, \mathcal{P}^{\perp}(v^{+} - U^{+})\right)_{2,\Omega} \\ &- \left(U, \left(\mathcal{S}(\tau)^{+} + \mathcal{T}(\tau)^{+}\right)(v^{+} - \mathcal{P}u^{+} - U^{+})\right)_{2,\Omega} \\ &- \left(U, \mathcal{S}(\tau)^{+}U^{+} - \varepsilon^{2}\bar{\mu}v^{+}\right)_{2,\Omega}. \end{aligned}$$

We use (4.36), with $v^+ - U^+$ in place of w^+ , to estimate, dropping the constants, the first expression on the right-hand side by

$$(4.38) \quad \left(\|\mathcal{D}_{2}(\tau)u\|_{1_{2},2,\Omega;\tau} + |\tau|\|u\|_{1_{2},2,\Omega;\tau} + \|\mathcal{D}_{1}U\|_{1_{2},2,\Omega;\tau} \right) \|\mathcal{D}_{1}(\tau)\mathcal{P}^{\perp}(v^{+}-U^{+})\|_{2,\Omega}.$$

The remaining terms, according to estimates similar to (4.21) and (4.22) as well as Poincaré's inequality, do not exceed

$$(4.39) \quad \|\mathcal{D}_{1}U\|_{1_{2},2,\Omega;\tau} \left(\|v^{+}-\mathcal{P}u^{+}-U^{+}\|_{1,2,\Omega;\tau}+\|\mathcal{D}_{1}U^{+}\|_{1_{2},2,\Omega;\tau}+|\tau|\|\mathcal{D}_{1}(\tau)v^{+}\|_{2,\Omega}\right)$$

(notice here that $\mathcal{PK}_{\mu}(\tau) = 0$). Further, by Poincaré's inequality,

$$(4.40) \quad \|\mathcal{D}_{1}(\tau)\mathcal{P}^{\perp}(v^{+}-U^{+})\|_{2,\Omega} \lesssim \|\mathcal{D}_{1}(\tau)(v^{+}-u^{+}-U^{+})\|_{2,\Omega} + \|\mathcal{D}_{1}(\tau)\mathcal{D}_{1}(\tau)u^{+}\|_{2,\Omega}$$

and

(4.41)
$$\|v^{+} - \mathcal{P}u^{+} - U^{+}\|_{1,2,\Omega;\tau} \lesssim \|\mathcal{D}_{1}(\tau)(v^{+} - u^{+} - U^{+})\|_{2,\Omega} + \|v^{+} - u^{+}\|_{1_{2},2,\Omega;\tau} + \|\mathcal{D}_{1}(\tau)u^{+}\|_{1,2,\Omega;\tau} + \|\mathcal{D}_{1}U^{+}\|_{1_{2},2,\Omega;\tau}.$$

If we combine (4.38) with (4.40) and (4.39) with (4.41) and apply Lemma 4.1 (for $\mathcal{A}_{\mu}(\tau)^{+}$), Lemma 4.2 (for $\mathcal{A}_{\mu}^{0}(\tau)$ and $\mathcal{A}_{\mu}^{0}(\tau)^{+}$), and Lemma 4.4 (for $\mathcal{K}_{\mu}(\tau)$ and $\mathcal{K}_{\mu}(\tau)^{+}$) and Theorems 4.3 and 4.5 (for $\mathcal{A}_{\mu}(\tau)^{+}$), then we obtain

$$|(\mathcal{V}_{\mu}(\tau)f,g)_{2,\Omega}| \lesssim ||f||_{2,\Omega} ||g||_{2,\Omega}.$$

This proves the theorem.

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