



A sharp relative-error bound for the Helmholtz *h*-FEM at high frequency

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Abstract

For the *h*-finite-element method (*h*-FEM) applied to the Helmholtz equation, the question of how quickly the meshwidth *h* must decrease with the frequency *k* to maintain accuracy as *k* increases has been studied since the mid 80's. Nevertheless, there still do not exist in the literature any *k*-explicit bounds on the *relative error* of the FEM solution (the measure of the FEM error most often used in practical applications), apart from in one dimension. The main result of this paper is the sharp result that, for the lowest fixed-order conforming FEM (with polynomial degree, *p*, equal to one), the condition " h^2k^3 sufficiently small" is sufficient for the relative error of the FEM solution in 2 or 3 dimensions to be controllably small (independent of *k*) for scattering of a plane wave by a nontrapping obstacle and/or a nontrapping inhomogeneous medium. We also prove relative-error bounds on the FEM solution for arbitrary fixed-order methods applied to scattering by a nontrapping obstacle, but these bounds are not sharp for $p \ge 2$. A key ingredient in our proofs is a result describing the oscillatory behaviour of the solution of the plane-wave scattering problem, which we prove using semiclassical defect measures.

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1 Introduction and informal statement of the main results

1.1 Introduction

When solving the Helmholtz equation $\Delta u + k^2 u = 0$ with the *h* version of the finiteelement method (where accuracy is increased by decreasing the meshwidth *h* while keeping the polynomial degree *p* constant), *h* must decrease faster than k^{-1} to maintain accuracy as *k* increases; this is the so-called "pollution effect" [4].

A thorough investigation of how quickly h must decrease with the frequency k to maintain accuracy as k increases was performed by Ihlenburg and Babuška in the mid 90's [70,71] on the 1-d model problem.

$$u'' + k^2 u = -f$$
 in (0, 1), $u(0) = 0$ and $u'(1) - iku(1) = 0$. (1.1)

An explicit expression for the discrete Green's function for this problem is available, and Ihlenburg and Babuška used this to prove the following two sets of results:

1. The *h*-FEM is quasi-optimal in the H^1 semi-norm, with quasi-optimality constant independent of *k*, if (hk^2/p) is sufficiently small; i.e. there exists c, C > 0, independent of *h*, *k*, and *p* such that, if $hk^2/p \le c$, then

$$\|\nabla(u-u_h)\|_{L^2(0,1)} \le C \min_{v_h \in \mathcal{H}_h} \|\nabla(u-v_h)\|_{L^2(0,1)},$$

where \mathcal{H}_h is the appropriate conforming subspace of $H^1(0, 1)$ of piecewise polynomials of degree p on meshes of width h, and u_h is the Galerkin solution; see [70, Theorem 3], [69, Theorem 4.13], [71, Theorem 3.5] (when p = 1 this result was proved earlier in [3, Theorem 3.2]). The numerical experiments in [70, Figures 8 and 9] then indicated that, when p = 1, the condition " hk^2 sufficiently small" for quasi-optimality is necessary.

Under an assumption on the data *f* (discussed below), the relative error in the *h*-FEM can be made arbitrarily small by, when *p* = 1, making *hk*^{3/2} sufficiently small and, when *p* ≥ 2 and the data is sufficiently smooth (see [69, Remark 4.28]), making *h*^{2p}*k*^{2p+1} sufficiently small. More precisely, [70, Equation 3.25], [71, Theorem 3.7], [69, Equation 4.5.15, §4.6.4, and Theorem 4.27] prove that there exists *C* > 0, independent of *h* and *k* (but dependent on *p*) such that, if *hk* is sufficiently small, then the Galerkin solution *u_h* exists and

$$\frac{\|u - u_h\|_{H^1_k(0,1)}}{\|u\|_{H^1_k(0,1)}} \le C\left(\left(\frac{hk}{p}\right)^p + k\left(\frac{hk}{p}\right)^{2p}\right),\tag{1.2}$$

where the weighted H^1 norm $\|\cdot\|_{H_k^1(0,1)}$ is defined by (3.2) below. The numerical experiments in [70, Figure 11], and [69, Figure 4.13] then indicated that, when p = 1, the condition " h^2k^3 sufficiently small" is necessary for the relative error to be bounded (in agreement with the earlier numerical experiments in [8] for small k).

A note on terminology: following [69–71], we call the regime in h, k, and p where the solution is quasi-optimal (with constant independent of k) the *asymptotic* regime, and the regime where the solution is not quasi-optimal the *preasymptotic* regime. For example, by the results in Points 1 and 2 above, when p = 1 the asymptotic regime is when hk^2 is sufficiently small and the preasymptotic regime is when $hk^2 \gg 1$.

The (asymptotic) quasi-optimality results in Point 1 above have since been generalised to Helmholtz problems in 2 and 3 dimensions (and improved in the case $p \ge 2$). Indeed, the fact that the *h*-FEM with p = 1 is quasi-optimal (with constant independent of *k*) in the full H_k^1 norm when hk^2 is sufficiently small was proved for the homogeneous Helmholtz equation on a bounded domain with impedance boundary conditions in [79, Proposition 8.2.7] (in the case of constant coefficients) and [61, Theorem 4.5 and Remark 4.6(ii)] (in the case of variable coefficients), and for scattering problems with variable coefficients in [50, Theorem 3]. The fact that the *h*-FEM for $p \ge 2$ is quasi-optimal when h^pk^{p+1} is sufficiently small was proved for a variety of constant coefficient Helmholtz problems in [80, Corollary 5.6], [81, Proof of Theorem 5.8], and [51, Theorem 5.1], and for a variety of problems including variable-coefficient Helmholtz problems in [25, Theorem 2.15]; the condition " h^pk^{p+1} sufficiently small" is indicated to be sharp for quasi-optimality by, e.g., the numerical experiments in [25, §4.4].

In contrast, the (preasymptotic) relative-error bound (1.2) in Point 2 above has *not* been obtained for any Helmholtz problem in 2 or 3 dimensions, even though numerical experiments indicate that the condition " $h^{2p}k^{2p+1}$ sufficiently small" is necessary and sufficient for the relative error to be controllably small; see, e.g., [32, Left-hand side of Figure 3]. The closest-available result is that, if $h^{2p}k^{2p+1}$ is sufficiently small, then

$$\|u - u_h\|_{H^1_k(D)} \le C\left((hk)^p + k(hk)^{2p}\right) \|f\|_{L^2(D)},$$
(1.3)

for the Helmholtz problem $\Delta u + k^2 u = -f$ posed in a domain *D* with either impedance boundary conditions on ∂D or a perfectly matched layer (PML). Indeed, for the PML problem, (1.3) is proved for p = 1 in [76, Theorem 4.4 and Remark 4.5(iv)] and [51, Theorem 5.4]. For the impedance problem, (1.3) is proved for p = 1 in [100, Theorem 6.1], for $p \ge 1$ in [32, Corollary 5.2] (following earlier work by [104]), and for $p \ge 1$ for the variable-coefficient Helmholtz equation $\nabla \cdot (A\nabla u) + k^2 nu = -f$ in [87, §2.3] (under a nontrapping condition on A and *n*).

We highlight that, while [32,51,76] all prove results of the form (1.3), all the numerical experiments in these papers consider the *relative error* (either in the H^1 norm [32,76], or the weighted H^1 norm (3.2) [51]), illustrating that relative error is indeed the quantity of interest in practice. An analogous situation is encountered in the preasymptotic error analyses of other Helmholtz FEMs in [14,18,33–35,44,101–103]: all these papers prove bounds on the error in terms of the data, as in (1.3), but all the numerical experiments in these papers concerning the error consider the *relative error*.

1.2 The main results of this paper and their novelty

The two main results are the following:

- (a) Theorem 4.1 proves the relative-error bound (1.2) when p = 1 for scattering of a plane wave by a nontrapping obstacle and/or a nontrapping inhomogeneous medium (modelled by the PDE ∇ · (A∇u) + k²nu = 0 with variable A and n) in 2 or 3 dimensions (see Definition 2.2 below for the precise definition of the boundary-value problems considered). As highlighted above, the numerical experiments in [8,69,70] show that "h²k³ sufficiently small" is necessary for the relative error of the *h*-FEM with p = 1 to be controllably small (independent of k), and so the result of Theorem 4.1 is the sharp bound to which the title of the paper refers.
- (b) Theorem 4.2 proves for $p \ge 2$ a slightly-weaker bound than (1.2), namely that

$$\frac{\|u - u_h\|_{H_k^1(\Omega_R)}}{\|u\|_{H_k^1(\Omega_R)}} \le C(hk + k(hk)^{p+1})$$
(1.4)

for scattering of a plane wave by a nontrapping obstacle in 2 or 3 dimensions, where C in (1.4) is independent of h and k but depends on p, with this dependence given explicitly in the theorem.

As discussed above, these are the first-ever frequency-explicit relative-error bounds on the Helmholtz *h*-FEM in 2 or 3 dimensions. We recall the interest (highlighted at the end of the previous subsection) from [14,18,32-35,44,51,76,100-104] in proving such bounds.

An additional novelty of Theorem 4.1 is that it applies to the variable-coefficient Helmholtz equation, and all the constants in the relative-error bound are explicit, not only in *k* and *h*, but also in the coefficients A and *n*. The only other coefficient-explicit, preasymptotic FEM error bound on the variable-coefficient Helmholtz equation in the literature appears in [87, Theorem 2.39], where the bound (1.3) is proved for the interior impedance problem when $h^{2p}k^{2p+1}$ is sufficiently small and A and *n* are nontrapping. The only other coefficient-explicit FEM error bounds for the Helmholtz equation with variable A and *n* are in [50,61]. Both prove quasi-optimality under the condition " hk^2 sufficient small" when p = 1, with [61, Theorem 3] proving this result for scattering by a nontrapping Dirichlet obstacle.

Our two main results, Theorems 4.1 and 4.2, are proved for a particular class of Helmholtz problems, namely those corresponding to scattering by a plane wave, and not for the equation $\Delta u + k^2 u = -f$ with general $f \in L^2$. We highlight that, for this latter class of problems, it is unreasonable to expect a relative-error bound such as (1.2) to hold, and thus the best one can do is prove bounds for a particular class of realistic data (as we do here). For example, consider the 1-d problem (1.1) with

$$f(x) := -\left[\exp(ik^n x)\chi(x)\right]'' - k^2 \left[\exp(ik^n x)\chi(x)\right],$$
 (1.5)

where χ has compact support in (0, 1). The solution to (1.1) is then $u(x) = \exp(ik^n x)\chi(x)$, which oscillates on a scale of k^{-n} , i.e., a smaller scale than k^{-1} when n > 1. The finite-element method with, say, p = 1 and $hk^{3/2}$ small (and independent of k) will therefore not resolve this solution, and hence a bound such as (1.2) does not hold. This example is nevertheless consistent with the previous results recalled in §1.1 since (i) the assumptions on the solution u in [70, First equation in §3.4] and [71, Definition 3.2] exclude such data f, and (ii) with f given by (1.5), $||f||_{L^2(0,1)} \sim k^{2n}$ and $||u||_{H^1_k(0,1)} \sim k^n$, so that $||f||_{L^2(0,1)} \gg ||u||_{H^1_k(0,1)}$, and the error estimate (1.3) holds in this case because, although the absolute error on left-hand side of (1.3) is large, the right-hand side of (1.3) is larger.

1.3 Discussion of these results in the context of using semiclassical analysis in the numerical analysis of the Helmholtz equation

In the last ~ 10 years, there has been growing interest in using results about the *k*-explicit analysis of the Helmholtz equation from *semiclassical analysis* (a branch of *microlocal analysis*) to design and analyse numerical methods for the Helmholtz equation.¹ The activity has so far occurred in, broadly speaking, five different directions:

- The use of the results in [83] (on the rigorous k → ∞ asymptotics of the solution of the Helmholtz equation in the exterior of a smooth convex obstacle with strictly positive curvature) to design and analyse k-dependent approximation spaces for integral-equation formulations [2,31,36,38,39,53,74,75],
- 2. The use of the results in [83], along with those in [72] on scattering from several convex obstacles, to analyse algorithms for multiple scattering problems [1,11, 37,40].
- 3. The use of bounds on the Helmholtz solution operator (also known as *resolvent estimates*) due to [86,99] (with the latter using the propagation of singularities results in [82]) to prove *k*-explicit bounds on both inverses of boundary-integral operators and the inf-sup constant of the domain-based variational formulation [7,22,23,91], and also to analyse preconditioning strategies [52].
- 4. The use of identities introduced in [86] to prove coercivity of boundary-integral operators [94] and to introduce new coercive formulations of Helmholtz problems [30,55,56,85,93].
- 5. The use of bounds on the restriction of quasimodes of the Laplacian to hypersurfaces from [17,27,64,95–97] to prove sharp *k*-explicit bounds on boundary integral operators [48], [63, Appendix A], [45,49], with these bounds then used to prove sharp *k*-explicit bounds on the number of iterations when GMRES is applied to boundary-integral equations [47].

The results of the present paper include a sixth direction. Namely, a key ingredient in our proofs of Theorems 4.1 and 4.2 (indeed, the ingredient that allows one to obtain a *relative-error* bound instead of a bound in terms of the data, such as (1.3)) is a

¹ A closely-related activity is the design and analysis of numerical methods for the Helmholtz equation based on proving *new* results about the $k \to \infty$ asymptotics of Helmholtz solutions for polygonal obstacles; see [20,21,58,65–67]

result describing the oscillatory behaviour of the solution of the plane-wave scattering problem, which we prove using *semiclassical defect measures*. These measures describe where the mass in phase space of a Helmholtz solution is concentrated in the high-frequency limit (see the discussion in §9.1 below), and were introduced in [57,77]; see [15] for more discussion on the history of defect measures.

2 Formulation of the problem

Assumption 2.1 (Assumptions on the domain and coefficients)

- (i) Ω_− ⊂ ℝ^d, d = 2, 3, is a bounded open Lipschitz set such that its open complement Ω₊ := ℝ^d \ Ω_− is connected.
- (ii) A ∈ C^{0,1}(Ω₊, SPD) (where SPD is the set of d × d real, symmetric, positive-definite matrices) is such that supp(I − A) is compact in ℝ^d and there exist 0 < A_{min} ≤ A_{max} < ∞ such that, for all ξ ∈ ℝ^d,

$$A_{\min}|\boldsymbol{\xi}|^2 \leq \boldsymbol{\xi}^T \left(\mathsf{A}(\mathbf{x})\boldsymbol{\xi} \right) \leq A_{\max}|\boldsymbol{\xi}|^2 \quad \text{for almost every } \mathbf{x} \in \Omega_+.$$
(2.1)

(iii) $n \in L^{\infty}(\Omega_+, \mathbb{R})$ is such that supp(1 - n) is compact in \mathbb{R}^d and there exist $0 < n_{\min} \le n_{\max} < \infty$ such that

$$n_{\min} \le n(\mathbf{x}) \le n_{\max}$$
 for almost every $\mathbf{x} \in \Omega_+$. (2.2)

Figure 1 shows a schematic of Ω_{-} and the supports of I - A and 1 - n. Let the scatterer Ω_{sc} be defined by $\Omega_{sc} := \Omega_{-} \cup \text{supp}(I - A) \cup \text{supp}(1 - n)$ (i.e., the union of the shaded areas in Fig. 1). Given R > 0 such that $\Omega_{sc} \subset B_R$, where B_R denotes the ball of radius R about the origin, let $\Omega_R := \Omega_{+} \cap B_R$. Let $\Gamma_R := \partial B_R$ and let $\Gamma := \partial \Omega_{-}$. Let **n** denote the outward-pointing unit normal vector field on both Γ and Γ_R . We denote by $\partial_{\mathbf{n}}$ the corresponding Neumann trace on Γ or Γ_R and $\partial_{\mathbf{n},A}$ the corresponding conormal-derivative trace. We denote by γu the Dirichlet trace on Γ or Γ_R .

Definition 2.2 (*Helmholtz plane-wave scattering problem*) Given k > 0 and $\mathbf{a} \in \mathbb{R}^d$ with $|\mathbf{a}| = 1$, let $u^I(\mathbf{x}) := e^{ik\mathbf{x}\cdot\mathbf{a}}$. Given Ω_- , A, and *n*, as in Assumption 2.1, we say $u \in H^1_{loc}(\Omega_+)$ satisfies the *Helmholtz plane-wave scattering problem* if

$$\nabla \cdot (A\nabla u) + k^2 n u = 0$$
 in Ω_+ , either $\gamma u = 0$ or $\partial_{\mathbf{n},A} u = 0$ on Γ , (2.3)

and $u^{S} := u - u^{I}$ satisfies the Sommerfeld radiation condition

$$\frac{\partial u^{S}}{\partial r}(\mathbf{x}) - \mathrm{i}ku^{S}(\mathbf{x}) = o\left(\frac{1}{r^{(d-1)/2}}\right)$$
(2.4)

as $r := |\mathbf{x}| \to \infty$, uniformly in $\widehat{\mathbf{x}} := \mathbf{x}/r$.

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Fig. 1 A schematic of Ω_- , the supports of I - A and 1 - n, and B_R



We call a solution of the Helmholtz equation satisfying the Sommerfeld radiation condition (2.4) an *outgoing* solution (so, in Definition 2.2, u^S is outgoing).

Define $DtN_k : H^{1/2}(\Gamma_R) \to H^{-1/2}(\Gamma_R)$ to be the Dirichlet-to-Neumann map for the equation $\Delta u + k^2 u = 0$ posed in the exterior of B_R with the Sommerfeld radiation condition (2.4). When $\Gamma_R = \partial B_R$, for some R > 0, the definition of DtN_k in terms of Hankel functions and polar coordinates (when d = 2)/spherical polar coordinates (when d = 3) is given in, e.g., [80, Equations 3.7 and 3.10]. Let

$$H^1_{0,D}(\Omega_R) := \{ v \in H^1(\Omega_R) : \gamma v = 0 \text{ on } \Gamma \}.$$

When Dirichlet boundary conditions are prescribed in (2.3), let

$$\mathcal{H} := H^1_{0,D}(\Omega_R); \tag{2.5}$$

when Neumann boundary conditions are prescribed, let

$$\mathcal{H} := H^1(\Omega_R). \tag{2.6}$$

Lemma 2.3 (Variational formulation of the Helmholtz plane-wave scattering problem) With u^I , Ω_- , A, n, Ω_R , and \mathcal{H} as above, define $\tilde{u} \in \mathcal{H}$ as the solution of the variational problem

find
$$\widetilde{u} \in \mathcal{H}$$
 such that $a(\widetilde{u}, v) = F(v)$ for all $v \in \mathcal{H}$, (2.7)

where

$$a(\widetilde{u}, v) := \int_{\Omega_R} \left((A\nabla \widetilde{u}) \cdot \overline{\nabla v} - k^2 n \widetilde{u} \overline{v} \right) - \left\langle \text{DtN}_k(\gamma \widetilde{u}), \gamma v \right\rangle_{\Gamma_R}, \quad and$$

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$$F(v) := \int_{\Gamma_R} \left(\partial_{\mathbf{n}} u^I - \mathrm{DtN}_k(\gamma u^I) \right) \overline{\gamma v}.$$
(2.8)

where $\langle \cdot, \cdot \rangle_{\Gamma_R}$ denotes the duality pairing on Γ_R that is linear in the first argument and antilinear in the second. Then $\tilde{u} = u|_{\Omega_R}$, where u is the solution of the Helmholtz plane-wave scattering problem of Definition 2.2.

For a proof of Lemma 2.3, see, e.g., [60, Lemma 3.3]. From here on we denote the solution of the variational problem (2.7) by u, so that u satisfies

$$a(u, v) = F(v)$$
 for all $v \in \mathcal{H}$. (2.9)

Lemma 2.4 The solution of the Helmholtz plane-wave scattering problem of Definition 2.2 exists and is unique.

Proof Uniqueness follows from the unique continuation principle; see [60, §1], [61, §2] and the references therein. Since $a(\cdot, \cdot)$ satisfies a Gårding inequality (see (10.6) below), Fredholm theory then gives existence.

The h finite-element method Let T_h be a family of triangulations of Ω_R (in the sense of, e.g., [28, Page 61]) that is shape regular (see, e.g., [12, Definition 4.4.13], [28, Page 128]). When Neumann boundary conditions are prescribed in (2.3), let

 $\mathcal{H}_h := \{ v \in C(\overline{\Omega_R}) : v |_K \text{ is a polynomial of degree } p \text{ for each } K \in \mathcal{T}_h \}; (2.10)$

when Dirichlet boundary conditions are prescribed we impose the additional condition that elements of \mathcal{H}_h are zero on Γ ; in both cases we then have $\mathcal{H}_h \subset \mathcal{H}$. The main results, Theorems 4.1 and 4.2 below require Γ to be at least $C^{1,1}$. For such Ω_R it is not possible to fit $\partial \Omega_R$ exactly with simplicial elements (i.e. when each element of \mathcal{T}_h is a simplex), and fitting $\partial \Omega_R$ with isoparametric elements (see, e.g., [28, Chapter VI]) or curved elements (see, e.g., [9]) is impractical. Some analysis of non-conforming error is therefore necessary, but since this is very standard (see, e.g., [12, Chapter 10]), we ignore this issue here.

The second main result, Theorem 4.2 (for $p \ge 2$ and analytic Γ), requires the triangulation \mathcal{T}_h to be quasi-uniform in the particular sense of [81, Assumption 5.1]. Triangulations satisfying this assumption can be constructed by refining a fixed triangulation that has analytic element maps; see [81, Remark 5.2].

The finite-element method for the variational problem (2.7) is the Galerkin method applied to the variational problem (2.7), i.e.

find
$$u_h \in \mathcal{H}_h$$
 such that $a(u_h, v_h) = F(v_h)$ for all $v_h \in \mathcal{H}_h$. (2.11)

Observe that setting $v = v_h$ in (2.9) and combining this with (2.11) we obtain the *Galerkin orthogonality* that

$$a(u - u_h, v_h) = 0 \quad \text{for all } v_h \in \mathcal{H}_h. \tag{2.12}$$

3 Definitions of quantities involved in the statement of the main results

Throughout the paper we assume that $R \ge R_0 > 0$ for some fixed $R_0 > 0$ and $k \ge k_0$ for some fixed $k_0 > 0$. For simplicity we assume throughout that

$$k_0 R_0 \ge 1 \quad \text{and} \quad hk \le 1. \tag{3.1}$$

Given a bounded open set D, we let the weighted H^1 norm, $\|\cdot\|_{H^1}$ be defined by

$$\|u\|_{H^{1}_{k}(D)}^{2} := \|\nabla u\|_{L^{2}(D)}^{2} + k^{2} \|u\|_{L^{2}(D)}^{2}.$$
(3.2)

We now define quantities $C_{\text{DtN}j}$, j = 1, 2, C_{sol} , C_{osc} , C_{PF} , C_{H^2} , C_{int} , and C_{MS} that appear in the main results (Theorems 4.1 and 4.2). All of these are dimensionless quantities, independent of k, h, and p, but dependent on one or more of A, n, Ω_{-} (indicated below).

 $C_{\text{DtN}j}$, j = 1, 2 By [80, Lemma 3.3], there exist $C_{\text{DtN}j} = C_{\text{DtN}j}(k_0R_0)$, j = 1, 2, such that

$$\left| \left\langle \mathrm{DtN}_{k}(\gamma u), \gamma v \right\rangle_{\Gamma_{R}} \right\rangle \right| \leq C_{\mathrm{DtN}1} \left\| u \right\|_{H^{1}_{k}(\Omega_{R})} \left\| v \right\|_{H^{1}_{k}(\Omega_{R})}$$
(3.3)

for all $u, v \in H^1(\Omega_R)$ and for all $k \ge k_0$, and

$$-\Re \langle \mathrm{DtN}_k \phi, \phi \rangle_{\Gamma_R} \ge C_{\mathrm{DtN}_2} R^{-1} \|\phi\|_{L^2(\Gamma_R)}^2$$
(3.4)

for all $\phi \in H^{1/2}(\Gamma_R)$ and for all $k \ge k_0$.

 C_{sol} We assume that A, n, and Ω_{-} are *nontrapping* in the sense that there exists $C_{\text{sol}} = C_{\text{sol}}(A, n, \Omega_{-}, R, k_0)$ such that, given $f \in L^2(\Omega_R)$, the solution of the boundary value problem (BVP)

$$\nabla \cdot (\mathsf{A}\nabla v) + k^2 n v = -f \text{ in } \Omega_+, \quad \text{either } \gamma v = 0 \text{ or } \partial_{\mathbf{n},\mathsf{A}} v = 0 \text{ on } \Gamma,$$

and v satisfies the Sommerfeld radiation condition (2.4) (with u^S replaced by v), satisfies the bound

$$\|v\|_{H^{1}(\Omega_{R})} \le C_{\text{sol}} R \|f\|_{L^{2}(\Omega_{+})} \quad \text{for all } k \ge k_{0};$$
(3.5)

observe that the factor R on the right-hand side makes C_{sol} dimensionless. (Remark 4.5 discusses the situation where this nontrapping assumption is removed and C_{sol} depends on k.) This assumption holds if the obstacle Ω_{-} and the coefficients A and n are non-trapping in the sense that all billiard trajectories (or, more precisely, Melrose–Sjöstrand generalized bicharacteristics [68, Section 24.3]) starting in an exterior neighbourhood of Ω_{-} and evolving according to the Hamiltonian flow defined by the symbol of (2.3) escape from that neighbourhood after some uniform time. For this flow to be well-defined, Γ must be C^{∞} , and A and n must be globally $C^{1,1}$ and C^{∞} in a neighbourhood

of Γ ; note that the flow may in general be set-valued rather than unique in cases where the boundary is permitted to be infinite-order flat. Assuming the uniqueness of the flow, an explicit expression for C_{sol} in terms of A, n, Ω_- , and R is then given in [50, Theorems 1 and 2, and Equation 6.32]. However, the bound (3.5) can be established in situations with much less smoothness; indeed, [60, Theorems 2.5, 2.7, and 2.19] establishes (3.5) for a Dirichlet C^0 star-shaped obstacle and L^{∞} A and n satisfying certain monotonicity assumptions. Furthermore, our arguments in the rest of the paper do not need the flow to be well-defined on $\Omega_{sc} := \Omega_- \cup \text{supp}(1-A) \cup \text{supp}(1-n)$, they only require that the bound (3.5) holds. We can therefore define nontrapping in this weaker sense, and work with scatterers of much lower smoothness than in standard microlocal-analysis settings.

 $C_{\rm osc}$ By Theorem 9.1 below, if A, n, and Ω_{-} are nontrapping then there exists $C_{\rm osc} = C_{\rm osc}(A, n, \Omega_{-})$ ('osc' standing for 'oscillation') such that for u a solution to the Helmholtz plane-wave scattering problem of Definition 2.2,

$$|u|_{H^{2}(\Omega_{R})} \leq C_{\text{osc}} k \, \|u\|_{H^{1}_{\mu}(\Omega_{R})}, \qquad (3.6)$$

where $|\cdot|_{H^2(\Omega_R)}$ denotes the H^2 semi-norm; i.e. $|u|_{H^2(\Omega_R)} := \sum_{|\alpha|=2} \int_{\Omega_R} |\partial^{\alpha} u|^2$. C_{PF} By [12, §5.3], [98, Corollary A.15], there exists $C_{\text{PF}} = C_{\text{PF}}(\Omega_-)$ ('PF' standing for 'Poincaré–Friedrichs') such that

$$R^{-2} \|v\|_{L^{2}(\Omega_{R})}^{2} \leq C_{\text{PF}} \left(R^{-1} \|\gamma v\|_{L^{2}(\Gamma_{R})}^{2} + \|\nabla v\|_{L^{2}(\Omega_{R})}^{2} \right)$$
(3.7)

for all $v \in H^1(\Omega_R)$.

 C_{H^2} By Theorem 6.1 below, there exists $C_{H^2} = C_{H^2}(A, \Omega_-)$ such that, if $f \in L^2(\Omega_R)$ and $v \in H^1(\Omega_R)$ satisfy

$$\nabla \cdot (\mathsf{A}\nabla v) = -f \text{ in } \Omega_R, \quad \partial_{\mathbf{n}} v = \mathrm{DtN}_k(\gamma v) \text{ on } \Gamma_R, \text{ and}$$
(3.8a)

either
$$\gamma v = 0$$
 or $\partial_{\mathbf{n}} v = 0$ on Γ , (3.8b)

then

$$\left|v\right|_{H^{2}(\Omega_{R})} \leq C_{H^{2}}\left(\left\|f\right\|_{L^{2}(\Omega_{R})} + R^{-1} \left\|\nabla v\right\|_{L^{2}(\Omega_{R})} + R^{-2} \left\|v\right\|_{L^{2}(\Omega_{R})}\right).$$
(3.9)

The key point in (3.9) is that, although v in (3.8) depends on k via the boundary condition on Γ_R , C_{H^2} is independent of k.

 C_{int} By, e.g., [12, Equation 4.4.28], [90, Theorem 4.1] the *nodal interpolant* I_h : $C(\overline{\Omega_R}) \to \mathcal{H}_h$ is well-defined for functions in $H^2(\Omega_R)$ (for d = 2, 3) and satisfies

$$\|v - I_h v\|_{L^2(\Omega_R)} + h \|\nabla (v - I_h v)\|_{L^2(\Omega_R)} \le C_{\text{int}} h^2 |v|_{H^2(\Omega_R)}, \qquad (3.10)$$

for all $v \in H^2(\Omega_R)$, for some C_{int} that depends only on the shape-regularity constant of the mesh. As a consequence of (3.10), the definition of $\|\cdot\|_{H^1_t(\Omega_R)}$ (3.2), and the

assumption that $hk \leq 1$ (3.1), we have

$$\|v - I_h v\|_{H^1_k(\Omega_R)} \le \sqrt{2C_{\text{int}} h} |v|_{H^2(\Omega_R)}.$$
(3.11)

 $C_{\rm MS}$ By [81, Lemma 3.4 and Proposition 5.3] there exists $C_{\rm MS} = C_{\rm MS}(\Omega_{-})$ ('MS' standing for 'Melenk–Sauter') so that, if Γ is analytic, A = I, n = 1, and Ω_{+} is nontrapping, then the bound (8.6) below holds.

In §1.2 we recalled that the only other frequency- and coefficient-explicit FEM error bounds for the variable-coefficient Helmholtz equation appear in [61, Theorems 4.2 and 4.5], [50, Theorem 3], and [87, Theorem 2.39]. We note here that the constants in these bounds are expressed in terms of analogous quantities to those defined above.

4 Statement and discussion of the main results

4.1 The main results

The first theorem holds for any $p \ge 1$, but is most relevant in the case p = 1.

Theorem 4.1 Let u be the solution of the Helmholtz plane-wave scattering problem (Definition 2.2). Assume that both Assumption 2.1 and (3.1) hold, Ω_{-} is $C^{1,1}$, and A, n, and Ω_{-} are nontrapping. If $p \ge 1$ and

$$h^2 k^3 \le C_1,$$
 (4.1)

then the Galerkin solution u_h to the variational problem (2.11) exists, is unique, and satisfies the bound

$$\frac{\|u - u_h\|_{H^1_k(\Omega_R)}}{\|u\|_{H^1_k(\Omega_R)}} \le C_2 hk + C_3 h^2 k^3, \tag{4.2}$$

where

$$C_{1} := \frac{1}{4(A_{\max} + C_{DtN1})n_{\max}(C_{H^{2}})^{2}(C_{int})^{2}C_{sol}R} \left(n_{\max} + \frac{1}{k_{0}R_{0}C_{sol}} + 2\right)^{-1} \\ \times \left(1 + \frac{\sqrt{2}}{\min\left\{C_{DtN2}(C_{PF})^{-1}, A_{\min}(1 + C_{PF})^{-1}\right\}}\right)^{-1}, \\ C_{2} := \frac{\sqrt{2}C_{int}C_{osc}}{A_{\min}} \left(\max\left\{A_{\max}, n_{\max}\right\} + C_{DtN1}\right),$$

and

$$C_3 := \frac{4\sqrt{2}}{\sqrt{A_{\min}}} \left(A_{\max} + C_{\text{DtN}1} \right) (C_{\text{int}})^2 C_{H^2} C_{\text{sol}} R C_{\text{osc}} \sqrt{n_{\max} + A_{\min}}$$

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$$\times \left(n_{\max} + \frac{1}{k_0 R_0 C_{\text{sol}}} + 2 \right)$$

Theorem 4.2 Let u be the solution of the Helmholtz plane-wave scattering problem (Definition 2.2). Assume that both Assumption 2.1 and (3.1) hold, A = I, n = 1, Ω_{-} is a nontrapping Dirichlet obstacle, Γ is analytic, and the triangulation T_h in the definition of \mathcal{H}_h (2.10) satisfies the quasi-uniformity assumption [81, Assumption 5.1]. If

$$\frac{(hk)^2}{p} + C_{\rm sol} R \frac{k(hk)^{p+1}}{p^p} \le \widetilde{C}_1$$
(4.3)

then the Galerkin solution u_h to the variational problem (2.11) exists, is unique, and satisfies the bound

$$\frac{\|u - u_h\|_{H^1_k(\Omega_R)}}{\|u\|_{H^1_k(\Omega_R)}} \le \left(\widetilde{C}_2 + \frac{\widetilde{C}_3 C_{\rm MS}}{p}\right) hk + \widetilde{C}_3 C_{\rm MS} C_{\rm sol} R \, \frac{k(hk)^{p+1}}{p^p}, \qquad (4.4)$$

where

$$\widetilde{C}_{1} := \frac{1}{2\sqrt{2}(1+C_{\text{DtN1}})C_{H^{2}}C_{\text{MS}}} \left(1 + \frac{\sqrt{2}}{\min\left\{C_{\text{DtN2}}(C_{\text{PF}})^{-1}, (1+C_{\text{PF}})^{-1}\right\}}\right)^{-1}, \\ \widetilde{C}_{2} := \sqrt{2}C_{\text{cont}}C_{\text{int}}C_{\text{osc}}, \quad and \quad \widetilde{C}_{3} := 4(1+C_{\text{DtN1}})C_{\text{int}}C_{\text{osc}}.$$

Observe that (i) the condition (4.3) is satisfied if $h^{p+1}k^{p+2}$ is sufficiently small, and (ii) the bound (4.4) is of the form (1.4).

The result of Theorem 4.2 might appear not to be a high-order result, since the lowest-order terms in (4.3) and (4.4) are h^2 and h, respectively. Nevertheless, for fixed p, if $k(hk)^{p+1}$ is sufficiently small, so that (4.3) is satisfied, then

$$h \sim k^{-1-1/(p+1)}$$
 so $hk \sim k^{-1/(p+1)} \ll 1$ as $k \to \infty$,

and the dominant term on the right-hand side of (4.4) is that involving $k(hk)^{p+1}$. We highlight that Theorem 4.2, along with the previous work discussed in §1.1, shows that high-order methods suffer less from the pollution effect than low-order methods.

4.2 How the main results are proved

Theorems 4.1 and 4.2 are proved using the so-called *elliptic-projection argument* or *modified duality argument*, used to prove the bound (1.3) on the solution in terms of the data. We first make some remarks about the history of this argument, and then outline our new contributions.

Recall that the classic duality argument, coming out of ideas introduced in [89], proves quasi-optimality of the Helmholtz FEM, and was used in, e.g., [3,24,25,50,51, 61,70,79–81,88]. The elliptic-projection argument is a modification of this argument

that allows one to prove results in the preasymptotic regime (as opposed to the asymptotic regime). The initial ideas were introduced in the Helmholtz context in [42,43] for interior-penalty discontinuous Galerkin methods, and then further developed for the standard FEM and continuous interior-penalty methods in [100,104]. The argument has been subsequently used by [6,24,32,51,76,101] (see, e.g., the literature review in [87, §2.3]).

We note that [43,100] also used an error-splitting argument (with this idea called "stability-error iterative improvement" in these papers), and that error splitting ideas were also used in [32], together with the idea of using discrete Sobolev norms in the duality argument. Although we do not use these ideas in this paper, one expects that they could be used to improve the *p* dependence in Theorem 4.2, but see [87, Remark 2.48] for a discussion on the challenges in doing this.

Our three new contributions to the elliptic-projection argument are (i) a rigorous proof, using semiclassical defect measures, of the bound (3.6) describing the oscillatory behaviour of the solution of the plane-wave scattering problem (see Theorem 9.1 below), (ii) the proof of H^2 regularity, with constant independent of k, of the solution of Poisson's equation with the boundary condition $\partial_{\mathbf{n}} v = \text{DtN}_k(\gamma v)$ (see (3.9) and Theorem 6.1), and (iii) determining how all the constants in the elliptic-projection argument depend on A, n, Ω_- , and R.

Regarding (i): oscillatory behaviour similar to (3.6) of Helmholtz solutions has been an assumption in many analyses of finite- and boundary-element methods; see, e.g., [70, First equation in §3.4], [71, Definition 3.2], [13, Definition 4.6], [5, Definition 3.5], [30, Assumption 3.4]. However, to our knowledge, the only existing rigorous results proving such behaviour are [59, Theorems 1.1 and 1.2] and [47, Theorem 1.11(c)]. These results concern the Neumann trace of the solution of the Helmholtz plane-wave scattering problem with A = I and n = 1, and are then used in [59] and [47] to analyse boundary-element methods applied to this problem. In common with (3.6), these results are obtained using semiclassical-analysis techniques.

Regarding (ii): the analogous result (H^2 regularity with constant independent of k) for Poisson's equation with the *impedance boundary condition* $\partial_{\mathbf{n}}v = ik\gamma v$ is central to the elliptic-projection argument for the Helmholtz equation with impedance boundary conditions. This result was explicitly assumed in [43, Lemma 4.3], implicitly assumed in [6,24,100,104], and recently proved in [26]. Our proof of (3.9) uses (and makes A-explicit) arguments from [26], which in turn use results from [62], adapting them to deal with the operator DtN_k, instead of *ik*, in the boundary condition.

Regarding (iii): while the standard duality argument applied to the Helmholtz equation discussed above has recently been made explicit in A, n, and Ω_{-} in [50,61] (as discussed in §1.2), the only places in the literature where the elliptic-projection argument is made explicit in A, n, and Ω_{-} are the present paper and [87, §2.3], leading to the coefficient-explicit preasymptotic error bounds on the Helmholtz FEM at highfrequency in Theorem 4.1 and [87, Theorem 2.39]. One area in which we expect these results to be applied is in the analysis of *uncertainty quantification (UQ)* algorithms for the high-frequency Helmholtz equation with random coefficients, as discussed in the following remark. *Remark 4.3* (The importance of coefficient-explicit FEM results for Helmholtz UQ) To analyse UQ algorithms that use the standard Helmholtz variational formulation, one needs to understand how existence and uniqueness of the Galerkin solution is affected by the randomness in the coefficients. One therefore needs coefficient-explicit existence and uniqueness results for the Galerkin solution for the Helmholtz equation with variable (deterministic) coefficients (such as in Theorem 4.1 and [87, Theorem 2.39]); this issue is highlighted (but not fully analysed) in the analysis of Monte Carlo and Multi-level Monte Carlo methods in [87, Chapter 5]; see [87, Assumption 5.1 and Remark 5.2].

The only other analyses of uncertainty quantification (UQ) algorithms for the highfrequency Helmholtz equation with random coefficients in the literature are [41,54] (concerning Monte Carlo and Quasi-Monte Carlo methods, respectively). Because of the issue described in the previous paragraph, these papers use formulations of the Helmholtz equation where existence and uniqueness of the Galerkin solution is established for all k, h, p, and for a class of (deterministic) coefficients ([41] uses the interior-penalty discontinuous-Galerkin method of [42,43,54] uses the coercive formulation of [56]). This then ensures that the Galerkin solution exists and is unique for all realisations of the random coefficients; see the discussion at the beginning of [41, §4].

4.3 Why does Theorem 4.2 not cover scattering by an inhomogeneous medium?

In both the elliptic-projection argument and the standard duality argument, a key role is played by the quantity $\eta(\mathcal{H}_h)$ defined by (8.3) below, which describes how well solutions of the (adjoint of the) Helmholtz equation can be approximated in \mathcal{H}_h .

In the case p = 1 we estimate $\eta(\mathcal{H}_h)$ using H^2 regularity of the solution (which holds when A and Ω_- satisfy the assumptions of Theorem 4.1), leading to the bound (8.5) below. When $p \ge 1$, A = I, n = 1, Ω_- is a Dirichlet obstacle, and Γ is analytic, [81] proved the bound (8.6) on $\eta(\mathcal{H}_h)$, and we use this result to prove Theorem 4.2. The bound (8.6) was proved via a judicious splitting of the solution [81, Theorem 4.20] into an analytic but oscillating part, and an H^2 part that behaves "well" for large frequencies, and this splitting is only available for the exterior Dirichlet problem with A = I and n = 1.

We highlight that an alternative splitting procedure valid for Helmholtz problems with variable coefficients was recently developed in [25], leading to an alternative proof of the bound on $\eta(\mathcal{H}_h)$ (8.6) [25, Lemma 2.13]. However, this alternative procedure requires that DtN_k be approximated by *ik* on Γ_R . Indeed, in [25, Proof of Lemma 2.13] the solution is expanded in powers of k, i.e. $u = \sum_{j=0}^{\infty} k^j u_j$, and then on Γ_R one has $\partial_n u_{j+1} = i\gamma u_j$; this relationship between u_{j+1} and u_j on Γ_R no longer holds if DtN_k is not approximated by *ik*.

4.4 Approximating DtN_k

Implementing the operator DtN_k is computationally expensive, and so in practice one seeks to approximate this operator by *either* imposing an absorbing boundary condition on Γ_R , or using a PML. In this paper we follow the precedent established in [80,81] of, when proving new results about the FEM for exterior Helmholtz problems, first assuming that DtN_k is realised exactly. We remark, however, that if the two key ingredients in §4.2 (a proof of the oscillatory behaviour (3.6) and H^2 -regularity, independent of k, of a Poisson problem) can be established when DtN_k is replaced by an absorbing boundary condition on Γ_R , then the result of Theorem 4.1 carry over to this case. When an impedance boundary condition (i.e. the simplest absorbing boundary condition) is imposed on Γ_R , the necessary Poisson H^2 -regularity result is proved in [26], but we discuss below in Remark 9.9 the difficulties in proving (3.6) in

4.5 Removing the nontrapping assumption

The only place in the proofs of Theorems 4.1 and 4.2 where the nontrapping assumption (i.e. the fact that C_{sol} in (3.5) is independent of k) is used is in the proof of the bound (3.6) (in Theorem 9.1 below). We sketch in Remark 9.10 below how (3.6) can be proved in the trapping case (i.e. when C_{sol} is not independent of k); the rest of the proofs of Theorems 4.1 and 4.2 then go through as before. In the case of Theorem 4.1, the requirement for the relative error to be bounded independently of k would then be that $h^2k^3C_{sol}$ be sufficiently small. Under the strongest form of trapping, C_{sol} can grow exponentially through a sequence of ks [10, §2.5], but is bounded polynomially in k if a set of frequencies of arbitrarily-small measure is excluded [73, Theorem 1.1]. However, it is not clear how sharp the requirement " $h^2k^3C_{sol}$ sufficiently small" for the relative error to be bounded is in these cases.

5 Outline of the proof

this case.

As highlighted in §4.2, one of the novelties of this paper is that it makes the ellipticprojection argument explicit in the coefficients A and *n*. However, this explicitness means that many of the expressions in the proofs are complicated (in the same way as the expressions in the results in Theorems 4.1 and 4.2 are complicated). In this section therefore, we give an outline of the proof, keeping track of the dependence on *k*, *h*, and *p*, but ignoring the dependence on A, *n*, Ω_- , and *R*. We use the notation $a \leq b$ when $a \leq Cb$ with *C* independent of *k*, *h*, and *p*, but dependent on A, *n*, Ω_- , and *R*.

As in the standard duality argument coming out of ideas introduced in [89] and then formalised in [88], our starting point is the fact that, since $a(\cdot, \cdot)$ satisfies the Gårding inequality (10.6), Galerkin orthogonality (2.12) and continuity of $a(\cdot, \cdot)$ (10.4) imply that, for any $v_h \in \mathcal{H}_h$,

$$A_{\min} \|u - u_{h}\|_{H^{1}_{k}(\Omega_{R})}^{2} \leq \Re a(u - u_{h}, u - v_{h}) + k^{2} (n_{\max} + A_{\min}) \|u - u_{h}\|_{L^{2}(\Omega_{R})}^{2}$$

$$\leq C_{\text{cont}} \|u - u_{h}\|_{H^{1}_{k}(\Omega_{R})} \|u - v_{h}\|_{H^{1}_{k}(\Omega_{R})} + k^{2} (n_{\max} + A_{\min}) \|u - u_{h}\|_{L^{2}(\Omega_{R})}^{2}.$$

(5.1)

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Recall (from, e.g., [88, Theorem 2.5], [80, Theorem 4.3], [92, Theorem 6.32]) that the standard duality argument (related to the Aubin-Nitsche trick) shows that

$$\|u - u_h\|_{L^2(\Omega_R)} \le C_{\text{cont}}\eta(\mathcal{H}_h) \|u - u_h\|_{H^1(\Omega_R)},$$
(5.2)

where $\eta(\mathcal{H}_h)$, defined by (8.3) below, describes how well solutions of the adjoint problem are approximated in the space \mathcal{H}_h . Inputting (5.2) into (5.1) one obtains quasioptimality, with constant independent of k, if $k\eta(\mathcal{H}_h)$ is sufficiently small. Lemma 8.2 below shows that $\eta(\mathcal{H}_h) \leq h + (hk)^p$, and thus the condition " $k\eta(\mathcal{H}_h)$ sufficiently small" is satisfied if $h^p k^{p+1}$ is sufficiently small.

In contrast, the elliptic-projection argument, which we follow, shows that

$$\|u - u_h\|_{L^2(\Omega_R)} \lesssim \eta(\mathcal{H}_h) \|u - w_h\|_{H^1_k(\Omega_R)} \quad \text{for all } w_h \in \mathcal{H}_h, \tag{5.3}$$

provided that $hk^2\eta(\mathcal{H}_h)$ is sufficiently small (see Lemma 10.1 below). Observe that (5.3) is a stronger bound than (5.2), since w_h on the right-hand side of (5.3) is arbitrary. The proof of (5.3) in our setting of the plane-wave scattering problem requires the new Poisson H^2 -regularity bound (3.9), which we prove in Theorem 6.1 below.

Inputting (5.3) into (5.1), choosing $w_h = v_h$, and using the inequality

$$2\alpha\beta \le \varepsilon\alpha^2 + \varepsilon^{-1}\beta^2 \quad \text{for all } \alpha, \beta, \varepsilon > 0, \tag{5.4}$$

on the first term on the right-hand side of (5.1), we obtain that, if $hk^2\eta(\mathcal{H}_h)$ is sufficiently small, then, for any $v_h \in \mathcal{H}_h$,

$$\|u-u_h\|_{H_k^1(\Omega_R)}^2 \lesssim (1+k^2(\eta(\mathcal{H}_h))^2) \|u-v_h\|_{H_k^1(\Omega_R)}^2$$

i.e. quasi-optimality. Assuming H^2 regularity of the solution, and using (3.11), we obtain that, if $hk^2\eta(\mathcal{H}_h)$ is sufficiently small, then

$$\|u - u_h\|_{H^1_k(\Omega_R)}^2 \lesssim \left(1 + k^2 (\eta(\mathcal{H}_h))^2\right) h^2 |u|_{H^2(\Omega_R)}^2.$$
(5.5)

In the standard elliptic-projection argument (see, e.g., [24, §5.5]) applied to the PDE $\Delta u + k^2 u = -f$, an H^2 -regularity bound similar to (3.5) and the nontrapping bound (3.5) are combined to give $|u|_{H^2(\Omega_R)} \leq k ||f||_{L^2(\Omega_R)}$, and combining this with both (5.5) and the bound $\eta(\mathcal{H}_h) \leq hk$ (see (8.5) below) proves the bound (1.3) with p = 1 on the Galerkin error in terms of the data when h^2k^3 is sufficiently small.

In contrast, in this paper we prove, using semiclassical defect measures, that the solution to the plane-wave scattering problem satisfies (3.6), i.e. $|u|_{H^2(\Omega_R)} \lesssim k ||u||_{H^1_k(\Omega_R)}$, (see Theorem 9.1 below), and using this in (5.5), along with the bounds on $\eta(\mathcal{H}_h)$ in Lemma 8.2, we obtain the relative-error bounds (4.2) and (4.4).

In summary, once one has proved the bound (3.6) (which we do via semiclassical analysis) and the Poisson H^2 -regularity bound (3.9) (which we do using results from [62] and properties of DtN_k), if one ignores the technicalities of making the argument

explicit in A, n, Ω_{-} , and R, then the proof of a preasymptotic relative-error bound follows via a straightforward modification of the elliptic-projection argument. Given the large and sustained interest (reviewed in §1.1) in preasymptotic relative-error bounds for the Helmholtz FEM, we believe this fact illustrates the advantage of approaching the numerical analysis of the Helmholtz equation from a perspective encompassing both numerical-analysis and semiclassical-analysis techniques.

6 Proof of the Poisson H²-regularity result (3.9)

Theorem 6.1 With A, Ω_- , Γ , and Ω_R as in §2, let $v \in H^1(\Omega_R)$ be the solution of the Poisson boundary value problem (3.8). If Γ is $C^{1,1}$, then $v \in H^2(\Omega_R)$ and the bound (3.9) holds.

We follow the recent proof of the related regularity result [26, Theorem 3.1] (where DtN_k is replaced by ik, A = I, and $\Omega_{-} = \emptyset$) and start by recalling results from [62].

Lemma 6.2 Let D be a bounded, convex, open set of \mathbb{R}^n with C^2 boundary. Then, for all $\mathbf{v} \in H^1(D; \mathbb{C}^d)$,

$$\int_{D} \left(|\nabla \cdot \mathbf{v}|^{2} - \sum_{i,j=1}^{n} \int_{D} \frac{\partial v_{i}}{\partial x_{j}} \frac{\partial v_{j}}{\partial x_{i}} \right) \geq -2\Re \langle (\gamma \mathbf{v})_{T}, \nabla_{T} (\gamma \mathbf{v} \cdot \mathbf{n}) \rangle_{\partial D}, \quad (6.1)$$

where ∇_T is the surface gradient on ∂D and $(\gamma \mathbf{v})_T := \gamma \mathbf{v} - \mathbf{n}(\gamma \mathbf{v} \cdot \mathbf{n})$ is the tangential component of $\gamma \mathbf{v}$.

Proof The result with v real follows from [62, Theorem 3.1.1.1] and the fact that the second fundamental form of ∂D (defined in, e.g., [62, §3.1.1]) is non-positive (see [62, Proof of Theorem 3.1.2.3]). The result with v complex follows in a straightforward way by repeating the argument in [62, Theorem 3.1.1.1] for complex v.

Lemma 6.3 ([62, Lemma 3.1.3.4]) If $A \in C^{0,1}(D, SPD)$ satisfies (2.1) (with Ω_+ replaced by D), then, for all $v \in H^2(D)$,

$$(A_{\min})^{2} \sum_{i,j=1}^{d} \left| \frac{\partial^{2} v}{\partial x_{i} \partial x_{j}} \right|^{2} \leq \sum_{i,j,\ell,m=1}^{d} A_{i\ell} A_{jm} \frac{\partial^{2} v}{\partial x_{j} \partial x_{\ell}} \frac{\partial^{2} \overline{v}}{\partial x_{i} \partial x_{m}}.$$
 (6.2)

As a first step to proving Theorem 6.1, we prove it in the case when $\Omega_{-} = \emptyset$.

Lemma 6.4 Let $A \in C^{0,1}(B_R, SPD)$ satisfy (2.1) (with Ω_+ replaced by B_R) and be such that supp $(I-A) \subset B_R$. Given $f \in L^2(B_R)$, let $v \in H^1(B_R)$ be the solution of

$$\nabla \cdot (A\nabla v) = -f \quad in \ B_R, \qquad \partial_{\mathbf{n}} v = \mathrm{DtN}_k(\gamma v) \ on \ \Gamma_R. \tag{6.3}$$

Then $v \in H^2(B_R)$ and

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$$\begin{aligned} \|v\|_{H^{2}(B_{R})}^{2} &\leq \frac{2}{(A_{\min})^{2}} \bigg[\|f\|_{L^{2}(B_{R})}^{2} + \left(d^{4} \|\nabla A\|_{L^{\infty}(B_{R})}^{2} \right. \\ &+ \frac{2}{(A_{\min})^{2}} d^{8} \|A\|_{L^{\infty}(B_{R})}^{2} \|\nabla A\|_{L^{\infty}(B_{R})}^{2} \bigg) \|\nabla v\|_{L^{2}(B_{R})}^{2} \bigg], \end{aligned}$$

where ∇A denotes the derivative of A.

Proof Let $w \in H^1(\mathbb{R}^d)$ be the outgoing solution of the following transmission problem

$$\nabla \cdot (\mathsf{A}\nabla w) = -f \quad \text{in } B_R, \qquad \Delta w + k^2 w = 0 \quad \text{in } \mathbb{R}^d \setminus \overline{B_R},$$
$$\gamma w_+ = \gamma w_- \text{ and } \partial_{\mathbf{n}} w_+ = \partial_{\mathbf{n}} w_- \quad \text{on } \Gamma_R,$$

where $w_- := w|_{B_R}$ and $w_+ := w|_{\mathbb{R}^d \setminus B_R}$. (Note that it is important here that A = I in a neighbourhood of Γ_R , so that $\partial_{\mathbf{n},A}w_- = \partial_{\mathbf{n}}w_-$.) By the definition of the operator DtN_k, $w_- = v$. Since Γ_R is C^2 , the regularity result [29, Theorem 5.2.1 and §5.4b] implies that $w_- \in H^2(B_R)$ and $w_+ \in H^2_{\text{loc}}(\mathbb{R}^d \setminus \overline{B_R})$; therefore $v \in H^2(B_R)$.

Since $v \in H^2(B_R)$ and A is Lipschitz, $A\nabla v \in H^1(B_R)$ and we can apply Lemma 6.2 with $\mathbf{v} := A\nabla v$. Since A = I near Γ_R , $\mathbf{v} = \nabla v$ near Γ_R and so the right-hand side of (6.1) becomes

$$-2\Re \langle \nabla_T(\gamma v), \nabla_T(\partial_{\mathbf{n}} v) \rangle_{\Gamma_R} = -2\Re \langle \nabla_T(\gamma v), \nabla_T(\mathrm{DtN}_k(\gamma v)) \rangle_{\Gamma_R}$$

where we have used the boundary condition in (6.3).

Now, DtN_k and ∇_T commute on Γ_R ; this can be seen either by rotation invariance, or by using the definition of DtN_k and ∇_T in terms of Fourier series on Γ_R . Therefore, the inequality (3.4) implies that the right-hand side of (6.1) is non-negative, hence

$$\sum_{i,j,\ell,m=1}^{d} \int_{B_{R}} \frac{\partial}{\partial x_{j}} \left(A_{i\ell} \frac{\partial v}{\partial x_{\ell}} \right) \frac{\partial}{\partial x_{i}} \left(A_{jm} \frac{\partial \overline{v}}{\partial x_{m}} \right) \le \|f\|_{L^{2}(B_{R})}^{2}.$$
(6.4)

The left-hand side of (6.4) equals

$$\sum_{i,j,\ell,m=1}^{d} \int_{\Omega} A_{i\ell} A_{jm} \frac{\partial^2 v}{\partial x_j \partial x_\ell} \frac{\partial^2 \overline{v}}{\partial x_i \partial x_m} + \sum_{i,j,\ell,m=1}^{d} \int_{\Omega} R_{i,j,\ell,m}, \quad (6.5)$$

where

$$R_{i,j,\ell,m} = \frac{\partial A_{i\ell}}{\partial x_j} \frac{\partial v}{\partial x_\ell} A_{jm} \frac{\partial^2 \overline{v}}{\partial x_i \partial x_m} + A_{i\ell} \frac{\partial^2 v}{\partial x_j \partial x_\ell} \frac{\partial A_{jm}}{\partial x_i} \frac{\partial \overline{v}}{\partial x_m} + \frac{\partial A_{i\ell}}{\partial x_j} \frac{\partial v}{\partial x_\ell} \frac{\partial A_{jm}}{\partial x_i} \frac{\partial \overline{v}}{\partial x_m} =: R_{i,j,\ell,m}^1 + R_{i,j,\ell,m}^2 + R_{i,j,\ell,m}^3.$$

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By the Cauchy-Schwarz inequality

$$\left| \int_{B_R} R^1_{i,j,\ell,m} \right| + \left| \int_{B_R} R^2_{i,j,\ell,m} \right| \le 2 \|\mathsf{A}\|_{L^{\infty}(B_R)} \|\nabla\mathsf{A}\|_{L^{\infty}(B_R)} \|\nabla v\|_{L^2(B_R)} |v|_{H^2(B_R)} \|v\|_{L^2(B_R)} \|v\|_$$

and

$$\left|\int_{B_R} R_{i,j,\ell,m}^3\right| \leq \|\nabla \mathsf{A}\|_{L^\infty(B_R)}^2 \|\nabla v\|_{L^2(B_R)}^2.$$

We therefore obtain

$$\left|\sum_{i,j,\ell,m=1}^{d} \int_{B_{R}} R_{i,j,\ell,m}\right| \leq 2d^{4} \|\mathbf{A}\|_{L^{\infty}(B_{R})} \|\nabla\mathbf{A}\|_{L^{\infty}(B_{R})} \|\nabla v\|_{L^{2}(B_{R})} |v|_{H^{2}(B_{R})} + d^{4} \|\nabla\mathbf{A}\|_{L^{\infty}(B_{R})}^{2} \|\nabla v\|_{L^{2}(B_{R})}^{2}.$$

Combining this with (6.2), (6.4), and (6.5), we obtain

$$(A_{\min})^{2} |v|_{H^{2}(B_{R})}^{2} \leq ||f||_{L^{2}(B_{R})}^{2} + 2d^{4} ||A||_{L^{\infty}(B_{R})} ||\nabla A||_{L^{\infty}(B_{R})} ||\nabla v||_{L^{2}(B_{R})} |v|_{H^{2}(B_{R})} + d^{4} ||\nabla A||_{L^{\infty}(B_{R})}^{2} ||\nabla v||_{L^{2}(B_{R})}^{2}.$$

Using (5.4) on the second term on the right-hand side, we obtain the result.

We now use Lemma 6.4 to prove Theorem 6.1.

Proof (Proof of Theorem 6.1) Let $0 < R_0 < R_1 < R$ be such that $\overline{\Omega_-} \subset B_{R_0}$, and let $\chi \in C^{\infty}(\mathbb{R}^d)$ be such that $0 \le \chi \le 1$ and

$$\chi = 0$$
 in B_{R_0} and $\chi = 1$ in $\mathbb{R}^d \setminus B_{R_1}$.

We decompose v as

$$v = \chi v + (1 - \chi)v =: v_1 + v_2.$$
(6.6)

Then $v_1 \in H^1(B_R)$ and satisfies

$$\nabla \cdot (\mathsf{A} \nabla v_1) = -\chi f + \nabla \chi \cdot (\mathsf{A} \nabla v) + \nabla v \cdot (\mathsf{A} \nabla \chi) + v \nabla \cdot (\mathsf{A} \nabla \chi) \quad \text{in } B_R,$$

and $\partial_{\mathbf{n}} v_1 = \text{DtN}_k(\gamma v_1)$ on Γ_R . Lemma 6.4 implies that $v_1 \in H^2(B_R)$ and that there exists $C_4 = C_4(A, d, \chi) > 0$ such that

$$|v_1|_{H^2(\Omega_R)} \le C_4 \left(\|f\|_{L^2(\Omega_R)} + R^{-1} \|\nabla v\|_{L^2(\Omega_R)} + R^{-2} \|v\|_{L^2(\Omega_R)} \right), \quad (6.7)$$

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where (i) we have used the fact that $\nabla \chi = 0$ in a neighbourhood of Ω_{-} to write all the norms as norms over Ω_R , and (ii) we have inserted the inverse powers of R on the right-hand side to keep C_4 a dimensionless quantity. On the other hand, v_2 satisfies

$$\nabla \cdot (\mathsf{A}\nabla v_2) = -(1-\chi)f - \nabla \chi \cdot (\mathsf{A}\nabla v) - \nabla v \cdot (\mathsf{A}\nabla \chi) - v \nabla \cdot (\mathsf{A}\nabla \chi) \text{ in } B_R,$$

 $v_2 = 0$ in $B_R \setminus B_{R_1}$, and either $\gamma v_2 = 0$ or $\partial_{\mathbf{n}} v_2 = 0$ on Γ .

Since A is Lipschitz, $A_{\min} > 0$, and both Γ and Γ_R are $C^{1,1}$, [62, Theorems 2.3.3.2, 2.4.2.5, and 2.4.2.7] imply that, if $w \in H^1(\Omega_R)$, $\nabla \cdot (A\nabla w) \in L^2(\Omega_R)$, and either $\gamma w = 0$ or $\partial_{\mathbf{n}} w = 0$ on $\partial \Omega_R$, then $w \in H^2(\Omega_-)$ and there exists $C_5 = C_5(A, \Omega_-, d, R) > 0$ such that

$$|w|_{H^{2}(\Omega_{R})} \leq C_{5} \left(\|\nabla \cdot (\mathsf{A}\nabla w) - w\|_{L^{2}(\Omega_{R})} + R^{-1} \|\nabla w\|_{L^{2}(\Omega_{R})} + R^{-2} \|w\|_{L^{2}(\Omega_{R})} \right).$$

Applying this with $w = v_2$, we obtain that

$$|v_2|_{H^2(\Omega_R)} \le C_6 \left(\|f\|_{L^2(\Omega_R)} + R^{-1} \|\nabla v\|_{L^2(\Omega_R)} + R^{-2} \|v\|_{L^2(\Omega_R)} \right), \quad (6.8)$$

and the bound (3.9) follows from combining (6.7) and (6.8) using (6.6).

7 The elliptic projection and associated results

Define the sesquilinear form $a_{\star}(\cdot, \cdot)$ by

$$a_{\star}(u,v) := \int_{\Omega_R} \mathsf{A}\nabla u \cdot \overline{\nabla v} - \langle \mathrm{DtN}_k \gamma u, \gamma v \rangle_{\Gamma_R}.$$
(7.1)

Recall from (2.5) and (2.6) that \mathcal{H} equals *either* $H^1_{0,D}(\Omega_R)$ (with Dirichlet conditions in (2.3)) *or* $H^1(\Omega_R)$ (with Neumann conditions).

Lemma 7.1 (Continuity and coercivity of $a_{\star}(\cdot, \cdot)$) For all $u, v \in \mathcal{H}$,

$$|a_{\star}(u, v)| \leq C_{\text{cont}\star} \|u\|_{H^{1}_{k}(\Omega_{R})} \|v\|_{H^{1}_{k}(\Omega_{R})} \quad and \quad \Re a_{\star}(v, v) \geq C_{\text{coer}\star} \|v\|_{H^{1}_{R}(\Omega_{R})}^{2},$$
(7.2)

where

$$C_{\text{cont}\star} := A_{\text{max}} + C_{\text{DtN}1}, \quad C_{\text{coer}\star} := \min\left\{C_{\text{DtN}2}(C_{\text{PF}})^{-1}, A_{\min}(1+C_{\text{PF}})^{-1}\right\}$$

and

$$\|v\|_{H^1_R(\Omega_R)}^2 := \|\nabla u\|_{L^2(\Omega_R)}^2 + \frac{1}{R^2} \|v\|_{L^2(\Omega_R)}^2.$$
(7.3)

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Proof The first inequality in (7.2) follows from the inequality (3.3) and the Cauchy–Schwarz inequality. The second inequality in (7.2) follows from (3.4) and (3.7). \Box

As a consequence of Lemma 7.1, we have

$$C_{\operatorname{corr} \star} \|v\|_{H^1_R(\Omega_R)}^2 \le \left|a_\star(v,v)\right| \le C_{\operatorname{corr} \star} \|v\|_{H^1_k(\Omega_R)}^2 \quad \text{for all } v \in \mathcal{H},$$
(7.4)

and we then define the new norm on \mathcal{H} ,

$$\|v\|_{\star} := \sqrt{a_{\star}(v, v)}.$$

Lemma 7.2 (Bounds on the solution of the variational problem associated with $a_{\star}(\cdot, \cdot)$) *The solution of the variational problem*

find
$$u \in \mathcal{H}$$
 such that $a_{\star}(u, v) = (f, v)_{L^2(\Omega_R)}$ for all $v \in \mathcal{H}$

satisfies

$$\|u\|_{H^{1}_{R}(\Omega_{R})} \leq \frac{R}{C_{\operatorname{coer}\star}} \|f\|_{L^{2}(\Omega_{R})} \quad and \quad |u|_{H^{2}(\Omega_{R})} \leq C_{H^{2}\star} \|f\|_{L^{2}(\Omega_{R})}, \quad (7.5)$$

where

$$C_{H^{2}\star} := C_{H^{2}} \left(1 + \sqrt{2} (C_{\text{coer}\star})^{-1} \right).$$

Proof Since $a_{\star}(\cdot, \cdot)$ is continuous and coercive in \mathcal{H} , the first bound in (7.5) follows from the Lax-Milgram theorem and the fact that

$$\sup_{v\in\mathcal{H}}\frac{\left|(f,v)_{L^{2}(\Omega_{R})}\right|}{\|v\|_{H^{1}_{R}(\Omega_{R})}}\leq R\,\|f\|_{L^{2}(\Omega_{R})}\,,$$

by the definition of $\|\cdot\|_{H^1_R(\Omega_R)}$ (7.3). The second bound in (7.5) follows from combining the first bound in (7.5) and the bound (3.9).

We now define the particular Galerkin projection known in the literature as the "elliptic projection" (see the discussion in §4.2).

Definition 7.3 (*Elliptic projection* \mathcal{P}_h) Given $u \in \mathcal{H}$, define $\mathcal{P}_h u \in \mathcal{H}_h$ by

$$a_{\star}(v_h, \mathcal{P}_h u) = a_{\star}(v_h, u) \text{ for all } v_h \in \mathcal{H}_h.$$

Since $a_{\star}(\cdot, \cdot)$ is continuous and coercive in \mathcal{H} by Lemma 7.1, the Lax–Milgram theorem implies that \mathcal{P}_h is well defined. The definition of \mathcal{P}_h then immediately implies the Galerkin-orthogonality property that

$$a_{\star}(v_h, u - \mathcal{P}_h u) = 0 \quad \text{for all } v_h \in \mathcal{H}_h.$$
(7.6)

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Lemma 7.4 (Approximation properties of \mathcal{P}_h) The elliptic projection \mathcal{P}_h satisfies

$$\|u - \mathcal{P}_h u\|_{\star} \le \sqrt{C_{\text{cont}\star}} \min_{v_h \in \mathcal{H}_h} \|u - v_h\|_{H^1_k(\Omega_R)} \quad and \quad (7.7)$$

$$\|u - \mathcal{P}_h u\|_{L^2(\Omega_R)} \le h\sqrt{2}C_{\mathrm{int}}C_{H^2} \sqrt{C_{\mathrm{cont}\star}} \|u - \mathcal{P}_h u\|_{\star}$$
(7.8)

for all $u \in \mathcal{H}$.

Proof By the Cauchy–Schwarz inequality $a_{\star}(\cdot, \cdot)$ is continuous in the $\|\cdot\|_{\star}$ norm, and by definition, $a_{\star}(\cdot, \cdot)$ is coercive in this norm. Therefore Céa's lemma implies that

$$\|u-\mathcal{P}_h u\|_{\star} \leq \min_{v_h\in\mathcal{H}_h} \|u-v_h\|_{\star},$$

and (7.7) follows from the norm equivalence (7.4).

To prove (7.8) we use the standard duality argument. Given $u \in \mathcal{H}$, let ξ be the solution of the variational problem

find
$$\xi \in \mathcal{H}$$
 such that $a_{\star}(\xi, v) = (u - \mathcal{P}_h u, v)_{L^2(\Omega_R)}$ for all $v \in \mathcal{H}$. (7.9)

Then, by Galerkin orthogonality (7.6) and continuity of $a_{\star}(\cdot, \cdot)$, for all $v_h \in \mathcal{H}_h$,

$$\|u - \mathcal{P}_{h}u\|_{L^{2}(\Omega_{R})}^{2} = a_{\star}(\xi, u - \mathcal{P}_{h}u) = a_{\star}(\xi - v_{h}, u - \mathcal{P}_{h}u)$$

$$\leq \|\xi - v_{h}\|_{\star} \|u - \mathcal{P}_{h}u\|_{\star}.$$
(7.10)

By the norm equivalence (7.4), the consequence (3.11) of the definition of C_{int} , the definition of ξ (7.9), and the second bound in (7.5),

$$\begin{split} \|\xi - I_h \xi\|_{\star} &\leq \sqrt{C_{\text{cont}\star}} \, \|\xi - I_h \xi\|_{H^1_k(\Omega_R)} \leq \sqrt{C_{\text{cont}\star}} \sqrt{2} C_{\text{int}} h |\xi|_{H^2(\Omega_R)}, \\ &\leq \sqrt{C_{\text{cont}\star}} \sqrt{2} C_{\text{int}} h C_{H^2\star} \, \|u - \mathcal{P}_h u\|_{L^2(\Omega_R)}, \end{split}$$

and the result (7.8) follows from combining this last inequality with (7.10).

8 Adjoint approximability

Definition 8.1 (Adjoint solution operator S^*) Given $f \in L^2(\Omega_R)$, let $S^* f$ be defined as the solution of the variational problem

find
$$\mathcal{S}^* f \in \mathcal{H}$$
 such that $a(v, \mathcal{S}^* f) = (v, f)_{L^2(\Omega_R)}$ for all $v \in \mathcal{H}$. (8.1)

 S^* is therefore the solution operator of the adjoint problem to the variational problem (2.7) with data in $L^2(\Omega_R)$.

Green's second identity applied to outgoing solutions of the Helmholtz equation implies that $\langle \text{DtN}_k \psi, \overline{\phi} \rangle_{\Gamma_R} = \langle \text{DtN}_k \phi, \overline{\psi} \rangle_{\Gamma_R}$ (see, e.g., [92, Lemma 6.13]); thus $a(\overline{v}, u) = a(\overline{u}, v)$ and so the definition (8.1) implies that

$$a(\overline{\mathcal{S}^*f}, v) = (\overline{f}, v)_{L^2(\Omega_R)} \quad \text{for all } v \in \mathcal{H};$$
(8.2)

i.e. S^*f is the complex-conjugate of an outgoing Helmholtz solution.

Following [88], we define the quantity $\eta(\mathcal{H}_h)$ by

$$\eta(\mathcal{H}_h) := \sup_{f \in L^2(\Omega_R)} \min_{v_h \in \mathcal{H}_h} \frac{\|\mathcal{S}^* f - v_h\|_{H^1_k(\Omega_R)}}{\|f\|_{L^2(\Omega_R)}};$$
(8.3)

observe that this definition implies that, given $f \in L^2(\Omega_R)$,

there exists
$$w_h \in \mathcal{H}_H$$
 such that $\|\mathcal{S}^*f - w_h\|_{H^1_k(\Omega_R)} \le \eta(\mathcal{H}_h) \|f\|_{L^2(\Omega_R)}$. (8.4)

Lemma 8.2 Assume that A, n, and Ω_{-} are nontrapping (and so (3.5) holds with C_{sol} independent of k).

(i) If $\Gamma \in C^{1,1}$ and $A \in C^{0,1}$, then

$$\eta(\mathcal{H}_h) \le hk \left[\sqrt{2}C_{\text{int}}C_{H^2}C_{\text{sol}}R\left(n_{\max} + \frac{1}{k_0R_0C_{\text{sol}}} + 2\right) \right].$$
(8.5)

(ii) If Ω_{-} is a Dirichlet obstacle (so that $\mathcal{H} = H^{1}_{0,D}(\Omega_{R})$), Γ is analytic, A = I, n = 1, and the triangulation \mathcal{T}_{h} in the definition of \mathcal{H}_{h} (2.10) satisfies the quasiuniformity assumption [81, Assumption 5.1], then there exists $C_{\rm MS} = C_{\rm MS}(\Omega_{-})$ such that

$$\eta(\mathcal{H}_h) \le C_{\rm MS} \left[\frac{h}{p} + C_{\rm sol} R \left(\frac{hk}{p} \right)^p \right].$$
(8.6)

Proof Part (ii) is proved in [81, Lemma 3.4 and Proposition 5.3]: see [81, Proof of Theorem 5.8], and observe that the nontrapping assumption implies that α in [81] equals zero. We now prove Part (i).

By the consequence (3.11) of the definition of C_{int} (3.10), there exists $v_h \in \mathcal{H}_h$ such that

$$\left\|\mathcal{S}^*f - v_h\right\|_{H^1_k(\Omega_R)} \le \sqrt{2}C_{\mathrm{int}}h|\mathcal{S}^*f|_{H^2(\Omega_R)}$$

(indeed, we can take $v_h = I_h(\mathcal{S}^* f)$). By (8.2), the BVP (3.8) is satisfied with $v := \mathcal{S}^* f$ and $\tilde{f} := f + k^2 n \mathcal{S}^* f$. Applying the bounds (3.9) and (3.5), we obtain

$$|\mathcal{S}^*f|_{H^2(\Omega_R)} \le C_{H^2} \left(k^2 n_{\max} \| \mathcal{S}^*f \|_{L^2(\Omega_R)} + \| f \|_{L^2(\Omega_R)} \right)$$

$$+ \frac{1}{R} \left\| \nabla(\mathcal{S}^* f) \right\|_{L^2(\Omega_R)} + \frac{1}{R^2} \left\| \mathcal{S}^* f \right\|_{L^2(\Omega_R)} \right)$$

$$\le C_{H^2} C_{\text{sol}} k R \left(n_{\max} + \frac{1}{kR C_{\text{sol}}} + \frac{1}{kR} + \frac{1}{(kR)^2} \right) \| f \|_{L^2(\Omega_R)} ,$$

and the result (8.5) follows from the assumption that $kR \ge k_0R_0 \ge 1$ (see (3.1)). \Box

9 Proof of the oscillatory-behaviour bound (3.6)

Theorem 9.1 If A, n, and Ω_{-} are nontrapping (in the sense that the bound (3.5) holds), then the bound (3.6) holds, i.e.,

$$|u|_{H^{2}(\Omega_{R})} \leq C_{\text{osc}} k \|u\|_{H^{1}_{\mu}(\Omega_{R})}.$$
(9.1)

Lemma 9.2 To prove Theorem 9.1, it is sufficient to prove that there exists $k_0 > 0$ and $C_{\text{mass}} = C_{\text{mass}}(A, n, \Omega_-, R) > 0$ such that

$$\|u\|_{L^{2}(\Omega_{R+1})} \le C_{\max} \|u\|_{L^{2}(\Omega_{R})} \quad for all \ k \ge k_{0}.$$
(9.2)

Proof We first claim that the map $k \mapsto u$ is continuous from $(1, \infty)$ to $H^2(\Omega_R)$; indeed, this follows from the well-posedness of the plane-wave scattering problem of Definition 2.2, H^2 regularity, and linearity. Therefore, the function $k \mapsto$ $\|u\|_{H^2(\Omega_R)} (k \|u\|_{H^1_k(\Omega_R)})^{-1}$ is continuous on $[1, \infty)$, and it is sufficient to prove that the bound (9.1) (i.e., (3.6)) holds for k sufficiently large.

Let $\chi \in C^{\infty}(\mathbb{R}^d)$ be such that $0 \le \chi \le 1$, $\chi = 1$ on Ω_R and $\chi = 0$ on $\mathbb{R}^d \setminus B_{R+1/2}$. Applying the H^2 -regularity results [62, Theorems 2.3.3.2, 2.4.2.5, and 2.4.2.7] to χu (with these results valid since A is Lipschitz, $A_{\min} > 0$, both Γ and Γ_R are $C^{1,1}$, and either $\gamma u = 0$ or $\partial_{\mathbf{n}} u = 0$ on Γ), we obtain, in a similar way to the proof of Theorem 6.1, that there exists $C_1 = C_1(A, n, \Omega_-, R) > 0$, such that

$$|u|_{H^2(\Omega_R)} \le C_1 k ||u||_{H^1_k(\Omega_{R+1})}.$$

Therefore to prove (9.1) (i.e., (3.6)), it is sufficient to prove that there exists $C_2 = C_2(A, n, \Omega_-, R) > 0$, such that

$$\|u\|_{H^{1}_{k}(\Omega_{R+1})} \le C_{2} \|u\|_{H^{1}_{k}(\Omega_{R})}.$$
(9.3)

We now need to show that we can prove (9.3) from (9.2). We claim that

$$\|\nabla u\|_{L^{2}(\Omega_{R+1})} \leq \sqrt{\frac{n_{\max}}{A_{\min}}} k \|u\|_{L^{2}(\Omega_{R+1})} \quad \text{for all } k > 0.$$
(9.4)

Indeed, applying Green's identity in Ω_R (which is justified by [78, Theorem 4.4] since $u \in H^1(\Omega_R)$) and recalling that either $\gamma u = 0$ or $\partial_n u = 0$ on Γ , we have that

$$\int_{\Omega_{R+1}} (\mathsf{A}\nabla u) \cdot \overline{\nabla u} - k^2 n |u|^2 = \Re \int_{\Gamma_{R+1}} \overline{u} \frac{\partial u}{\partial r}.$$

By (3.4), the right-hand side is ≤ 0 , and (9.4) follows using the inequalities (2.1) and (2.2). Therefore, using (9.4) and (9.2),

$$\|u\|_{H^{1}_{k}(\mathcal{Q}_{R+1})} \leq \sqrt{\frac{n_{\max}}{A_{\min}} + 1} \ k \ \|u\|_{L^{2}(\mathcal{Q}_{R+1})} \leq C_{\max} \sqrt{\frac{n_{\max}}{A_{\min}} + 1} \ k \ \|u\|_{L^{2}(\mathcal{Q}_{R})}$$

which implies the bound (9.3), and the result follows.

9.1 Overview of the ideas used in the rest of this section to prove (9.2)

We have therefore reduced proving the oscillatory-behaviour bound (3.6)/(9.1) to proving the bound (9.2), which we prove using *defect measures*. The precise definition of a defect measure is given in Theorem 9.3 below, but the idea is that the defect measure of a Helmholtz solution describes where the mass of the solution in phase space (i.e. the set of positions **x** and momenta $\boldsymbol{\xi}$) is concentrated in the high-frequency limit. Two examples of this feature are

- (i) the defect measure of the plane wave u¹(**x**) := exp(ik**x** · **a**) is the product of a delta function at **ξ** = **a** and Lebesgue measure in **x** (see (9.8) below), reflecting the fact that, at high frequency (and in fact at any frequency), all the mass in phase space of the plane wave is travelling in the direction **a**, and
- (ii) the defect measure of an outgoing solution of the Helmholtz equation is zero on the so-called "directly incoming set" (see Lemma 9.8 below), where this set is defined in (9.20) below as points in phase space that don't hit the scatterer when propagated backwards along the flow.

A key feature of the defect measure of a Helmholtz solution is that it is invariant under the Hamiltonian flow defined by the symbol of the PDE, as long as the flow doesn't encounter the scatterer (see Theorem 9.6 below) This is analogous to results about propagation of singularities of the wave equation, where singularities travel along the trajectories of the flow (the *bicharacteristics*), and the projection of these trajectories in space are the *rays*.

The main ingredients to our proof of (9.2) are Points (i) and (ii) above, invariance under the flow (away from the scatterer), and then geometric arguments about the rays, using the fact that away from the scatterer the rays are straight lines and the flow has constant speed along the rays (see (9.12) below).

To conclude this overview, we direct the reader to [105, Chapter 5] for extensive discussion of defect measures in \mathbb{R}^d , to [16,50,84] for material on defect measures on manifolds with boundary, and to [15] for discussion on the history of defect measures.

9.2 Recap of results about defect measures

9.2.1 Symbols and quantisation

Before defining defect measures, we need to define the functions on phase space (i.e. the set of positions **x** and momenta $\boldsymbol{\xi}$) that the defect measure can act upon by dual pairing. These functions are called *symbols*, defined as functions on the *cotangent* bundle $T^*\Omega_+$. Recall the definition of the cotangent bundle of \mathbb{R}^d :

 $T^* \mathbb{R}^d := \mathbb{R}^d \times (\mathbb{R}^d)^*;$

for our purposes, we can consider $T^*\mathbb{R}^d$ as $\{(\mathbf{x}, \boldsymbol{\xi}) : \mathbf{x} \in \mathbb{R}^d, \boldsymbol{\xi} \in \mathbb{R}^d\}$, i.e. the set of positions \mathbf{x} and momenta $\boldsymbol{\xi}$. On $T^*\mathbb{R}^d$, the *quantisation* of a symbol $b(\mathbf{x}, \boldsymbol{\xi}) \in C^{\infty}_{\text{comp}}(T^*\mathbb{R}^d)$ is defined by

$$b(\mathbf{x}, (\mathbf{i}k)^{-1}\partial_{\mathbf{x}})u(\mathbf{x}) := \frac{k^d}{(2\pi)^d} \int_{\mathbb{R}^d} \int_{\mathbb{R}^d} e^{\mathbf{i}k(\mathbf{x}-\mathbf{y})\cdot\boldsymbol{\xi}} b(\mathbf{x}, \boldsymbol{\xi})u(\mathbf{y}) \,\mathrm{d}\mathbf{y} \,\mathrm{d}\boldsymbol{\xi}; \qquad (9.5)$$

see, e.g., [105, §4]. The same definition holds for symbols supported away from the boundary of $\overline{\Omega}_+$. We omit the analogous definition near the boundary since it is more involved; see [16, §4.2] (where it involves the so-called *compressed cotangent bundle* of Ω_+ , $T_b^*\overline{\Omega_+}$) and [84, §1.2]. We will not, in any event, require any specifics of the measure at the boundary in proving Theorem 9.1.

9.2.2 Existence of defect measures

Theorem 9.3 (Existence of defect measures [105, Theorem 5.2], [16, §4.2]) Suppose $\{v(k)\}_{k_0 \le k < \infty}$ is a collection of functions that is uniformly locally bounded in $L^2(\Omega_+)$, *i.e. given* $\chi \in C^{\infty}_{\text{comp}}(\mathbb{R}^d)$ there exists C > 0, depending on χ and k_0 but independent of k, such that

$$\|\chi v(k)\|_{L^2(\Omega_+)} \le C \quad \text{for all } k \ge k_0.$$
 (9.6)

Then there exists a sequence $k_{\ell} \to \infty$ and a non-negative Radon measure μ on $T_{b}^{*}\overline{\Omega_{+}}$ (depending on k_{ℓ}) such that, for any symbol $b(\mathbf{x}, \boldsymbol{\xi}) \in C_{\text{comp}}^{\infty}(T_{b}^{*}\overline{\Omega_{+}})$

$$\left\langle b\left(\mathbf{x}, (\mathbf{i}k_{\ell})^{-1}\partial_{\mathbf{x}}\right)v(k_{\ell}), v(k_{\ell})\right\rangle_{\Omega_{+}} \longrightarrow \int b \, \mathrm{d}\mu \quad as \, \ell \to \infty.$$
 (9.7)

In the case of a plane wave $u^{I}(\mathbf{x}) := \exp(ik\mathbf{x} \cdot \mathbf{a})$ with $|\mathbf{a}| = 1$, a direct calculation using (9.5) and the definition of the Fourier transform shows that, for all k,

$$\langle b \, u^{I}, u^{I} \rangle_{\mathbb{R}^{d}} := \frac{k^{d}}{(2\pi)^{d}} \int_{\mathbb{R}^{d}} \int_{\mathbb{R}^{d}} \int_{\mathbb{R}^{d}} e^{ik(\mathbf{x}-\mathbf{y})\cdot\boldsymbol{\xi}} e^{ik\mathbf{y}\cdot\mathbf{a}} e^{-ik\mathbf{x}\cdot\mathbf{a}} b(\mathbf{x},\boldsymbol{\xi}) \,\mathrm{d}\boldsymbol{\xi} \,\mathrm{d}\mathbf{y} \,\mathrm{d}\mathbf{x}$$
$$= \int_{\mathbb{R}^{d}} b(\mathbf{x},\mathbf{a}) \,\mathrm{d}\mathbf{x};$$
(9.8)

i.e. for any sequence $k_{\ell} \to \infty$, the corresponding defect measure of u^{I} is the product of the Lebesgue measure in **x** by a delta measure at $\boldsymbol{\xi} = \mathbf{a}$; we therefore talk about *the* (as opposed to *a*) defect measure of u^{I} .

The next lemma proves that, if *u* is the solution of the plane-wave scattering problem and χ is an arbitrary cut-off function, then χu is uniformly bounded in *k* (on compact subsets of Ω_+); existence of a defect measure of *u* then follows from Theorem 9.3. In the rest of this section, to emphasise the *k*-dependence of *u*, we write u = u(k).

Lemma 9.4 Let u(k) be the solution of the plane-wave scattering problem of Definition 2.2. Assume that A, n, and Ω_{-} are nontrapping. Then there exists $C(A, n, \Omega_{-}, R, k_0) > 0$ such that

$$\|u(k)\|_{L^2(\Omega_P)} \le C \quad \text{for all } k \ge k_0.$$
 (9.9)

Proof Let $\chi \in C_{\text{comp}}^{\infty}(\mathbb{R}^d)$ be such that $\chi = 1$ in a neighbourhood of the scatterer Ω_{sc} . Let $v := u^S + \chi u^I$, so that $u = (1 - \chi)u^I + v$. Since $||u^I(k)||_{L^2(\Omega_R)} \leq C_1(R)$ for all k > 0, the result (9.9) will follow if we prove a uniform bound on $||v(k)||_{L^2(\Omega_R)}$. The definition of v implies that v satisfies the Sommerfeld radiation condition, either $\gamma v = 0$ or $\partial_{\mathbf{n}} v = 0$ on Γ , and, with $\mathcal{L}_{\mathsf{A},n} w := \nabla \cdot (\mathsf{A} \nabla w) + k^2 n w$ and [A, B] := AB - BA,

$$\mathcal{L}_{\mathsf{A},n}v = -\mathcal{L}_{\mathsf{A},n}\big((1-\chi)u^{I}\big) = \big[\mathcal{L}_{\mathsf{A},n},\chi\big]u^{I} - (1-\chi)\mathcal{L}_{\mathsf{A},n}u^{I} = \big[\mathcal{L}_{\mathsf{A},n},\chi\big]u^{I},$$

since $\mathcal{L}_{A,n}u^I = 0$ when $1 - \chi \neq 0$. By explicit calculation, using the fact that $u^I(\mathbf{x}) = \exp(ik\mathbf{x} \cdot \mathbf{a})$,

$$\left\| \left[\mathcal{L}_{\mathsf{A},n}, \chi \right] u^{I} \right\|_{L^{2}(\Omega_{R})} \leq C_{1} k,$$

where C_1 depends on $\|A\|_{L^{\infty}(\Omega_R)}$, $\|\nabla A\|_{L^{\infty}(\Omega_R)}$, and χ , but is independent of k. The nontrapping bound (3.5) then implies that $\|v(k)\|_{L^2(\Omega_R)} \leq C_2$ with C_2 independent of k, and the result follows.

9.2.3 Support and invariance properties of defect measures

Recall that the semi-classical principal symbol of the Helmholtz equation (2.3) is given by

$$p(\mathbf{x}, \boldsymbol{\xi}) := \sum_{i=1}^{d} \sum_{j=1}^{d} A_{ij}(\mathbf{x}) \xi_i \xi_j - n(\mathbf{x})$$
(9.10)

(see, e.g., [105, Page 281]). In our arguments below we only consider points $(\mathbf{x}, \boldsymbol{\xi})$ in phase space when p = 0; this is because of the following result.

Theorem 9.5 (Support of defect measure [105, Theorem 5.4], [16, Equation 3.17]) Suppose u(k) satisfies (9.9), and let μ be any defect measure of u(k). Then supp $\mu \subset \{(\mathbf{x}, \boldsymbol{\xi}) : p(\mathbf{x}, \boldsymbol{\xi}) = 0\}$.

As an illustration of this, the plane wave $u^{I}(\mathbf{x}) := \exp(i\mathbf{k}\mathbf{x} \cdot \mathbf{a})$ with $|\mathbf{a}| = 1$ is solution of the Helmholtz equation (2.3) with A = I and n = 1, and hence $p = |\boldsymbol{\xi}|^{2} - 1$ in this case. By (9.8), the defect measure of u^{I} is the product of Lebesgue measure in \mathbf{x} and a delta function at $\boldsymbol{\xi} = \mathbf{a}$, and thus is supported in $|\boldsymbol{\xi}| = 1$, i.e., p = 0, as expected from Theorem 9.5.

The final result about defect measures that we need is their invariance under the flow (away from the scatterer). This result is Theorem 9.6 below; to state it, we first need to define the flow.

Away from Γ , and provided that *A* and *n* are both $C^{1,1}$, the flow φ_t is defined as follows: given $\rho = (\mathbf{x}_0, \xi_0), \varphi_t(\rho) := (\mathbf{x}(t), \boldsymbol{\xi}(t))$ where $(\mathbf{x}(t), \boldsymbol{\xi}(t))$ is the solution of the Hamiltonian system

$$\dot{x}_{i}(t) = \partial_{\xi_{i}} p(\mathbf{x}(t), \boldsymbol{\xi}(t)), \qquad \dot{\xi}_{i}(t) = -\partial_{x_{i}} p(\mathbf{x}(t), \boldsymbol{\xi}(t)), \qquad (9.11)$$

with initial condition $(\mathbf{x}(0), \boldsymbol{\xi}(0)) = (\mathbf{x}_0, \boldsymbol{\xi}_0)$, where the Hamiltonian equals p defined by (9.10). Near both Γ and places where A and n are not $C^{1,1}$, the definition of φ_t is more involved – this is to account for reflection or refraction. However, we do not need this definition in what follows, since our arguments take place away from these regions. In fact our arguments take place away from the scatterer Ω_{sc} . Outside Ω_{sc} , A = I, and n = 1; thus $p(\mathbf{x}, \boldsymbol{\xi}) = |\boldsymbol{\xi}|^2 - 1$. From (9.11), the flow satisfies $\dot{x}_i = 2\xi_i$ and $\dot{\xi}_i = 0$ and is therefore given by the straight-line motion

$$\mathbf{x} = \mathbf{x}_0 + 2t\boldsymbol{\xi}_0, \quad \boldsymbol{\xi} = \boldsymbol{\xi}_0.$$
 (9.12)

The arguments below consider the flow with speed 2 (i.e. with $|\xi_0| = 1$). This is without loss of generality, since away from Ω_{sc} Theorem 9.5 implies that μ is only non-zero when $|\xi| = 1$.

Both in the next result and later, we let π_x denote projection in the x variables, i.e. $\pi_x((x, \xi)) = x$.

Theorem 9.6 (Invariance of defect measure under the flow away from the scatterer) Suppose that u(k) satisfies (9.9), and let μ be any defect measure of u(k). If $A \subset T^* \mathbb{R}^d$ is such that $\pi_{\mathbf{x}}(\varphi_s(A)) \cap \Omega_{sc} = \emptyset$ for s between 0 and t, (i.e. the flow acting on A doesn't hit the scatterer from time 0 to time t), then

$$\mu(\varphi_t(A)) = \mu(A). \tag{9.13}$$

Proof In the absence of the scatterer, invariance of the measure under the flow is the statement that, for $b \in C^{\infty}_{\text{comp}}(T^*\mathbb{R}^d)$,

$$\partial_s \left(\int (b \circ \varphi_{-s})(\rho) \, \mathrm{d}\mu \right) = 0 \quad \text{for all } s,$$
 (9.14)

and this is proved in [105, Theorem 5.4], [16, Proposition 4.4]. For this result to hold in the presence of the scatterer in a time interval $0 \le s \le t$, we need the spatial projection of the integrand in (9.14) to not be supported during this time interval on Ω_{sc} , i.e., we need the condition that

$$\pi_{\mathbf{x}}(\operatorname{supp}(b \circ \varphi_{-s})) \cap \Omega_{\operatorname{sc}} = \emptyset \quad \text{for } 0 \le s \le t.$$
(9.15)

Under this condition, (9.14) implies that

$$\int b(\rho) \,\mathrm{d}\mu = \int (b \circ \varphi_{-s})(\rho) \,\mathrm{d}\mu \quad \text{for all } 0 \le s \le t.$$
(9.16)

Let 1_A denote the indicator function of a set *A*. By approximating 1_A by smooth symbols, (9.16) holds with $b(\rho) = 1_A(\rho)$, provided that the condition (9.15) holds. Since $\varphi_{-s}(\rho) \in A$ iff $\rho \in \varphi_s(A)$, we have

$$\pi_{\mathbf{x}}(\operatorname{supp}(1_{A} \circ \varphi_{-s})) = \pi_{\mathbf{x}}(\operatorname{supp}(1_{\varphi_{s}(A)})) = \pi_{\mathbf{x}}(\varphi_{s}(A)),$$

and thus (9.15) holds by the assumption in the statement of the theorem.

Therefore, (9.16) implies that, for all $0 \le s \le t$,

$$\int \mathbf{1}_A(\rho) \, \mathrm{d}\mu = \int \mathbf{1}_A(\varphi_{-s}(\rho)) \, \mathrm{d}\mu = \int \mathbf{1}_{\varphi_s(A)}(\rho) \, \mathrm{d}\mu,$$

i.e.

$$\mu(A) = \mu(\varphi_s(A)) \quad \text{for all } 0 \le s \le t,$$

which implies (9.13).

9.3 Proof of (9.2) using defect measures

The following lemma reduces proving the bound (9.2) to proving a statement about defect measures.

Lemma 9.7 Let $0 < R_0 < R$ be such that $\Omega_{sc} \subset B_{R_0}$. If every defect measure of u is non-zero and there exists $C_{R,R_0} > 0$ such that, for every defect measure μ of u,

$$\mu(T^*\Omega_{R+2}) \le C_{R,R_0}\mu(T^*\Omega_{R_0}),\tag{9.17}$$

then the bound (9.2) holds.

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Proof We prove the contrapositive. Suppose (9.2) fails; we aim to exhibit a defect measure associated to *u* for which (9.17) fails. Then, for any $C_1 > 0$, there exists a sequence $(k_n)_{n=1}^{\infty}$, with $k_n \to \infty$, such that

$$\|u(k_n)\|_{L^2(\Omega_{R+1})} \ge C_1 \|u(k_n)\|_{L^2(\Omega_R)};$$
(9.18)

we choose $C_1 := 2C_{R,R_0}$. By Lemma 9.4, the sequence $\{u(k_n)\}_{n=1}^{\infty}$ is locally uniformly bounded and Theorem 9.3 implies that, by passing to a subsequence, there exists a defect measure μ of u associated to the subsequence, which we again denote k_n . Let $\chi_0, \chi_1 \in C^{\infty}(\mathbb{R}^d)$ be such that $0 \le \chi_0, \chi_1 \le 1$, and

$$\operatorname{supp}\chi_1 \subset B_{R+2}, \quad \chi_1 = 1 \text{ in } B_{R+1}, \quad \operatorname{supp}\chi_0 \subset B_R, \quad \chi_0 = 1 \text{ in } B_{R_0}.$$

The bound (9.18) then implies that

$$\|\chi_1 u(k_n)\|_{L^2(\Omega_+)} \ge 2C_{R,R_0} \|\chi_0 u(k_n)\|_{L^2(\Omega_+)}.$$
(9.19)

Passing to the limit $n \to \infty$ and using the property of defect measure (9.7), we obtain that

$$\int \chi_1^2 \,\mathrm{d}\mu \geq 2C_{R,R_0} \int \chi_0^2 \,\mathrm{d}\mu.$$

The definitions of χ_0 and χ_1 imply that

$$\int \chi_0^2 \,\mathrm{d}\mu \ge \int \mathbf{1}_{T^* \Omega_{R_0}} \,\mathrm{d}\mu = \mu(T^* \Omega_{R_0})$$

and

$$\int \chi_1^2 \,\mathrm{d}\mu \leq \int \mathbb{1}_{T^* \mathfrak{Q}_{R+2}} \,\mathrm{d}\mu = \mu(T^* \mathfrak{Q}_{R+2});$$

hence

$$\mu(T^*\Omega_{R+2}) \ge 2C_{R,R_0}\mu(T^*\Omega_{R_0}),$$

contradicting (9.17).

Before using Lemma 9.7 to prove (9.2), we prove a result (Lemma 9.8 below) about the structure of μ , exploiting the fact that $u = u^{I} + u^{S}$ with u^{S} is outgoing (in the sense that it satisfies the Sommerfeld radiation condition (2.4)). To make use of this outgoing property, we need to define appropriate notions of incoming and outgoing for elements of phase space. Let \mathcal{I} denote the *directly incoming set* defined by

$$\mathcal{I} := \left\{ \rho \in T^*(\Omega_+ \backslash \Omega_{\mathrm{sc}}), \text{ s.t. } \pi_{\mathbf{x}} \left(\bigcup_{t \ge 0} \varphi_{-t}(\rho) \right) \cap \Omega_{\mathrm{sc}} = \emptyset \right\};$$
(9.20)

where recall that π_x denotes projection in the x variables. That is, \mathcal{I} is everything that never hits the scatterer under backward flow. Let

$$\Gamma_+ := (T^* \Omega_+) \backslash \mathcal{I}.$$

These definitions of \mathcal{I} and Γ_+ do not require the generalized bicharacteristic flow φ_t to be defined in $T^*\Omega_{sc}$, but when the flow is defined everywhere, Γ_+ is the forward generalized bicharacteristic flowout of Ω_{sc} , that is

$$\Gamma_{+} = \left\{ \bigcup_{t \ge 0} \varphi_{t}(\rho) : \rho \in T^{*} \Omega_{sc} \right\} \text{ when } \varphi_{t} \text{ is defined everywhere.}$$

The following lemma uses outgoingness of u^S to show that, given a set E in phase space, the mass of u lying over E is *either* in the forward flowout Γ_+ or associated to the incident wave u^I .

Lemma 9.8 For any Borel set $E \subset T^*\Omega$, $\mu(E \setminus \Gamma_+) = \mu^I(E \setminus \Gamma_+)$, where μ is any defect measure of u, and μ^I is the defect measure of u^I .

Proof Let k_{ℓ} be the sequence associated to the particular defect measure of u. By Lemma 9.4, $u^{S}(k_{\ell})$ is uniformly locally bounded, and so there exists a subsequence k_{ℓ_m} and a defect measure associated to u^{S} , denoted by μ^{S} . Then, by linearity and (9.7), $\mu = \mu^{S} + \mu^{I}$. It is therefore sufficient to prove that $\mu^{S}(E \setminus \Gamma_{+}) = 0$. But, by the definition of Γ_{+} , $E \setminus \Gamma_{+} \subset \mathcal{I}$, and $\mu^{S}(\mathcal{I}) = 0$ by [16, Proposition 3.5], [50, Lemma 3.4], since u^{S} is outgoing.

Proof of Theorem 9.1 By Lemmas 9.2 and 9.7 it is sufficient to prove the bound (9.17) (observe that the hypothesis in Lemma 9.7 that every defect measure of *u* is non-zero holds by Lemma 9.8 since $\mu^{I}(\mathcal{I}) \neq 0$). Let $R_{sc} := \max_{\mathbf{x} \in \Omega_{sc}} |\mathbf{x}|$. We claim that it is sufficient to show that, for any $\rho > R_{sc}$ there exists $\varepsilon = \varepsilon(R_{sc}, \rho)$, with $\varepsilon(R_{sc}, \rho)$ is an increasing function of ρ , and $C = C(\rho, \varepsilon) > 0$ such that

$$\mu(T^*(B_{\rho+\varepsilon} \setminus B_{\rho})) \le C(\rho, \varepsilon)\mu(T^*\Omega_{\rho}).$$
(9.21)

Indeed, we now show that the bound (9.17) then follows by using (9.21) repeatedly. Since $\varepsilon(R_{sc}, \rho)$ is an increasing function of ρ , if $\varepsilon^* := \varepsilon(R_{sc}, R_0)$, then (9.21) implies, with $C(\rho) := C(\rho, \varepsilon(R_{sc}, \rho))$,

$$\mu(T^*(B_{\rho+\varepsilon^*} \setminus B_{\rho})) \le C(\rho)\,\mu(T^*\Omega_{\rho}) \quad \text{for all } \rho \ge R_0. \tag{9.22}$$

The bound (9.17) then follows by applying (9.22) with $\rho = R_0$, $\rho = R_0 + \varepsilon^*$, ..., $\rho = R_0 + m\varepsilon^*$, where $m = \lceil (R + 2 - R_0)/\varepsilon^* \rceil$.

It is therefore sufficient to prove the bound (9.21); we introduce the notation that $A := B_{\rho+\varepsilon} \setminus B_{\rho}$, and observe that (9.21) then reads $\mu(T^*A) \leq C(\rho, \varepsilon)\mu(T^*\Omega_{\rho})$. We prove this bound by combining the following three inequalities:

$$\mu(T^*A) \le \mu(T^*A \cap \Gamma_+) + \mu_I(T^*A) = \mu(T^*A \cap \Gamma_+) + |A|$$
(9.23)

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(where $|\cdot|$ denotes Lebesgue measure in \mathbb{R}^d),

$$\mu(T^*A \cap \Gamma_+) \le \mu(T^*(B_\rho \setminus B_{\rho_0})) \le \mu(T^*\Omega_\rho), \tag{9.24}$$

where $\rho_0 := (\rho + R_{sc})/2$, and

$$\mu(T^*\Omega_\rho) \ge \delta|\Omega_\rho| \tag{9.25}$$

for some $\delta > 0$. Indeed, using (9.23), (9.24), and (9.25), we have

$$\mu(T^*A) \le \left(1 + |A|(\delta|\Omega_{\rho}|)^{-1}\right) \mu(T^*\Omega_{\rho}),$$

which is (9.21). We prove (9.23) and (9.25) using Lemma 9.8 and the structure of μ^{I} , and (9.24) using invariance of defect measures under the flow outside of $T^{*}\Omega_{sc}$ (i.e. Theorem 9.6).

Proof of (9.23) Lemma 9.8 implies that

$$\mu(T^*A) = \mu(T^*A \cap \Gamma_+) + \mu(T^*A \setminus \Gamma_+) \le \mu(T^*A \cap \Gamma_+) + \mu_I(T^*A).$$

By (9.8), μ_I is a δ -measure on $\boldsymbol{\xi} = \mathbf{a}$ times Lebesgue measure in \mathbf{x} , so $\mu_I(T^*A) = |A|$, (where $|\cdot|$ denotes Lebesgue measure in \mathbb{R}^d) and (9.23) follows.

Proof of (9.24) Recall that, for $X \subset \mathbb{R}^d \setminus \overline{\Omega_{sc}}$,

$$S^*X := \{ (\mathbf{x}, \boldsymbol{\xi}) : \mathbf{x} \in X, \, \boldsymbol{\xi} \in \mathbb{R}^d \text{ with } |\boldsymbol{\xi}| = 1 \},\$$

and observe that, by Theorem 9.5, $\mu(T^*A \cap \Gamma_+) = \mu(S^*A \cap \Gamma_+)$ and $\mu(T^*(B_\rho \setminus B_{\rho_0})) = \mu(S^*(B_\rho \setminus B_{\rho_0}))$; we therefore only need to prove that

$$\mu(S^*A \cap \Gamma_+) \le \mu(S^*(B_\rho \setminus B_{\rho_0})). \tag{9.26}$$

We first introduce some notation that allows us to bound $\mu(S^*A \cap \Gamma_+)$ using only the invariance of defect measure (9.13) in the exterior of Ω_{sc} . Given $\mathbf{b} \in \mathbb{R}^d$ with $|\mathbf{b}| = 1$ and $\tilde{\rho} > R_{sc}$, let $\Omega_{sc,\tilde{\rho},\mathbf{b}} \subset \mathbb{R}^d$ and $\Lambda_{sc,\tilde{\rho},\mathbf{b}} \subset S^*\Omega_+$ be defined by

$$\Omega_{\mathrm{sc},\widetilde{\rho},\mathbf{b}} := \left(\bigcup_{t\geq 0} \left(\Omega_{\mathrm{sc}} + t\mathbf{b}\right)\right) \cap \Omega_{\widetilde{\rho}} \quad \text{and} \quad \Lambda_{\mathrm{sc},\rho+\varepsilon,\mathbf{b}} := \Omega_{\mathrm{sc},\widetilde{\rho},\mathbf{b}} \times \{\mathbf{b}\};$$

i.e. $\Omega_{sc,\tilde{\rho},\mathbf{b}}$ equals the union of all possible translations of Ω_{sc} in the direction **b**, intersected with $\Omega_{\tilde{\rho}}$, and $\Lambda_{sc,\tilde{\rho},\mathbf{b}}$ equals these points paired with the direction **b**. By (9.12), the spatial projections of the flow outside Ω_{sc} are straight lines, and thus

$$\Gamma_{+} \cap S^{*} \Omega_{\widetilde{\rho}} \cap \{ \boldsymbol{\xi} = \mathbf{b} \} = \Big\{ (\mathbf{x}, \mathbf{b}) \in S^{*} \Omega_{\widetilde{\rho}} : \exists s \geq 0 \text{ s.t. } \mathbf{x} - s\mathbf{b} \in \Omega_{sc} \Big\}.$$

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Therefore

$$\Gamma_{+} \cap S^{*} \Omega_{\widetilde{\rho}} \cap \{ \boldsymbol{\xi} = \mathbf{b} \} \subset \Lambda_{\mathrm{sc}, \widetilde{\rho}, \mathbf{b}}, \qquad \Gamma_{+} \cap S^{*} \Omega_{\widetilde{\rho}} \subset \bigcup_{\mathbf{b} \in \mathbb{R}^{d}, |\mathbf{b}|=1} \Lambda_{\mathrm{sc}, \widetilde{\rho}, \mathbf{b}}, \quad (9.27)$$

and thus, for any $\varepsilon > 0$,

$$S^*A \cap \Gamma_+ = S^*A \cap S^*\Omega_{\rho+\varepsilon} \cap \Gamma_+ \subset S^*A \cap \left(\bigcup_{\mathbf{b} \in \mathbb{R}^d, |\mathbf{b}|=1} \Lambda_{\mathrm{sc},\rho+\varepsilon,\mathbf{b}}\right).$$
(9.28)

Recall that $\rho_0 := (\rho + R_{sc})/2$. Let

$$t_0 := \frac{\rho_0 - R_{\rm sc}}{4} = \frac{\rho - R_{\rm sc}}{8} \tag{9.29}$$

and

$$\varepsilon := -\rho + \sqrt{R_{\rm sc}^2 + \left(\frac{\rho - R_{\rm sc}}{4} + \sqrt{\rho^2 - R_{\rm sc}^2}\right)^2}; \tag{9.30}$$

observe that $\varepsilon > 0$ and ε is an increasing function of ρ , as claimed underneath (9.21). We now claim that, with these definitions of t_0 and ε ,

$$\bigcup_{0 \le t \le t_0} \varphi_t \left(S^*(B_\rho \setminus B_{\rho_0}) \right) \cap \Omega_{\rm sc} = \emptyset$$
(9.31)

(i.e., the forward flowout of the annulus $B_{\rho} \setminus B_{\rho_0}$ does not hit the scatterer for $0 \le t \le t_0$) and

$$S^*A \cap \left(\bigcup_{\mathbf{b} \in \mathbb{R}^d, |\mathbf{b}|=1} \Lambda_{\mathrm{sc},\rho+\varepsilon,\mathbf{b}}\right) \subset \varphi_{t_0} \left(S^*(B_{\rho} \setminus B_{\rho_0})\right).$$
(9.32)

(Since $S^*A \cap \Gamma_+$ is contained in the left-hand side of (9.32) by (9.28), (9.32) says that the forward flowout of $B_\rho \setminus B_{\rho_0}$ in time t_0 covers all points in S^*A that are ever reached by flowout from $T^*\Omega_{sc}$.) Outside Ω_{sc} the flow has speed 2 and its spatial projections are straight lines. Therefore (9.31) is ensured if $t_0 < (\rho_0 - R_{sc})/2$, which is ensured by (9.29).

We now show that (9.32) holds. Since

$$(\mathbf{x}, \mathbf{b}) = (\mathbf{x} - 2t_0\mathbf{b} + 2t_0\mathbf{b}, \mathbf{b}) = \varphi_{t_0}(\mathbf{x} - 2t_0\mathbf{b}, \mathbf{b}),$$

(9.32) follows from showing that $(\mathbf{x}-2t_0\mathbf{b}, \mathbf{b}) \in S^*(B_{\rho} \setminus B_{\rho_0})$, i.e. $\mathbf{x}-2t_0\mathbf{b} \in B_{\rho} \setminus B_{\rho_0}$, for all (\mathbf{x}, \mathbf{b}) belonging to the left-hand side of (9.32). For such (\mathbf{x}, \mathbf{b}) , by definition,

$$\rho \le |\mathbf{x}| \le \rho + \epsilon, \text{ and } \mathbf{x} - s\mathbf{b} \in \Omega_{sc}$$
 (9.33)



Fig. 2 Figure showing the lengths L_1 and L_2 defined by (9.34)

for some $s \ge 0$. We now claim that for such (\mathbf{x}, \mathbf{b}) ,

$$\mathbf{x} - \ell \mathbf{b} \in B_{\rho} \setminus B_{\rho_0}$$
 for all $L_1 < \ell \leq L_2$,

where

$$L_1 := \sqrt{(\rho + \varepsilon)^2 - R_{\rm sc}^2} - \sqrt{\rho^2 - R_{\rm sc}^2}, \qquad L_2 := \rho - \rho_0. \tag{9.34}$$

This is because, on the one hand, a ray of length > L_1 starting from a point **x** in a direction $-\mathbf{b}$, with (**x**, **b**) satisfying (9.33), will automatically enter B_{ρ} . Indeed, the longest such ray that does not intersect B_{ρ} has length L_1 , as shown in Fig. 2. On the other hand, a ray of length $\leq L_2$ starting from a point **x** in a direction $-\mathbf{b}$, with (**x**, **b**) satisfying (9.33), will not intersect B_{ρ_0} . Indeed, the shortest such ray that enters $\overline{B_{\rho_0}}$ has length L_2 , as shown in Fig. 2. It is then straightforward to check that $L_1 < 2t_0 \leq L_2$ when t_0 is given by (9.29) and ε is given by (9.30), so that (9.32) holds.

We now prove the bound (9.26) on $\mu(S^*A \cap \Gamma_+)$ using (9.31) and (9.32). Because of (9.31), we can use (9.13) to find that

$$\mu(\varphi_{t_0}(S^*(B_{\rho} \setminus B_{\rho_0}))) = \mu(S^*(B_{\rho} \setminus B_{\rho_0}));$$

using this with (9.28) and (9.32), we obtain (9.26), and thus (9.24).

Proof of (9.25) Using Lemma 9.8 and the structure of μ_I , we have

$$\mu(T^*\Omega_{\rho}) \geq \mu(T^*\Omega_{\rho} \setminus \Gamma_+) = \mu_I(T^*\Omega_{\rho} \setminus \Gamma_+)$$

= $\mu_I((T^*\Omega_{\rho} \setminus \Gamma_+) \cap \{\boldsymbol{\xi} = \mathbf{a}\}) + \mu_I((T^*\Omega_{\rho} \setminus \Gamma_+) \cap \{\boldsymbol{\xi} \neq \mathbf{a}\})$
= $\left|\pi_{\mathbf{x}}((T^*\Omega_{\rho} \setminus \Gamma_+) \cap \{\boldsymbol{\xi} = \mathbf{a}\})\right|.$ (9.35)

Since

$$\pi_{\mathbf{x}}\Big((T^*\mathcal{Q}_{\rho}\setminus\Gamma_+)\cap\{\boldsymbol{\xi}=\mathbf{a}\}\Big)\cup\pi_{\mathbf{x}}\Big((T^*\mathcal{Q}_{\rho}\cap\Gamma_+)\cap\{\boldsymbol{\xi}=\mathbf{a}\}\Big)\supset\mathcal{Q}_{\rho}$$

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we obtain

$$\left|\pi_{\mathbf{x}}\Big((T^*\mathcal{Q}_{\rho}\setminus\Gamma_+)\cap\{\boldsymbol{\xi}=\mathbf{a}\}\Big)\right|\geq |\mathcal{Q}_{\rho}|-\left|\pi_{\mathbf{x}}\Big((T^*\mathcal{Q}_{\rho}\cap\Gamma_+)\cap\{\boldsymbol{\xi}=\mathbf{a}\}\Big)\right|.$$
(9.36)

By the first inclusion in (9.27),

$$\left|\pi_{\mathbf{x}}\left((T^*\Omega_{\rho}\cap\Gamma_+)\cap\{\boldsymbol{\xi}=\mathbf{a}\}\right)\right|\leq \left|\Omega_{\mathrm{sc},R,\mathbf{a}}\right|,\tag{9.37}$$

with this inequality expressing the fact that any parts of the scattered wave travelling in direction **a** must lie in $\Omega_{sc,R,a}$. Combining (9.36) with (9.37) yields

$$\left|\pi_{\mathbf{X}}\left((T^*\mathcal{Q}_{\rho}\setminus\Gamma_{+})\cap\{\boldsymbol{\xi}=\mathbf{a}\}\right)\right|\geq|\mathcal{Q}_{\rho}|-|\mathcal{Q}_{\mathrm{sc},R,\mathbf{a}}|.$$
(9.38)

Since $\Omega_{\text{sc},R,\mathbf{a}} \subseteq \Omega_{\rho}$, there exists $\delta > 0$ such that $|\Omega_{\rho}| - |\Omega_{\text{sc},R,\mathbf{a}}| \ge \delta |\Omega_{\rho}|$, and thus (9.35) and (9.38) imply that (9.25) holds; the proof is complete.

Remark 9.9 (What if impedance boundary conditions are imposed on Γ_R ?) If the impedance boundary condition $\partial_{\mathbf{n}} u^S - iku^S = 0$ is imposed on Γ_R (as an approximation of DtN_k), then there are additional reflections on Γ_R [84], [46, §2] μ^S has support on the incoming set, and Lemma 9.8 no longer holds.

Remark 9.10 (Proving Theorem 9.1 in the trapping case) In the trapping case, $||u(k)||_{L^2(\Omega_P)}$ may no longer be uniformly bounded, as it is in Lemma 9.4, since (3.5) no longer holds with C_{sol} bounded independently of k. If a subsequence of k's exists along which $||u(k)||_{L^2(\Omega_R)}$ is uniformly bounded, we may obtain a contradiction by the same argument as above by considering this subsequence. Thus, we can assume, without loss of generality, that $||u(k)||_{L^2(\Omega_R)} \to \infty$. Now instead of defining defect measures of u(k), one can instead define defect measures of $u(k)/||u(k)||_{L^2(\Omega_P)}$. If R is sufficiently large, then the bound in [19, Theorem 1.1] (i.e. the fact that the nontrapping cut-off resolvent estimate holds, even under trapping, if the supports of the cut-offs on both sides are sufficiently far away from the scatterer) implies that $v(k) := u(k)/||u(k)||_{L^2(\Omega_R)}$ satisfies (9.6). Any defect measure of v(k) is then immediately non-zero, since $\mu(\chi^2) \ge 1$ for any χ with supp $\chi \supset B_R$. Lemma 9.7 goes through as before after multiplying both sides of (9.19) by $||u(k)||_{L^2(\Omega_R)}^{-2}$. The main change needed to the rest of the proof is to take into account the fact that a defect measure of $u^{I}(k)/||u(k)||_{L^{2}(\Omega_{R})}$ is zero when $||u(k)||_{L^{2}(\Omega_{R})}$ grows through the sequence k_{ℓ} associated with that measure. In this situation, however, the bound (9.23) becomes $\mu(T^*A) \leq \mu(T^*A \cap \Gamma_+)$; combining this with (9.24) we obtain $\mu(T^*A) \leq 2\mu(T^*\Omega_R)$, from which the key bound (9.21) (and hence the result of the theorem) follows.

10 Proof of Theorems 4.1 and 4.2

Lemma 10.1 (Aubin-Nitsche analogue via elliptic projection) Assuming that the Galerkin solution u_h to the variational problem (2.11) exists, if

$$hk^2\eta(\mathcal{H}_h) \le \mathcal{C}_1, \quad where \quad \mathcal{C}_1 := \frac{1}{2\sqrt{2}C_{\operatorname{cont}\star}C_{H^2\star}C_{\operatorname{int}}n_{\max}},$$
 (10.1)

then

$$\|u - u_h\|_{L^2(\Omega_R)} \le 2C_{\operatorname{cont}\star}\eta(\mathcal{H}_h) \|u - w_h\|_{H^1_k(\Omega_R)} \quad \text{for all } w_h \in \mathcal{H}_h.$$

Proof Let $\xi = S^*(u - u_h)$; i.e. ξ is the solution of variational problem

find
$$\xi \in \mathcal{H}$$
 such that $a(v, \xi) = (v, u - u_h)_{L^2(\Omega_R)}$ for all $v \in \mathcal{H}$.

Then, by Galerkin orthogonality (7.6) and the definition of $a_{\star}(\cdot, \cdot)$ (7.1), for all $v_h \in \mathcal{H}_h$,

$$\|u - u_h\|_{L^2(\Omega_R)}^2 = a(u - u_h, \xi) = a(u - u_h, \xi - v_h),$$

= $a_{\star}(u - u_h, \xi - v_h)_{L^2(\Omega_R)} - k^2(n(u - u_h), \xi - v_h)_{L^2(\Omega_R)}.$
(10.2)

We choose $v_h = \mathcal{P}_h \xi$, and then use (in the following order) (i) the Galerkin orthogonality (7.6), (ii) continuity of $a_{\star}(\cdot, \cdot)$, (iii) the bound (7.8), (iv) the upper bound in the norm equivalence (7.4) and the bound (7.7), and (v) the consequence (8.4) of the definition of η to obtain that, for all $w_h \in \mathcal{H}_H$,

$$\begin{split} \|u - u_{h}\|_{L^{2}(\Omega_{R})}^{2} &= a_{\star}(u - w_{h}, \xi - \mathcal{P}_{h}\xi)_{L^{2}(\Omega_{R})} - k^{2}(n(u - u_{h}), \xi - \mathcal{P}_{h}\xi)_{L^{2}(\Omega_{R})} \\ &\leq \|u - w_{h}\|_{\star} \|\xi - \mathcal{P}_{h}\xi\|_{\star} + k^{2}n_{\max} \|u - u_{h}\|_{L^{2}(\Omega_{R})} \|\xi - \mathcal{P}_{h}\xi\|_{L^{2}(\Omega_{R})} \\ &\leq \left(\|u - w_{h}\|_{\star} + hk^{2}\sqrt{2}C_{\operatorname{int}}C_{H^{2}\star}\sqrt{C_{\operatorname{cont}\star}}n_{\max} \|u - u_{h}\|_{L^{2}(\Omega_{R})}\right) \|\xi - \mathcal{P}_{h}\xi\|_{\star} \\ &\leq \left(\sqrt{C_{\operatorname{cont}\star}} \|u - w_{h}\|_{H^{1}_{k}} + hk^{2}\sqrt{2}C_{\operatorname{int}}C_{H^{2}\star}\sqrt{C_{\operatorname{cont}\star}}n_{\max} \|u - u_{h}\|_{L^{2}}\right) \\ &\times \sqrt{C_{\operatorname{cont}\star}} \min_{v_{h}\in\mathcal{H}_{h}} \|\xi - v_{h}\|_{H^{1}_{k}(\Omega_{R})} \\ &\leq \left(\sqrt{C_{\operatorname{cont}\star}} \|u - w_{h}\|_{H^{1}_{k}} + hk^{2}\sqrt{2}C_{\operatorname{int}}C_{H^{2}\star}\sqrt{C_{\operatorname{cont}\star}}n_{\max} \|u - u_{h}\|_{L^{2}}\right) \\ &\times \sqrt{C_{\operatorname{cont}\star}}\eta(\mathcal{H}_{h}) \|u - u_{h}\|_{L^{2}(\Omega_{R})}; \end{split}$$
(10.3)

the result then follows.

Remark 10.2 (Advantage of elliptic-projection over standard duality argument) Comparing (10.2) and (10.3) we see the advantage of the elliptic-projection argument over the standard duality argument: in (10.3), Galerkin orthogonality for $a_{\star}(\cdot, \cdot)$ has

allowed us to obtain $u - w_h$ (with w_h arbitrary) as opposed to $u - u_h$ in the first argument of the sesquilinear form on the right-hand side, leading to the bound (5.3) instead of (5.2). The price for this is that we have an additional L^2 inner product on the right-hand side of (10.3), and controlling this leads to the condition (10.1).

Recall that, by the Cauchy–Schwarz inequality and the inequality (3.3), $a(\cdot, \cdot)$ is continuous, i.e., for all $u, v \in \mathcal{H}$,

$$|a(u, v)| \le C_{\text{cont}} \|u\|_{H^{1}_{k}(\Omega_{R})} \|v\|_{H^{1}_{k}(\Omega_{R})}, \qquad (10.4)$$

where $C_{\text{cont}} := \max \{A_{\max}, n_{\max}\} + C_{\text{DtN1}}$.

Lemma 10.3 Assuming that the Galerkin solution u_h to the variational problem (2.11) exists, if (10.1) holds, then

$$\|u - u_h\|_{H^1_k(\Omega_R)} \le \left(C_2 hk + C_3 hk^2 \eta(\mathcal{H}_h) \right) \|u\|_{H^1_k(\Omega_R)},$$
(10.5)

where

$$\mathcal{C}_2 := \frac{\sqrt{2}C_{\text{cont}}C_{\text{int}}C_{\text{osc}}}{A_{\min}} \quad and \quad \mathcal{C}_3 := \frac{4C_{\text{cont}\star}C_{\text{int}}C_{\text{osc}}\sqrt{n_{\max} + A_{\min}}}{\sqrt{A_{\min}}}$$

Proof Since DtN_k satisfies the inequality (3.4), and A and *n* satisfy the inequalities (2.1) and (2.2), $a(\cdot, \cdot)$ (2.8) satisfies the Gårding inequality

$$\Re a(v,v) \ge A_{\min} \|v\|_{H^1_k(\Omega_R)}^2 - k^2 (n_{\max} + A_{\min}) \|v\|_{L^2(\Omega_R)}^2.$$
(10.6)

Using Galerkin orthogonality (2.12) and continuity of $a(\cdot, \cdot)$ (10.4), we find that that (5.1) holds for any $v_h \in \mathcal{H}_h$. Using first the inequality (5.4) with $\alpha = ||u - u_h||_{H_k^1(\Omega_R)}$, $\beta = C_{\text{cont}} ||u - v_h||_{H_k^1(\Omega_R)}$, $\varepsilon = A_{\min}$, and then Lemma 10.1, we find that if (10.1) holds, then, for any $v_h \in \mathcal{H}_h$,

$$\frac{A_{\min}}{2} \|u - u_{h}\|_{H_{k}^{1}(\Omega_{R})}^{2} \\
\leq \frac{(C_{\text{cont}})^{2}}{2A_{\min}} \|u - v_{h}\|_{H_{k}^{1}(\Omega_{R})}^{2} + k^{2} (n_{\max} + A_{\min}) \|u - u_{h}\|_{L^{2}(\Omega_{R})}^{2} \\
\leq \left[\frac{(C_{\text{cont}})^{2}}{2A_{\min}} + 4k^{2} (n_{\max} + A_{\min}) (C_{\text{cont}\star})^{2} (\eta(\mathcal{H}_{h}))^{2} \right] \|u - v_{h}\|_{H_{k}^{1}(\Omega_{R})}^{2}, \quad (10.7)$$

By the consequence (3.11) of the definition of C_{int} and the bound (3.6)/(9.1),

$$\|u - I_h u\|_{H^1_k(\Omega_R)} \le \sqrt{2}hC_{\text{int}}|u|_{H^2(\Omega_R)} \le \sqrt{2}hkC_{\text{int}}C_{\text{osc}} \|u\|_{H^1_k(\Omega_R)}.$$
 (10.8)

Choosing $v_h = I_h u$ in (10.7), using (10.8), taking the square root and using the inequality $\sqrt{a^2 + b^2} \le a + b$ for all a, b > 0, we find the result (10.5).

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Proof of (9.25) (Proof of Theorem 4.1) Under the assumption that the Galerkin solution u_h exists, the fact that the bound (4.2) holds under the condition (4.1) follows from combining Lemma 10.3 with the bound (8.5) on η . To prove that u_h exists under the condition (4.1), recall that, since the variational problem (2.11) is equivalent to a linear system of equations in a finite-dimensional space, existence of a solution follows from uniqueness. Suppose that there exists a $\tilde{u}_h \in \mathcal{H}_h$ such that $a(\tilde{u}_h, v_h) = 0$ for all $v_h \in \mathcal{H}_h$; to prove uniqueness, we need to show that $\tilde{u}_h = 0$. Let \tilde{u} be such that $a(\tilde{u}, v) = 0$ for all $v \in \mathcal{H}$, so that \tilde{u}_h is the Galerkin approximation to \tilde{u} . Repeating the argument in the first part of the proof we see that the condition (4.1) holds then the bound (4.2) holds (with u replaced by \tilde{u} and u_h replaced by \tilde{u}_h). By Lemma 2.4, $\tilde{u} = 0$, so (4.2) implies that $\tilde{u}_h = 0$ and the proof is complete.

Proof of (9.25) (Proof of Theorem 4.2) This is very similar to the proof of Theorem 4.1, except that we use the bound (8.6) on $\eta(\mathcal{H}_h)$ instead of (8.5).

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