# Wavenumber-explicit convergence of the $h p$-FEM for the full-space heterogeneous Helmholtz equation with smooth coefficients 

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## A R T I C L E I N F O

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High frequency
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#### Abstract

A convergence theory for the $h p$-FEM applied to a variety of constant-coefficient Helmholtz problems was pioneered in the papers [35], [36], [15], [34]. This theory shows that, if the solution operator is bounded polynomially in the wavenumber $k$, then the Galerkin method is quasioptimal provided that $h k / p \leq C_{1}$ and $p \geq C_{2} \log k$, where $C_{1}$ is sufficiently small, $C_{2}$ is sufficiently large, and both are independent of $k, h$, and $p$. The significance of this result is that if $h k / p=C_{1}$ and $p=C_{2} \log k$, then quasioptimality is achieved with the total number of degrees of freedom proportional to $k^{d}$; i.e., the $h p$-FEM does not suffer from the pollution effect. This paper proves the analogous quasioptimality result for the heterogeneous (i.e. variable-coefficient) Helmholtz equation, posed in $\mathbb{R}^{d}, d=2,3$, with the Sommerfeld radiation condition at infinity, and $C^{\infty}$ coefficients. We also prove a bound on the relative error of the Galerkin solution in the particular case of the plane-wave scattering problem. These are the first ever results on the wavenumber-explicit convergence of the $h p$-FEM for the Helmholtz equation with variable coefficients.


## 1. Introduction

### 1.1. Context

Over the last 10 years, a wavenumber-explicit convergence theory for the $h p$-FEM applied to the Helmholtz equation
$\Delta u+k^{2} u=-f$
was established in the papers [35], [36], [15], [34]. This theory is based on decomposing solutions of the Helmholtz equation into two components:
(i) an analytic component, satisfying bounds with the same $k$ dependence as those satisfied by the full Helmholtz solution, and
(ii) a component with finite regularity, satisfying bounds with improved $k$-dependence compared to those satisfied by the full Helmholtz solution.

Such a decomposition was obtained for

- the Helmholtz equation (1.1) posed in $\mathbb{R}^{d}, d=2$, 3 , with compactlysupported $f$, and with the Sommerfeld radiation condition
$\frac{\partial u}{\partial r}(x)-\mathrm{i} k u(x)=o\left(\frac{1}{r^{(d-1) / 2}}\right)$
as $r:=|x| \rightarrow \infty$, uniformly in $\hat{x}:=x / r$ [35, Lemma 3.5],
- the Helmholtz exterior Dirichlet problem where the obstacle has analytic boundary [36, Theorem 4.20],
- the Helmholtz interior impedance problem where the domain is either smooth $(d=2,3)$ [36, Theorem 4.10], [34, Theorem 4.5], or polygonal [36, Theorem 4.10], [15, Theorem 3.2].

This decomposition was then used to prove quasioptimality of the $h p$ FEM applied to the standard Helmholtz variational formulation in [35], [36], [15], and applied to a discontinuous Galerkin formulation in [34]. Indeed, for the standard variational formulation (defined for the fullspace problem in Definition 2.2 below) applied to the boundary value problems above, if the solution operator of the problem is bounded polynomially in $k$ (see Definition 2.6 below), then there exist $C_{1}, C_{2}$, and $C_{\text {qo }}$ (independent of $k, h$, and $p$ ) such that if
$\frac{h k}{p} \leq C_{1} \quad$ and $\quad p \geq C_{2} \log k$
then the Galerkin solution $u_{N}$ exists, is unique, and satisfies

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$\left\|u-u_{N}\right\|_{H_{k}^{1}} \leq C_{\mathrm{qo}} \min _{v_{N} \in V_{N}}\left\|u-v_{N}\right\|_{H_{k}^{1}}$,
where $V_{N}$ is the $h p$ approximation space and the norm $\|\cdot\|_{H_{k}^{1}}$ is the standard weighted $H^{1}$ norm (defined by (2.7) below). Since the total number of degrees of freedom of the approximation space is proportional to $(p / h)^{d}$, the significance of this result is that it shows there is a choice of $h$ and $p$ such that the Galerkin solution is quasioptimal, with quasioptimality constant (i.e. $C_{\text {qo }}$ ) independent of $k$, and with the total number of degrees of freedom proportional to $k^{d}$; thus, with these choices of $k$ and $p$, the $h p$-FEM does not suffer from the pollution effect [2].

Over the last few years, there has been increasing interest in the numerical analysis of the heterogeneous Helmholtz equation, i.e. the Helmholtz equation with variable coefficients
$\nabla \cdot(\mathrm{A} \nabla u)+k^{2} n u=-f ;$
see, e.g., [8], [3], [10], [18], [38], [21], [16], [29], [19]. However there do not yet exist in the literature analogous results to those in [35], [36], [15], [34] for the variable-coefficient Helmholtz equation.

### 1.2. Informal statement and discussion of the main results

The main results This paper considers the variable-coefficient Helmholtz equation (1.4) with $C^{\infty}$ coefficients posed in $\mathbb{R}^{d}, d=2,3$, with the Sommerfeld radiation condition at infinity. We obtain analogous results to those obtained in [35] for this scenario with constant coefficients. That is, we prove quasioptimality of the $h p$-FEM under the conditions (1.3) and provided that the solution operator is polynomially bounded in $k$; see Theorem 3.4 below.

We obtain this result by decomposing the solution $u$ to (1.4) into two components:
$\left.u\right|_{B_{R}}=u_{H^{2}}+u_{\mathcal{A}}$
where $u_{H^{2}} \in H^{2}\left(B_{R}\right)$ and $u_{\mathcal{A}}$ is analytic in $B_{R}$, where $B_{R}$ denotes the ball of radius $R$ centred at the origin (and $R$ is arbitrary); see Theorem 3.1 below. This is exactly analogous to the decomposition obtained in [35], except that now $u$ satisfies the variable-coefficient equation (1.4) instead of (1.1).

Overview of the ideas behind the decomposition and subsequent bounds The idea in [35] was to decompose the data $f$ in (1.1) into "low-" and "high-" frequency components, with $u_{\mathcal{A}}$ the Helmholtz solution for the low-frequency component of $f$ and $u_{H^{2}}$ the Helmholtz solution for the high-frequency component of $f$. The frequency cut-offs were defining using the indicator function
$1_{B_{\lambda k}}(\zeta):= \begin{cases}1 & \text { for }|\zeta| \leq \lambda k, \\ 0 & \text { for }|\zeta| \geq \lambda k,\end{cases}$
with $\lambda$ a free parameter (see [35, Equation 3.31] and the surrounding text). In [35] the frequency cut-off (1.5) was then used with (a) the expression for $u$ as a convolution of the fundamental solution and the data $f$, and (b) the fact that the fundamental solution is known explicitly when $\mathrm{A}=\mathrm{I}$ and $n=1$, to obtain the appropriate bounds on $u_{\mathcal{A}}$ and $u_{H^{2}}$ using explicit calculation.

In this paper we use the same idea as in [35] of decomposing into low- and high-frequency components, but apply frequency cut-offs to the solution $u$ as opposed to the data $f$. Then, given any cut-off function that is zero for $|\zeta| \geq C k$, bounding the corresponding low-frequency component $u_{\mathcal{A}}$ is relatively straightforward using basic properties of the Fourier-transform (namely the expression for the Fourier transform of a derivative and Parseval's theorem). Indeed, in Fourier space each derivative corresponds to a power of the Fourier variable $\zeta$, and the frequency cut-off means that $|\zeta| \leq C k$ for $u_{\mathcal{A}}$; i.e. every derivative of $u_{\mathcal{A}}$ brings down a power of $k$ compared to $u_{\mathcal{A}}$ (see $\S 5.3$ below). The
main difficulty therefore is showing that the high-frequency component $u_{H^{2}}$ satisfies a bound with one power of $k$ improvement over the bound satisfied by $u$.

The main idea of the present paper is that the high-frequency cut-off can be chosen so that the (scaled) Helmholtz operator
$P_{k}:=-\left(k^{-2} \nabla \cdot(\mathrm{~A} \nabla \cdot)+n\right)$
is semiclassically elliptic on the support of the high-frequency cut-off. Furthermore, choosing the cut-off function to be smooth (as opposed to discontinuous, as in (1.5)) then allows us to use basic facts about the "nice" behaviour of elliptic semiclassical pseudodifferential operators (namely, they are invertible up to a small error) to prove the required bound on $u_{H^{2}}$. (Recall that semiclassical pseudodifferential operators are just pseudodifferential operators with a large/small parameter; in this case the large parameter is $k$.)

We now discuss further the frequency cut-offs and the bound on $u_{H^{2}}$ via ellipticity.

The frequency cut-offs In contrast to (1.5), we choose $\chi_{\mu} \in C_{\text {comp }}^{\infty}\left(\mathbb{R}^{d}\right)$ such that
$\chi_{\mu}\left(k^{-2}|\zeta|^{2}\right)= \begin{cases}1 & \text { for }|\zeta| \leq \sqrt{\mu} k, \\ 0 & \text { for }|\zeta| \geq \sqrt{2 \mu} k,\end{cases}$
where the parameter $\mu$ is chosen later in the argument. With the Fourier transform and its inverse defined by
$\mathcal{F} \varphi(\zeta):=\int_{\mathbb{R}^{d}} \exp (-\mathrm{i} x \cdot \zeta) \varphi(x) \mathrm{d} x$ and
$\mathcal{F}^{-1} \psi(x):=(2 \pi)^{-d} \int_{\mathbb{R}^{d}} \exp (\mathrm{i} x \cdot \zeta) \psi(\zeta) \mathrm{d} \zeta$,
we define the low-frequency cut-off $\Pi_{L}$ by
$\Pi_{L} v(x):=\mathcal{F}^{-1}\left(\chi_{\mu}\left(k^{-2}|\zeta|^{2}\right) \mathcal{F} v(\zeta)\right)$,
and the high-frequency cut-off $\Pi_{H}$ by
$\Pi_{H} v(x):=\mathcal{F}^{-1}\left(\left(1-\chi_{\mu}\left(k^{-2}|\zeta|^{2}\right)\right) \mathcal{F} v(\zeta)\right)$,
so that $\Pi_{L}+\Pi_{H}=I$. We let $\varphi \in C_{c}^{\infty}$ be equal to one on $B_{R+1}$ and vanish outside $B_{R+2}$, and then
$u_{\mathcal{A}}:=\left.\Pi_{L}(\varphi u)\right|_{B_{R}} \quad$ and $\quad u_{H^{2}}:=\left.\Pi_{H}(\varphi u)\right|_{B_{R}}$.
The bound on the high-frequency component $u_{H^{2}}$ via ellipticity Recall that a PDE is elliptic if its principal symbol is non-zero. The concept of ellipticity for semiclassical differential operators (or, more generally, semiclassical pseudodifferential operators) is analogous, except that it now involves the semiclassical principal symbol (see (4.17) below). The semiclassical principal symbol of $P_{k}(1.6)$ is
$\langle\mathrm{A} \xi, \xi\rangle-n$,
where $\langle\cdot, \cdot\rangle$ denotes the $\ell^{2}$ inner product and $\xi=k^{-1} \zeta$ (see (4.12) below and the surrounding text).

If the parameter $\mu$ in the cut-off function $\chi_{\mu}(1.7)$ is chosen to be a certain function of $A$ and $n$ (see (5.7) below), then the symbol (1.12) is bounded away from zero when $k^{-2}|\zeta|^{2} \geq \mu$, i.e. in the region of Fourier space where $\Pi_{H}$ is non-zero; one therefore describes $P_{k}$ as "microlocally elliptic", where the adjective "microlocal" indicates that we have ellipticity on just a region of phase space (rather than on all of phase space in the more familiar global ellipticity).

These ellipticity properties are then used with the standard microlocal elliptic estimate for pseudodifferential operators, appearing in the semiclassical setting in, e.g., [14, Appendix E], and stated in this setting
as Theorem 4.3 below. The whole point is that a semiclassical pseudodifferential operator that is elliptic in some region of phase space can be inverted (up to some small error) in that region, and the norm of the inverse is bounded uniformly in the large parameter (here $k$ ) as long as one uses weighted norms (analogous to the familiar $H_{k}^{1}$ norm (2.7)).

The result is that $u_{H^{2}}$ satisfies a bound with one power of $k$ improvement over the bound satisfied by $u$ (compare (3.1) and (2.12)). To give a simple illustration of how ellipticity can give this improved $k$-dependence, we contrast the solutions of
$P_{k} u:=-\left(\Delta+k^{2}\right) u=f \quad$ and $\quad \widetilde{P}_{k} v:=-\left(\Delta-k^{2}\right) v=f$,
with both equations posed in $\mathbb{R}^{d}$ with compactly-supported $f$, and with $u$ satisfying the Sommerfeld radiation condition (1.2) and $v$ satisfying boundedness at infinity. The $L^{2} \rightarrow L^{2}$ bounds that are sharp in terms of $k$-dependence are
$\|u\|_{L^{2}\left(B_{R}\right)} \lesssim k^{-1}\|f\|_{L^{2}\left(\mathbb{R}^{d}\right)} \quad$ and $\quad\|v\|_{L^{2}\left(\mathbb{R}^{d}\right)} \lesssim k^{-2}\|f\|_{L^{2}\left(\mathbb{R}^{d}\right)}$,
with the former given by Part (i) of Theorem 2.7, and the latter following from the Lax-Milgram theorem. The operator $P_{k}$ is not semiclassically elliptic on all of phase space (its semiclassical principal symbol is $|\xi|^{2}-1$ ), whereas $\widetilde{P}_{k}$ is semiclassically elliptic on all of phase space (its semiclassical principal symbol is $|\xi|^{2}+1$ ); we therefore see that ellipticity has resulted in the solution operator having improved $k$-dependence. The proof of the bound on $u_{H^{2}}$ is more technical, but the idea - that the improvement in $k$-dependence comes from ellipticity - is the same.

The assumption that the solution operator is polynomially bounded in $k$ We need to assume that the solution operator is polynomially bounded in $k$ (in sense of Definition 2.6 below), both in proving the bound on $u_{H^{2}}$, and in proving quasi-optimality of the $h p$-FEM.

The $k$-dependence of the Helmholtz solution operator depends on whether the problem is trapping or nontrapping. For the heterogeneous Helmholtz equation (1.4) posed in $\mathbb{R}^{d}$ (i.e. with no obstacle), trapping can be created by the coefficients A and $n$; see, e.g., [39]. If the problem is nontrapping, then the Helmholtz solution operator (measured in the natural norms) is bounded in $k$. However, under the strongest form of trapping, the Helmholtz solution operator can grow exponentially in $k$ [39]. Nevertheless, it has recently been proved that, if a set of frequencies of arbitrarily small measure is excluded, then the solution operator is polynomially bounded under any type of trapping [28]. Therefore, the result that the $h p$-FEM is quasi-optimal holds for a wide class of Helmholtz problems; see Corollary 3.5 below.

Why do we need $C^{\infty}$ coefficients? As highlighted above, our proof of the decomposition relies on standard results about semiclassical pseudodifferential operators (recapped in §4). These results are usually stated for $C^{\infty}$ symbols, and thus to fit into this framework A and $n$ must be $C^{\infty}$. However, examining the results we use, we see that we only need the symbol of the PDE to be in $C^{L}$ where $L$ depends only on the dimension $d$ and on the exponent $M$ appearing in the assumption that the solution operator is polynomially bounded (see Definitions 2.5 and 2.6 below). Therefore, while we consider $\mathrm{A}, n \in C^{\infty}$ to easily use results about semiclassical pseudodifferential operators from [52], [14, Appendix E], our results hold for $\mathrm{A} \in C^{L}$ and $n \in C^{L}$, where $L=L(d, M)$.

Extending the decomposition result to the solution of other PDEs Our proof of the decomposition result only relies on the principal symbol of the differential operator being bounded below at infinity (in the sense of (3.8) below). Therefore, the decomposition result Theorem 3.1 is valid for a much larger class of PDEs (and indeed pseudodifferential operators) than (1.4); see Remark 3.7 below for more details.

In the follow-up paper [27], we use the ideas of the present paper combined with much more sophisticated tools of semiclassical and microlocal analysis (namely the black-box scattering framework of Sjöstrand-Zworski [45], the Helffer-Sjöstrand functional calculus [23],
and associated results by Helffer, Robert, and Sjöstrand [22], [40], [44]) to prove analogous decompositions for a wide variety of scattering problems (albeit with slightly weaker estimates on $u_{\mathcal{A}}$ ). In particular, the main result of the present paper, Theorem 3.1, is rederived in this more general context as [27, Theorem 1.16].

We also note that, as announced in the abstract [4], Bernkopf, Chaumont-Frelet, and Melenk are also studying the question of $k$ explicit convergence of the $h p$-FEM for the Helmholtz equation with variable coefficients.

Outline of the paper $\S 2$ gives the definitions of the boundary-value problem and the finite-element method. §3 states the main results. §4 recaps results about semiclassical pseudodifferential operators, with [52] and [14, Appendix E] as the main references. §5 proves the result about the decomposition $\left.u\right|_{B_{R}}=u_{H^{2}}+u_{\mathcal{A}}$ (Theorem 3.1). §6 proves the result about quasioptimality of the $h p$-FEM (Theorem 3.4).

## 2. Formulation of the problem

### 2.1. The boundary value problem

Assumption 2.1 (Assumptions on the coefficients). $\mathrm{A} \in C^{\infty}\left(\mathbb{R}^{d}, \mathrm{SPD}\right)$ (where SPD is the set of $d \times d$ real, symmetric, positive-definite matrices) is such that $\operatorname{supp}(\mathrm{I}-\mathrm{A})$ is compact in $\mathbb{R}^{d}$ and there exist $0<A_{\min } \leq A_{\max }<\infty$ such that, in the sense of quadratic forms,
$A_{\text {min }} \leq \mathrm{A}(x) \leq A_{\max } \quad$ for all $x \in \mathbb{R}^{d}$.
$n \in C^{\infty}\left(\mathbb{R}^{d}, \mathbb{R}\right)$ is such that $\operatorname{supp}(1-n)$ is compact in $\mathbb{R}^{d}$ and there exist $0<n_{\text {min }} \leq n_{\text {max }}<\infty$ such that
$n_{\text {min }} \leq n(x) \leq n_{\text {max }} \quad$ for all $x \in \mathbb{R}^{d}$.

Let $R>0$ be such that $\operatorname{supp}(\mathrm{I}-\mathrm{A}) \cup \operatorname{supp}(1-n) \Subset B_{R}$, where $B_{R}$ denotes the ball of radius $R$ about the origin and $\Subset$ denotes compact containment. Let $\gamma$ and $\partial_{n}$ denote the Dirichlet and Neumann traces, respectively, on $\partial B_{R}$, where the normal vector for the Neumann trace points out of $B_{R}$.

Define $\operatorname{DtN}_{k}: H^{1 / 2}\left(\partial B_{R}\right) \rightarrow H^{-1 / 2}\left(\partial B_{R}\right)$ to be the Dirichlet-toNeumann map for the equation $\Delta u+k^{2} u=0$ posed in the exterior of $B_{R}$ with the Sommerfeld radiation condition (1.2). The definition of $\mathrm{DtN}_{k}$ in terms of Hankel functions and polar coordinates (when $d=2$ )/spherical polar coordinates (when $d=3$ ) is given in, e.g., [35, Equations 3.7 and 3.10].

Definition 2.2 (Heterogeneous Helmholtz Problem on $\mathbb{R}^{d}$ ). Given $A$ and $n$ satisfying Assumption $2.1, R>0$ such that $\operatorname{supp}(\mathrm{I}-\mathrm{A}) \cup \operatorname{supp}(1-n) \Subset$ $B_{R}, k>0$, and $F \in\left(H^{1}\left(B_{R}\right)\right)^{*}, u \in H^{1}\left(B_{R}\right)$ satisfies the Heterogeneous Helmholtz Problem on $\mathbb{R}^{d}$ if $u$ satisfies the variational problem
find $u \in H^{1}\left(B_{R}\right)$ such that $\quad a(u, v)=F(v) \quad$ for all $v \in H^{1}\left(B_{R}\right)$,
where
$a(u, v):=\int_{B_{R}}\left((\mathrm{~A} \nabla u) \cdot \overline{\nabla v}-k^{2} n u \bar{v}\right)-\left\langle\operatorname{DtN}_{k}(\gamma u), \gamma v\right\rangle_{\partial B_{R}}$,
where $\langle\cdot, \cdot\rangle_{\partial B_{R}}$ denotes the duality pairing on $\partial B_{R}$ that is linear in the first argument and antilinear in the second.

Lemma 2.3 (Helmholtz boundary value problems included in Definition 2.2).
(i) If

$$
\begin{equation*}
F(v):=\int_{B_{R}} f \bar{v} \tag{2.5}
\end{equation*}
$$

with $f \in L^{2}\left(B_{R}\right)$, then the solution $u$ to (2.3) equals $\left.\widetilde{u}\right|_{B_{R}}$, where $\widetilde{u} \in$ $H_{\mathrm{loc}}^{1}\left(\mathbb{R}^{d}\right)$ is the solution to

$$
\nabla \cdot(\mathrm{A} \nabla \widetilde{u})+k^{2} n \tilde{u}=-f \quad \text { in } \mathbb{R}^{d}
$$

and $\tilde{u}$ satisfies the Sommerfeld radiation condition (1.2).
(ii) If

$$
\begin{equation*}
F(v):=\int_{\partial B_{R}}\left(\partial_{n} u^{I}-\operatorname{DtN}_{k}\left(\gamma u^{I}\right)\right) \overline{\gamma v} \quad \text { with } \quad u^{I}(x):=\exp (\mathrm{i} k x \cdot a) \tag{2.6}
\end{equation*}
$$

where $a \in \mathbb{R}^{d}$ with $|a|=1$, then the solution $u$ to (2.3) equals $\left.\widetilde{u}\right|_{B_{R}}$, where $\tilde{u} \in H_{\mathrm{loc}}^{1}\left(\mathbb{R}^{d}\right)$ is the solution of the Helmholtz plane-wave scattering problem; i.e.
$\nabla \cdot(\mathrm{A} \nabla \tilde{u})+k^{2} n \tilde{u}=0 \quad$ in $\mathbb{R}^{d}$,
and $\widetilde{u}^{S}:=\tilde{u}-u^{I}$ satisfies the Sommerfeld radiation condition (1.2).
Part (i) of Lemma 2.3 is proved in, e.g., [20, Lemma 3.3]; the proof of Part (ii) is similar.

Let the weighted $H^{1}$ norm, $\|\cdot\|_{H_{k}^{1}\left(B_{R}\right)}$, be defined by
$\|u\|_{H_{k}^{1}\left(B_{R}\right)}^{2}:=\|\nabla u\|_{L^{2}\left(B_{R}\right)}^{2}+k^{2}\|u\|_{L^{2}\left(B_{R}\right)}^{2}$.
Lemma 2.4. The solution of the Heterogeneous Helmholtz Problem on $\mathbb{R}^{d}$ (defined in Definition 2.2) exists, is unique, and there exists $C(k, A, n, R)>0$ such that
$\|u\|_{H_{k}^{1}\left(B_{R}\right)} \leq C\|F\|_{\left(H_{k}^{1}\left(B_{R}\right)\right)^{*}} \quad$ for all $k>0$.
Proof. Uniqueness follows from the unique continuation principle; see [20, §1], [21, §2] and the references therein. Since $a(\cdot, \cdot)$ satisfies a Gårding inequality (see (6.4) below), Fredholm theory then gives existence and the bound (2.8).

Properties of $\mathrm{DtN}_{k}$ and $a(\cdot, \cdot)$ We use later the following two properties of $\mathrm{DtN}_{k}$ : given $k_{0}, R_{0}>0$, there exists $C_{\mathrm{DtN}}=C_{\mathrm{DtN}}\left(k_{0} R_{0}\right)$ such that, for all $k \geq k_{0}$ and $R \geq R_{0}$,
$\left|\left\langle\operatorname{DtN}_{k}(\gamma u), \gamma v\right\rangle_{\partial B_{R}}\right\rangle \mid \leq C_{\mathrm{DtN} 1}\|u\|_{H_{k}^{1}\left(B_{R}\right)}\|v\|_{H_{k}^{1}\left(B_{R}\right)}$
for all $u, v \in H^{1}\left(B_{R}\right)$, and
$-\Re\left\langle\mathrm{DtN}_{k} \phi, \phi\right\rangle_{\partial B_{R}} \geq 0 \quad$ for all $\phi \in H^{1 / 2}\left(\partial B_{R}\right)$.
For a proof of (2.9), see [35, Lemma 3.3]. For a proof of (2.10), see [37, Theorem 2.6.4] (for $d=3$ ) and [7, Corollary 3.1] or [35, Lemma 3.10] (for $d=2,3$ ).

Let $C_{\text {cont }}=C_{\text {cont }}\left(\mathrm{A}, n, R, k_{0}\right)$ be the continuity constant of the sesquilinear form $a(\cdot, \cdot)$ (defined in (2.4)) in the norm $\|\cdot\|_{H_{k}^{1}\left(B_{R}\right)}$; i.e.
$a(u, v) \leq C_{\text {cont }}\|u\|_{H_{k}^{1}\left(B_{R}\right)}\|v\|_{H_{k}^{1}\left(B_{R}\right)} \quad$ for all $u, v \in H^{1}\left(B_{R}\right)$ and $k \geq k_{0}$.
By the Cauchy-Schwarz inequality and (2.9),
$C_{\text {cont }} \leq \max \left\{A_{\max }, n_{\max }\right\}+C_{\mathrm{DtN} 1}$.

### 2.2. The behaviour of the solution operator for large $k$

Definition $2.5\left(C_{\text {sol }}\right)$. Given $f \in L^{2}\left(B_{R}\right)$, let $u$ be the solution of the heterogeneous Helmholtz equation (1.4) with the Sommerfeld radiation condition (1.2) (i.e. $u$ is the solution of the variational problem (2.3) with $F(v)$ given by (2.5)). Given $k_{0}>0$, let $C_{\text {sol }}=C_{\text {sol }}\left(k, \mathrm{~A}, n, R, k_{0}\right)>0$ be such that
$\|u\|_{H_{k}^{1}\left(B_{R}\right)} \leq C_{\text {sol }}\|f\|_{L^{2}\left(B_{R}\right)} \quad$ for all $k>0$.
$C_{\text {sol }}$ exists by Lemma 2.4; indeed, with $C$ given by (2.8), $C_{\text {sol }}:=C / k$.

How $C_{\text {sol }}$ depends on $k$ is crucial to the analysis below, and to emphasise this we write $C_{\text {sol }}=C_{\text {sol }}(k)$. Below we consider $C_{\text {sol }}$ with different values of $R$, and we then write, e.g., $C_{\text {sol }}(k ; R)$ (as in the bound (3.2) below).

A key assumption in the analysis of the Helmholtz $h p$-FEM is that $C_{\text {sol }}(k)$ is polynomially bounded in $k$ in the following sense.

Definition 2.6 ( $C_{\text {sol }}$ is polynomially bounded in $k$ ). Given $k_{0}$ and $K \subset$ $\left[k_{0}, \infty\right), C_{\text {sol }}(k)$ is polynomially bounded for $k \in K$ if there exists $C>0$ and $M>0$ such that
$C_{\text {sol }}(k) \leq C k^{M}$ for all $k \in K$,
where $C$ and $M$ are independent of $k$ (but depend on $k_{0}$ and possibly also on $K, \mathrm{~A}, n, d, R)$.

There exist $C^{\infty}$ coefficients A and $n$ such that $C_{\text {sol }}\left(k_{j}\right) \geq c_{1} \exp \left(c_{2} k_{j}\right)$ for $0<k_{1}<k_{2}<\ldots$ with $k_{j} \rightarrow \infty$ as $j \rightarrow \infty$, see [39], but this exponential growth is the worst-possible, since $C_{\text {sol }}(k) \leq c_{3} \exp \left(c_{4} k\right)$ for all $k \geq k_{0}$ by [5, Theorem 2]. We now recall results on when $C_{\text {sol }}(k)$ is polynomially bounded in $k$.

Theorem 2.7 (Conditions under which $C_{\text {sol }}(k)$ is polynomially bounded in $k)$.
(i) A and $n$ are $C^{\infty}$ and nontrapping (i.e. all the trajectories of the Hamiltonian flow defined by the symbol of (1.4) starting in $B_{R}$ leave $B_{R}$ after a uniform time), then $C_{\text {sol }}(k)$ is independent of $k$ for all $k$, i.e., (2.13) holds for all $k$ with $M=0$.
(ii) If $n=1$ and A is $C^{0,1}$ then, given $k_{0}>0$ and $\delta>0$ there exists a set $J \subset\left[k_{0}, \infty\right)$ with $|J| \leq \delta$ such that

$$
\begin{equation*}
C_{\mathrm{sol}}(k) \leq C k^{5 d / 2+1+\varepsilon} \quad \text { for all } k \in\left[k_{0}, \infty\right) \backslash J \tag{2.14}
\end{equation*}
$$

for any $\varepsilon>0$, where $C$ depends on $\delta, \varepsilon, d, k_{0}$, and A . If A is $C^{1, \sigma}$ for some $\sigma>0$ then the exponent is reduced to $5 d / 2+\varepsilon$.

References for the proof. (i) is proved using either (a) the propagation of singularities results of [13] combined with either the parametrix argument of [48, Theorem 3]/[49, Chapter 10, Theorem 2] or LaxPhillips theory [30], or (b) the defect-measure argument of [6, Theorem 1.3 and §3]. It has recently been proved that, for this situation, $C_{\text {sol }}$ is proportional to the length of the longest trajectory in $B_{R}$; see $[16$, Theorems 1 and 2, and Equation 6.32].
(ii) is proved in [28, Theorem 1.1 and Corollary 3.6].

### 2.3. The finite-element method

Let $\left(V_{N}\right)_{N=0}^{\infty}$ be a sequence of finite-dimensional subspaces of $H^{1}\left(B_{R}\right)$ that converge to $H^{1}\left(B_{R}\right)$ in the sense that, for all $v \in H^{1}\left(B_{R}\right)$,
$\lim _{N \rightarrow \infty}\left(\min _{v_{N} \in V_{N}}\left\|v-v_{N}\right\|_{H^{1}\left(B_{R}\right)}\right)=0$.
Later we specialise to the triangulations described in [35, §5], which allow curved elements and thus fit $\partial B_{R}$ exactly.

The finite-element method for the variational problem (2.3) is the Galerkin method applied to the variational problem (2.3), i.e.
find $u_{N} \in V_{N}$ such that $a\left(u_{N}, v_{N}\right)=F\left(v_{N}\right)$ for all $v_{N} \in V_{N}$.

## 3. Statement of the main results

Theorem 3.1 (Decomposition of the solution). Let A and n satisfy Assumption 2.1 and let $R>0$ be such that $\operatorname{supp}(\mathrm{I}-\mathrm{A}) \cup \operatorname{supp}(1-n) \Subset B_{R}$. Given $f \in L^{2}\left(B_{R}\right)$, let $u$ satisfy $\nabla \cdot(\mathrm{A} \nabla u)+k^{2} n u=-f$ in $\mathbb{R}^{d}$ and the Sommerfeld radiation condition (1.2).

If $C_{\text {sol }}(k)$ is polynomially bounded (in the sense of Definition 2.6) for $k \in K \subset\left[k_{0}, \infty\right)$, then there exist $C_{3}, C_{4}, C_{5}>0$ such that
$\left.u\right|_{B_{R}}=u_{H^{2}}+u_{\mathcal{A}}$
where $u_{H^{2}} \in H^{2}\left(B_{R}\right)$ with
$\left\|\partial^{\alpha} u_{H^{2}}\right\|_{L^{2}\left(B_{R}\right)} \leq C_{3} k^{|\alpha|-2}\|f\|_{L^{2}\left(B_{R}\right)}$
for all $|\alpha| \leq 2$ and for all $k \in K \subset\left[k_{0}, \infty\right)$,
and $u_{\mathcal{A}} \in C^{\infty}\left(B_{R}\right)$ with

$$
\begin{align*}
& \left\|\partial^{\beta} u_{\mathcal{A}}\right\|_{L^{2}\left(B_{R}\right)} \leq C_{\text {sol }}(k ; R+2) C_{4}\left(C_{5} k\right)^{|\beta|-1}\|f\|_{L^{2}\left(B_{R}\right)} \\
& \quad \text { for all } \beta \text { and for all } k \in K \subset\left[k_{0}, \infty\right), \tag{3.2}
\end{align*}
$$

where $C_{3}, C_{4}$, and $C_{5}$ depend on $\mathrm{A}, n, d$, and $k_{0}$, but are independent of $k$, $f, \alpha$, and $\beta$.

Remark 3.2 ( $u_{\mathcal{A}}$ is analytic). Since $C_{4}$ and $C_{5}$ are independent of $\beta$, the bound (3.2) implies that $u_{\mathcal{A}}$ is in the class of analytic functions on $B_{R}$, $\mathcal{A}\left(\boldsymbol{B}_{R}\right)$, defined by
$\mathcal{A}\left(B_{R}\right):=\left\{v \in \bigcap_{n \in \mathbb{N}} H^{n}\left(B_{R}\right): \exists c_{0}, c_{1}>0\right.$, independent of $n$,

$$
\text { such that } \left.|u|_{H^{n}\left(B_{R}\right)} \leq c_{1} c_{0}^{n} n!\right\}
$$

where $|u|_{H^{n}}^{2}:=\sum_{|\alpha|=n}\left\|\partial^{\alpha} u\right\|_{L^{2}}^{2}$. See, e.g., [11, §1.1.b], both for this definition, and for how the definition implies convergence of the Taylor series of elements of $\mathcal{A}\left(B_{R}\right)$ at every point in $\overline{B_{R}}$.

Remark 3.3 (The bounds of Theorem 3.1 written with the notation $\nabla^{n}$ ). The analogous bounds to (3.1) and (3.2) in [35], [36] are written using the notation
$\left|\nabla^{n} u(x)\right|^{2}:=\sum_{|\alpha|=n} \frac{n!}{\alpha!}\left|\partial^{\alpha} u(x)\right|^{2}$.
Since $\sum_{|\alpha|=n}(n!/ \alpha!)=d^{n}$,
if $\left\|\partial^{\alpha} u\right\|_{L^{2}\left(B_{R}\right)} \leq C_{1}\left(C_{2}\right)^{|\alpha|}$ for all $\alpha$ with $|\alpha|=n$,
then $\left\|\nabla^{n} u\right\|_{L^{2}\left(B_{R}\right)} \leq C_{1}\left(C_{2} \sqrt{d}\right)^{n}$,
and so the bounds (3.1) and (3.2) can also be written as bounds on $\left\|\nabla^{n} u_{H^{2}}\right\|_{L^{2}\left(B_{R}\right)}$ and $\left\|\nabla^{n} u_{\mathcal{A}}\right\|_{L^{2}\left(B_{R}\right)}$ respectively.

The following result about quasioptimality of the $h p$-FEM is then obtained by combining Theorem 3.1, well-known results about the convergence of the Galerkin method based on duality arguments (recapped in Lemma 6.4 below), and results about the $h p$ approximation spaces in [35, §5] (used in Lemma 6.5 below).

Theorem 3.4 (Quasioptimality of the $h p$-FEM if $C_{\text {sol }}(k)$ is polynomially bounded). Let $d=2$ or 3, and let $k_{0}>0$. Let $\left(V_{N}\right)_{N=0}^{\infty}$ be the piecewisepolynomial approximation spaces described in [35, §5] (where, in particular, the triangulations are quasi-uniform), and let $u_{N}$ be the Galerkin solution defined by (2.15).

If $C_{\text {sol }}(k)$ is polynomially bounded (in the sense of Definition 2.6) for $k \in K \subset\left[k_{0}, \infty\right)$ then there exist $C_{1}, C_{2}>0$, depending on $\mathrm{A}, n, R$, and $d$, and $k_{0}$, but independent of $k$, $h$, and $p$, such that if (1.3) holds, then, for all $k \in K$, the Galerkin solution exists, is unique, and satisfies the quasi-optimal error bound
$\left\|u-u_{N}\right\|_{H_{k}^{1}\left(B_{R}\right)} \leq C_{\mathrm{qo}} \min _{v_{N} \in V_{N}}\left\|u-v_{N}\right\|_{H_{k}^{1}\left(B_{R}\right)}$,
with
$C_{\text {qo }}:=\frac{2\left(\max \left\{A_{\max }, n_{\max }\right\}+C_{\mathrm{DtN} 1}\right)}{A_{\min }}$

Combining Theorem 3.4 with the results on $C_{\text {sol }}(k)$ recapped in Theorem 2.7, we obtain the following specific examples of coefficients A and $n$ when quasioptimality holds.

Corollary 3.5 (Quasioptimality under specific conditions on A and $n$ ). Let $d=2$ or 3 , and let $k_{0}>0$.
(i) If A and $n$ are nontrapping, then there exist $C_{1}, C_{2}>0$, depending on $\mathrm{A}, n, R$, and $d$, and $k_{0}$, but independent of $k$, $h$, and $p$, such that if (1.3) holds then, for all $k \geq k_{0}$, the Galerkin solution exists, is unique, and satisfies the quasi-optimal error bound (3.3) with $C_{\mathrm{qo}}$ given by (3.4).
(ii) If A is $C^{\infty}$ and $n=1$ then, given $\delta>0$, there exist a set $J$ with $|J| \leq \delta$ and constants $\widetilde{C}_{1}, \widetilde{C}_{2}>0$, with all three depending on $\mathrm{A}, n, R, d$, and $k_{0}$, but independent of $k$, and $\widetilde{C}_{2}$ additionally depending on $\delta$ and $k_{0}$ such that, for all $k \in\left[k_{0}, \infty\right) \backslash J$, if (1.3) holds (with $C_{1}, C_{2}$ replaced by $\widetilde{C}_{1}, \widetilde{C}_{2}$ ) then the Galerkin solution exists, is unique, and satisfies (3.3) with $C_{\text {qo }}$ given by (3.4).

For the plane-wave scattering problem (i.e. for $F(v)$ given by (2.6)), the regularity result
$|u|_{H^{2}\left(B_{R}\right)} \leq C_{\text {osc }} k\|u\|_{H_{k}^{1}\left(B_{R}\right)}$
was recently proved in [29, Theorem 9.1 and Remark 9.10], where $C_{\text {osc }}$ depends on A, $n, d$, and $R$, but is independent of $k$. The polynomial approximation bounds in [35, §B] imply that, for the sequence of approximation spaces $\left(V_{N}\right)_{N=0}^{\infty}$ described in [35, §5],
$\min _{v_{N} \in V_{N}}\left\|u-v_{N}\right\|_{H_{k}^{1}\left(B_{R}\right)} \leq C_{6} \frac{h}{p}\left(1+\frac{k h}{p}\right)|u|_{H^{2}\left(B_{R}\right)}$
where $C_{6}$ only depends on the constants in [35, Assumption 5.2] (which depend on the element maps from the reference element). Using (3.6) and (3.5) to bound the right-hand side of (3.3), we obtain the following bound on the relative error of the Galerkin solution.

Corollary 3.6 (Bound on the relative error of the Galerkin solution). Let the assumptions of Theorem 3.4 hold and, furthermore, let $F(v)$ be given by (2.6) (so that $u$ is the solution of the plane-wave scattering problem). If $C_{\text {sol }}(k)$ is polynomially bounded (in the sense of Definition 2.6) for $k \in K \subset$ $\left[k_{0}, \infty\right)$, then there exists $C_{6}>0$, independent of $k$, $h$, and $p$, such that if (1.3) holds, then, for all $k \in K$,
$\frac{\left\|u-u_{N}\right\|_{H_{k}^{1}\left(B_{R}\right)}}{\|u\|_{H_{k}^{1}\left(B_{R}\right)}} \leq C_{\mathrm{qo}} C_{6} C_{\mathrm{osc}} C_{1}\left(1+C_{1}\right)$,
with $C_{\mathrm{qo}}$ given by (3.4); i.e. the relative error can be made arbitrarily small by making $C_{1}$ smaller.

Remark 3.7 (Theorem 3.1 is valid for solutions of a much larger class of PDEs). Inspecting the proof of Theorem 3.1 below, we see that the conclusion, i.e. the decomposition $u=u_{H^{2}}+u_{\mathcal{A}}$ with $u_{H^{2}}$ and $u_{\mathcal{A}}$ satisfying the bounds (3.1) and (3.2) respectively, holds under much weaker assumptions. Indeed, the conclusion still holds under the following three assumptions only.
(i) $P_{k}$ is a family of properly-supported second-order pseudo-differential operators, with principal symbol $p_{k}(x, \zeta)$,
(ii) $p_{k}(x, \zeta)$ is coercive at infinity in the sense that
$\liminf _{|\xi| \rightarrow \infty, x \in \mathbb{R}^{d}}\langle k \xi\rangle^{-2} p_{k}(x, k \xi) \geq c>0$,
where $c>0$ does not depend on $k$, and
(iii) the solution to $P_{k} u=-f$, posed in $\mathbb{R}^{d}$ with $\operatorname{supp} f \subset B_{R}$ and $f \in$ $L^{2}\left(B_{R}\right)$, satisfies the bound

$$
\|u\|_{L^{2}\left(B_{R+2}\right)} \leq C k^{M}\|f\|_{L^{2}\left(B_{R}\right)}
$$

with $C$ and $M$ independent of $k, u$, and $f$. (In fact, the 2 in the $R+2$ on the left-hand side of the bound can be replaced by any number $>0$.)

In particular, no assumption is made about lower-order terms of $P_{k}$, or the behaviour of $u$ at infinity (such as a radiation condition).

## 4. Recap of relevant results about semiclassical pseudodifferential operators

The proof of Theorem 3.1 relies on standard results about semiclassical pseudodifferential operators. We review these here, with our default references being [52] and [14, Appendix E]. Homogeneous - as opposed to semiclassical - versions of the results in this section can be found in, e.g., [47, Chapter 7], [41, Chapter 7], [25, Chapter 6]. ${ }^{1}$

While the use of homogeneous pseudodifferential operators in numerical analysis is well established, see, e.g., [41], [25], there has been less use of semiclassical pseudodifferential operators. However, these are ideally-suited for studying the high-frequency behaviour of Helmholtz solutions. Indeed, semiclassical pseudodifferential operators are just pseudodifferential operators with a large/small parameter, and behaviour with respect to this parameter is then explicitly kept track of in the associated calculus.

The semiclassical parameter $\hbar=k^{-1} \quad$ Instead of working with the parameter $k$ and being interested in the large- $k$ limit, the semiclassical literature usually works with a parameter $h:=k^{-1}$ and is interested in the small- $h$ limit. So that we can easily recall results from this literature, we also work with the small parameter $k^{-1}$, but to avoid a notational clash with the meshwidth of the FEM, we let $\hbar:=k^{-1}$ (the notation $\hbar$ comes from the fact that the semiclassical parameter is related to Planck's constant, which is written as $2 \pi \hbar$; see, e.g., [52, §1.2], [14, Page 82], [32, Chapter 1]). In this notation, the Helmholtz equation $\nabla \cdot(\mathrm{A} \nabla u)+k^{2} n u=-f$ becomes
$P_{\hbar} u=\hbar^{2} f, \quad$ where $\quad P_{\hbar}:=-\hbar^{2} \nabla \cdot(\mathrm{~A} \nabla \cdot)-n$.
While some results in semiclassical analysis are valid in the limit $\hbar$ small, the results we recap in this section are valid for all $0<\hbar \leq \hbar_{0}$ with $\hbar_{0}<\infty$ arbitrary.

The semiclassical Fourier transform $\mathcal{F}_{\hbar}$ The semiclassical Fourier transform is defined for $\hbar>0$ by
$\mathcal{F}_{\hbar} \phi(\xi):=\int_{\mathbb{R}^{d}} \exp (-\mathrm{i} x \cdot \xi / \hbar) \phi(x) \mathrm{d} x$,
and its inverse by
$\mathcal{F}_{\hbar}^{-1} \psi(x):=(2 \pi \hbar)^{-d} \int_{\mathbb{R}^{d}} \exp (\mathrm{i} x \cdot \xi / \hbar) \psi(\xi) \mathrm{d} \xi ;$
see [52, §3.3]. Then
$\mathcal{F}_{\hbar}\left((-\mathrm{i} \hbar \partial)^{\alpha} \phi\right)=\xi^{\alpha} \mathcal{F}_{\hbar} \phi$
and
$\|\phi\|_{L^{2}\left(\mathbb{R}^{d}\right)}=\frac{1}{(2 \pi \hbar)^{d / 2}}\left\|\mathcal{F}_{\hbar} \phi\right\|_{L^{2}\left(\mathbb{R}^{d}\right)}$.

[^1]Semiclassical Sobolev spaces In the same way that it is convenient to work with the weighted $H^{1}$ norm (2.7) when studying the Helmholtz equation with parameter $k$, it is convenient to use norms weighted with $\hbar$ when studying (4.1). Therefore on the space
$H_{\hbar}^{s}\left(\mathbb{R}^{d}\right):=\left\{u \in L^{2}\left(\mathbb{R}^{d}\right),\langle\xi\rangle^{s} \mathcal{F}_{\hbar} u \in L^{2}\left(\mathbb{R}^{d}\right)\right\}$,
where $\langle\xi\rangle:=\left(1+|\xi|^{2}\right)^{1 / 2}, \quad s \in \mathbb{R}$,
we use the norm
$\|u\|_{H_{\hbar}^{s}\left(\mathbb{R}^{d}\right)}^{2}:=(2 \pi \hbar)^{-d} \int_{\mathbb{R}^{d}}\langle\xi\rangle^{2 s}\left|\mathcal{F}_{\hbar} u(\xi)\right|^{2} \mathrm{~d} \xi ;$
see $[52, \S 8.3],[14, \S E .1 .8]$. We abbreviate $H_{\hbar}^{s}\left(\mathbb{R}^{d}\right)$ to $H_{\hbar}^{s}$ and $L^{2}\left(\mathbb{R}^{d}\right)$ to $L^{2}$.

We record for later the fact that, by (4.3) and (4.4), for multiindices $\alpha$,

$$
\begin{align*}
\hbar^{|\alpha|}\left\|\partial^{\alpha} \phi\right\|_{L^{2}} & =\left\|(-\mathrm{i} \hbar \partial)^{\alpha} \phi\right\|_{L^{2}}=\frac{1}{(2 \pi \hbar)^{d / 2}}\left\|\xi^{\alpha} \mathcal{F}_{\hbar} \phi\right\|_{L^{2}} \\
& \leq \frac{1}{(2 \pi \hbar)^{d / 2}}\left\|\langle\xi\rangle^{|\alpha|} \mathcal{F}_{\hbar} \phi\right\|_{L^{2}}=\|\phi\|_{H_{\hbar}^{|\alpha|}} \tag{4.6}
\end{align*}
$$

Phase space The set of all possible positions $x$ and momenta (i.e. Fourier variables) $\xi$ is denoted by $T^{*} \mathbb{R}^{d}$; this is known informally as "phase space". Strictly, $T^{*} \mathbb{R}^{d}:=\mathbb{R}^{d} \times\left(\mathbb{R}^{d}\right)^{*}$, but for our purposes, we can consider $T^{*} \mathbb{R}^{d}$ as $\left\{(x, \xi): x \in \mathbb{R}^{d}, \xi \in \mathbb{R}^{d}\right\}$.

To deal with the behaviour of functions on phase space uniformly near $\xi=\infty$ (so-called fiber infinity), we consider the radial compactification in the $\xi$ variable of $T^{*} \mathbb{R}^{d}$. This is defined by
$\bar{T}^{*} \mathbb{R}^{d}:=\mathbb{R}^{d} \times B^{d}$,
where $B^{d}$ denotes the closed unit ball, considered as the closure of the image of $\mathbb{R}^{d}$ under the radial compactification map
$\mathrm{RC}: \xi \mapsto \xi /(1+\langle\xi\rangle) ;$
see $\left[14\right.$, §E.1.3]. Near the boundary of the ball, $|\xi|^{-1} \circ R C^{-1}$ is a smooth function, vanishing to first order at the boundary, with $\left(|\xi|^{-1} \circ \mathrm{RC}^{-1}, \widehat{\xi} \circ \mathrm{RC}^{-1}\right)$ thus giving local coordinates on the ball near its boundary. The boundary of the ball should be considered as a sphere at infinity consisting of all possible directions of the momentum variable. More generally, we denote $\bar{T}^{*} X:=X \times B^{d}$ for $X \subset \mathbb{R}^{d}$, and where appropriate (e.g., in dealing with finite values of $\xi$ only), we abuse notation by dropping the composition with RC from our notation and simply identifying $\mathbb{R}^{d}$ with the interior of $B^{d}$.

Symbols, quantisation, and semiclassical pseudodifferential operators A symbol is a function on $T^{*} \mathbb{R}^{d}$ that is also allowed to depend on $\hbar$, and thus can be considered as an $\hbar$-dependent family of functions. Such a family $a=\left(a_{\hbar}\right)_{0<\hbar \leq \hbar_{0}}$, with $a_{\hbar} \in C^{\infty}\left(T^{*} \mathbb{R}^{d}\right)$, is a symbol of order $m$, written as $a \in S^{m}\left(\mathbb{R}^{d}\right)$, if for any multiindices $\alpha, \beta$
$\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} a(x, \xi)\right| \leq C_{\alpha, \beta}\langle\xi\rangle^{m-|\beta|} \quad$ for all $(x, \xi) \in T^{*} \mathbb{R}^{d}$ and for all $0<\hbar \leq \hbar_{0}$,
where $C_{\alpha, \beta}$ does not depend on $\hbar$, $x$, or $\xi$; see [52, p. 207], [14, §E.1.2]. In this paper, we only consider these symbol classes on $\mathbb{R}^{d}$, and so we abbreviate $S^{m}\left(\mathbb{R}^{d}\right)$ to $S^{m}$.

For $a \in S^{m}$, we define the semiclassical quantisation of $a, \operatorname{Op}_{\hbar}(a)$ : $\mathcal{S}\left(\mathbb{R}^{d}\right) \rightarrow \mathcal{S}\left(\mathbb{R}^{d}\right)$, by
$\left(\mathrm{Op}_{\hbar}(a) v\right)(x):=(2 \pi \hbar)^{-d} \int_{\mathbb{R}^{d}} \int_{\mathbb{R}^{d}} \exp (\mathrm{i}(x-y) \cdot \xi / \hbar) a(x, \xi) v(y) \mathrm{d} y \mathrm{~d} \xi$
for $v \in \mathcal{S}\left(\mathbb{R}^{d}\right)$; [52, §4.1] [14, Page 543]. The integral in (4.8) need not converge, and can be understood either as an oscillatory integral in the
sense of [52, §3.6], [24, §7.8], or as an iterated integral, with the $y$ integration performed first; see [14, Page 543].

Conversely, if $A$ can be written in the form above, i.e. $A=\mathrm{Op}_{\hbar}(a)$ with $a \in S^{m}$, we say that $A$ is a semiclassical pseudo-differential operator of order $m$ and we write $A \in \Psi_{\hbar}^{m}$. We use the notation $a \in \hbar^{l} S^{m}$ if $\hbar^{-l} a \in S^{m}$; similarly $A \in \hbar^{l} \Psi_{\hbar}^{m}$ if $\hbar^{-l} A \in \Psi_{\hbar}^{m}$.

Theorem 4.1 (Composition and mapping properties of semiclassical pseudodifferential operators [52, Theorem 8.10], [14, Proposition E. 17 and Proposition E.19]). If $A \in \Psi_{\hbar}^{m_{1}}$ and $B \in \Psi_{\hbar}^{m_{2}}$, then
(i) $A B \in \Psi_{\hbar}^{m_{1}+m_{2}}$,
(ii) $[A, B]:=A B-B A \in \hbar \Psi_{\hbar}^{m_{1}+m_{2}-1}$,
(iii) For any $s \in \mathbb{R}, A$ is bounded uniformly in $\hbar$ as an operator from $H_{\hbar}^{s}$ to $H_{\hbar}^{s-m_{1}}$.

Residual class We say that $A=O\left(\hbar^{\infty}\right)_{\Psi-\infty}$ if, for any $s>0$ and $N \geq 1$, there exists $C_{s, N}>0$ so that
$\|A\|_{H_{\hbar}^{-s} \rightarrow H_{\hbar}^{s}} \leq C_{N, s} \hbar^{N} ;$
i.e. $A \in \Psi_{\hbar}^{-\infty}$ and furthermore all of its operator norms are bounded by any algebraic power of $\hbar$.

Principal symbol $\sigma_{\hbar}$ Let the quotient space $S^{m} / \hbar S^{m-1}$ be defined by identifying elements of $S^{m}$ that differ only by an element of $\hbar S^{m-1}$. For any $m$, there is a linear, surjective map
$\sigma_{\hbar}^{m}: \Psi_{\hbar}^{m} \rightarrow S^{m} / \hbar S^{m-1}$,
called the principal symbol map, such that, for $a \in S^{m}$,
$\sigma_{\hbar}^{m}\left(\mathrm{Op}_{\hbar}(a)\right)=a \quad \bmod \hbar S^{m-1} ;$
see [52, Page 213], [14, Proposition E.14] (observe that (4.10) implies that $\left.\operatorname{ker}\left(\sigma_{\hbar}^{m}\right)=\hbar \Psi_{\hbar}^{m-1}\right)$.

When applying the map $\sigma_{\hbar}^{m}$ to elements of $\Psi_{\hbar}^{m}$, we denote it by $\sigma_{\hbar}$ (i.e. we omit the $m$ dependence) and we use $\sigma_{\hbar}(A)$ to denote one of the representatives in $S^{m}$ (with the results we use then independent of the choice of representative). Key properties of the principal symbol that we use below are that
$\sigma_{\hbar}(A B)=\sigma_{\hbar}(A) \sigma_{\hbar}(B)$,
$\sigma_{\hbar}\left(P_{\hbar}\right)=\langle\mathrm{A} \xi, \xi\rangle-n$,
where $\langle\cdot, \cdot\rangle$ denotes the $\ell^{2}$ inner product on $\mathbb{R}^{d}$. The property (4.11) is proved in [14, Proposition E.17], (4.12) follows from (4.10) since $P_{\hbar}=\mathrm{Op}_{\hbar}\left(\langle\mathrm{A} \xi, \xi\rangle-n-\mathrm{i} \hbar \xi_{\ell} \partial_{j} A_{j \ell}\right)$ (where we sum over the indices $j$ and $\ell$ ).

Operator wavefront set $\mathrm{WF}_{\hbar}$ We say that $\left(x_{0}, \xi_{0}\right) \in \bar{T}^{*} \mathbb{R}^{d}$ is not in the semiclassical operator wavefront set of $A=\mathrm{Op}_{\hbar}(a) \in \Psi_{\hbar}^{m}$, denoted by $\mathrm{WF}_{\hbar} A$, if there exists a neighbourhood $U$ of $\left(x_{0}, \xi_{0}\right)$ such that for all multiindices $\alpha, \beta$ and all $N \geq 1$ there exists $C_{\alpha, \beta, U, N}>0$ (independent of $\hbar$ ) so that, for all $0<\hbar \leq \hbar_{0}$,
$\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} a(x, \xi)\right| \leq C_{\alpha, \beta, U, N} \hbar^{N}\langle\xi\rangle^{-N} \quad$ for all $(x, \mathrm{RC}(\xi)) \in U ;$
i.e. outside its semiclassical operator wavefront set an operator vanishes faster than any algebraic power of both $\hbar$ and $\langle\xi\rangle^{-1}$; see [52, Page 194], [14, Definition E.27]. Three properties of the semiclassical operator wavefront set that we use below are
$\mathrm{WF}_{\hbar}(A B) \subset \mathrm{WF}_{\hbar} A \cap \mathrm{WF}_{\hbar} B$
(see [52, §8.4], [14, E.2.5]),
$\mathrm{WF}_{\hbar}\left(\mathrm{Op}_{\hbar}(a)\right) \subset \operatorname{supp} a$
(since $(\operatorname{supp} a)^{c} \subset\left(\mathrm{WF}_{\hbar}\left(\mathrm{Op}_{\hbar}(a)\right)\right)^{c}$ by (4.13)), and
$\mathrm{WF}_{\hbar} A=\emptyset \Longleftrightarrow A=O\left(\hbar^{\infty}\right)_{\Psi-\infty}$
(see [14, E.2.2]).
Compactly-supported operators We say that $A$ is compactly supported if its Schwartz kernel is compactly supported in some set $K \Subset \mathbb{R}^{d} \times \mathbb{R}^{d}$, for all $0<\hbar \leq \hbar_{0}$. We recall that if $\mathcal{D}\left(\mathbb{R}^{d}\right):=C_{\text {comp }}^{\infty}\left(\mathbb{R}^{d}\right)$ (i.e. the set of test functions) and $\mathcal{D}^{\prime}\left(\mathbb{R}^{d}\right)$ denote the set of linear functionals on $\mathcal{D}\left(\mathbb{R}^{d}\right)$ (i.e. the set of distributions), given a bounded, sequentially-continuous operator $A: \mathcal{D} \rightarrow \mathcal{D}^{\prime}$ there exists a Schwartz kernel $\mathcal{K}_{A} \in \mathcal{D}^{\prime}\left(\mathbb{R}^{d} \times \mathbb{R}^{d}\right)$ such that
$A v(x)=\int_{\mathbb{R}^{d}} \mathcal{K}_{A}(x, y) v(y) \mathrm{d} y$,
in the sense of distributions; see, e.g., [24, Theorem 5.2.1], [14, §A.7]. We use below the facts that

- $A$ is compactly supported iff there exist $\chi_{1}, \chi_{2} \in \mathcal{D}$ such that $A=$ $\chi_{1} A \chi_{2}$, thus
- if $\chi_{1}, \chi_{2} \in \mathcal{D}$ are compactly supported functions, then $\chi_{1} A \chi_{2}$ is compactly supported, and
- if $P$ is a differential operator and $\chi \in \mathcal{D}$, then both $\chi P$ and $P \chi$ are compactly supported.

Ellipticity We say that $B \in \Psi_{\hbar}^{m}$ is elliptic on $X \subset \bar{T}^{*} \mathbb{R}^{d}$ if there exists $c>0$, independent of $\hbar$, such that
$\langle\xi\rangle^{-m}\left|\sigma_{\hbar}(B)(x, \xi)\right| \geq c, \quad$ for all $(x, \mathrm{RC}(\xi)) \in X$ and for all $0<\hbar \leq \hbar_{0}$.

A key feature of elliptic operators is that they are microlocally invertible; this is reflected in the following result.

Proposition 4.2 (Elliptic parametrix [14, Proposition E.32]). ${ }^{2}$ Let $A \in \Psi_{\hbar}^{m}$ and $B \in \Psi_{\hbar}^{\ell}$ be such that $B$ is elliptic on $\mathrm{WF}_{\hbar}(A)$. Then there exist $Q, Q^{\prime} \in$ $\Psi_{\hbar}^{m-\ell}$ such that

$$
A=B Q+O\left(\hbar^{\infty}\right)_{\Psi^{-\infty}}=Q^{\prime} B+O\left(\hbar^{\infty}\right)_{\Psi^{-\infty}}
$$

Theorem 4.3 (Elliptic estimate [14, Theorem E.33]). ${ }^{2}$ Let $A \in \Psi_{\hbar}^{m_{1}}, B_{1} \in$ $\Psi_{\hbar}^{m_{2}}$, and $P \in \Psi_{\hbar}^{\ell}$ be so that $B_{1} P$ is elliptic on $\mathrm{WF}_{\hbar}(A)$.
(i) Given $s, N>0$, and $M>0$, if $v \in \mathcal{D}^{\prime}$ and $B_{1} P v \in H^{s-m_{2}-\ell}$ then $A v \in$ $H^{s-m_{1}}$ and there exists $C_{s}>0, C_{N, M, s}>0$ (independent of $v$ and $\hbar$ ) such that

$$
\begin{equation*}
\|A v\|_{H_{\hbar}^{s-m_{1}}} \leq C_{s}\left\|B_{1} P v\right\|_{H_{\hbar}^{s-m_{2}-\ell}}+C_{N, M, s} \hbar^{M}\|v\|_{H_{\hbar}^{-N}} \tag{4.18}
\end{equation*}
$$

(ii) If, in addition, $A$ and $B_{1} P$ are compactly supported, then there exists $\tilde{\chi} \in C_{\text {comp }}^{\infty}$ so that

$$
\begin{equation*}
\|A v\|_{H_{\hbar}^{s-m_{1}}} \leq C_{s}\left\|B_{1} P v\right\|_{H_{\hbar}^{s-m_{2}-\ell}}+C_{N, M, s} \hbar^{M}\|\tilde{\chi} v\|_{H_{\hbar}^{-N}} \tag{4.19}
\end{equation*}
$$

Part (i) of Theorem 4.3 is proved by using Proposition 4.2 with $B=B_{1} P \in \Psi_{\hbar}^{m_{2}+\ell}$, applying the resulting operator equation to $v$, and taking norms. The operator $Q^{\prime} \in \Psi_{\hbar}^{m_{1}-m_{2}-\ell}$ and the constant $C_{s}$ is then $\left\|Q^{\prime}\right\|_{H_{\hbar}^{s-m_{2}-\ell} \rightarrow H_{\hbar}^{s-m_{1}}}$. The proof of Part (ii) is similar, using that, since $A$ and $B_{1} P$ are both compactly supported, there exists $\tilde{\chi} \in C_{\text {comp }}^{\infty}$ such that $\left(A-B_{1} P\right) v=\left(A-B_{1} P\right) \widetilde{\chi} v$.

[^2]
## 5. Proof of Theorem 3.1

In the notation introduced in §4, Theorem 3.1 becomes the following.

Theorem 5.1. Let A and $n$ satisfy Assumption 2.1 and let $R>0$ be such that $\operatorname{supp}(\mathrm{I}-\mathrm{A}) \cup \operatorname{supp}(1-n) \Subset B_{R}$. Given $f \in L^{2}\left(B_{R}\right)$, let u satisfy $P_{\hbar} u=\hbar^{2} f$ in $\mathbb{R}^{d}$ and the Sommerfeld radiation condition (1.2). Assume that, given $k_{0}>0, C_{\mathrm{sol}}(k)$ is polynomially bounded (in the sense of Definition 2.6) for $k \in K \subset\left[k_{0}, \infty\right)$. Given $k_{0}>0$, let $\hbar_{0}:=k_{0}^{-1}$, and let $H:=\left\{k^{-1}: k \in K\right\} \subset$ $\left(0, \hbar_{0}\right]$.

Then there exist $C_{3}, C_{4}, C_{5}>0$ such that
$\left.u\right|_{B_{R}}=u_{H^{2}}+u_{\mathcal{A}}$
where $u_{H^{2}} \in H_{\hbar}^{2}\left(B_{R}\right)$ with

$$
\begin{align*}
& \left\|\partial^{\alpha} u_{H^{2}}\right\|_{L^{2}\left(B_{R}\right)} \leq C_{3} \hbar^{2-|\alpha|}\|f\|_{L^{2}\left(B_{R}\right)} \\
& \quad \text { for all }|\alpha| \leq 2 \text { and for all } \hbar \in H \subset\left(0, \hbar_{0}\right] \tag{5.1}
\end{align*}
$$

and $u_{\mathcal{A}} \in C^{\infty}\left(B_{R}\right)$ with
$\left\|\partial^{\beta} u_{\mathcal{A}}\right\|_{L^{2}\left(B_{R}\right)} \leq C_{\text {sol }}\left(\hbar^{-1} ; R+2\right) C_{4}\left(\frac{\hbar}{C_{5}}\right)^{1-|\beta|}\|f\|_{L^{2}\left(B_{R}\right)}$
for all $\beta$ and for all $\hbar \in H \subset\left(0, \hbar_{0}\right]$,
where $C_{3}, C_{4}$, and $C_{5}$ depend on $\mathrm{A}, n, d$, and $\hbar_{0}$, but are independent of $\hbar$, $f, \alpha$, and $\beta$.

### 5.1. Step 0: restatement of bounds on the solution operator in semiclassical notation

The definition of $C_{\text {sol }}$ (Definition 2.5) implies that, in semiclassical notation,
$\|u\|_{H_{\hbar}^{1}\left(B_{R}\right)} \leq \hbar C_{\text {sol }}\left(\hbar^{-1}\right)\|f\|_{L^{2}\left(B_{R}\right)} \quad$ for all $\hbar>0$.
It is convenient to record here in semiclassical notation the bound on the solution operator when $C_{\text {sol }}$ is polynomially bounded.

Lemma 5.2 (Polynomial boundedness rewritten in terms of $\hbar$ ). Given $f \in$ $L_{\text {comp }}^{2}\left(\mathbb{R}^{d}\right)$, let $u \in H_{\text {loc }}^{1}\left(\mathbb{R}^{d}\right)$ be the solution to
$P_{\hbar} u=\hbar^{2} f$
satisfying the Sommerfeld radiation condition (1.2) (with $k=\hbar^{-1}$ ).
If $C_{\mathrm{sol}}(k)$ is polynomially bounded for $k \in K \subset\left[k_{0}, \infty\right)$ (in the sense of Definition 2.6), then there exists $M>0$ (independent of $\hbar$ ) such that, given $\chi \in C_{\mathrm{comp}}^{\infty}\left(\mathbb{R}^{d}\right)$, there exists $C>0$ (independent of $\hbar$ but dependent on $\chi$ ) such that
$\|\chi u\|_{L^{2}} \leq C \hbar^{1-M}\|f\|_{L^{2}} \quad$ for $\hbar \in H \subset\left(0, \hbar_{0}\right]$,
where $\hbar_{0}:=k_{0}^{-1}$ and $H:=\left\{k^{-1}: k \in K\right\}$.
The bound (5.4) also holds with $\|\chi u\|_{L^{2}}$ replaced by $\|\chi u\|_{H_{\hbar}^{1}}$, but we only need it in the form (5.4) for what follows.

### 5.2. Step 1: the definitions of $u_{\mathcal{A}}$ and $u_{H^{2}}$

The cut-off functions $\chi$ and $\chi_{\mu}$ Let $\chi \in C_{\text {comp }}^{\infty}\left(\mathbb{R}^{d} ;[0,1]\right)$ be such that
$\chi= \begin{cases}1 & \text { in } B_{1} \\ 0 & \text { outside } B_{2} .\end{cases}$
For $\mu>0$, let
$\chi_{\mu}(\cdot):=\chi\left(\frac{\dot{\mu}}{\mu}\right)$.


Fig. 5.1. The locations of $\mathrm{WF}_{\hbar}\left(\Pi_{H}\right)$ and $\mathrm{WF}_{\hbar}\left(\Pi_{L}\right)$, the regions where the principal symbols of $\Pi_{H}$ and $\Pi_{L}$ equal one, and the region where $P_{\hbar}$ is elliptic.

We define $\mu_{0}=\mu_{0}(\mathrm{~A}, n)$ by
$\mu_{0}(\mathrm{~A}, n):=\left(1+\frac{2 n_{\max }}{A_{\min }}\right)$.
The reason for this definition is that it implies that
if $\quad|\xi|^{2} \geq \mu_{0} \quad$ then $\quad\langle\xi\rangle^{-2} \sigma_{\hbar}(P) \geq \frac{A_{\min }}{2}>0$.
Indeed, by (4.12),
$\langle\xi\rangle^{-2} \sigma_{\hbar}(P) \geq \frac{A_{\min }|\xi|^{2}-n_{\max }}{1+|\xi|^{2}}=\frac{A_{\min }}{2}+\left(\frac{A_{\min }}{2}\right)\left(\frac{|\xi|^{2}-1-2 n_{\max } / A_{\min }}{1+|\xi|^{2}}\right)$,
and (5.8) follows. The importance of the property (5.8) is explained at the end of this subsection.

The frequency cut-offs $\Pi_{L}$ and $\Pi_{H}$ We define $\Pi_{L}$ and $\Pi_{H}$, the projections on low and high frequencies respectively, by (1.9) and (1.10). The definition of the quantisation $\mathrm{Op}_{\hbar}$ (4.8) and the change of variable $\zeta=\xi / \hbar$ imply that
$\Pi_{L}=\mathrm{Op}_{\hbar}\left(\chi_{\mu}\left(|\xi|^{2}\right)\right)$
and
$\Pi_{H}=I-\Pi_{L}$.
These definitions and the definition of $\Psi_{\hbar}^{m}\left(\mathbb{R}^{d}\right)$ in $\S 4$ imply that $\Pi_{L} \in$ $\Psi_{\hbar}^{-\infty}\left(\mathbb{R}^{d}\right)$ and $\Pi_{H} \in \Psi_{\hbar}^{0}\left(\mathbb{R}^{d}\right)$.

The locations of the wavefront sets of the frequency cut-offs, and the regions where their symbols equal one In Fig. 5.1 we show, as functions of $|\xi|^{2}$, the locations of $\mathrm{WF}_{\hbar}\left(\Pi_{H}\right)$ and $\mathrm{WF}_{\hbar}\left(\Pi_{L}\right)$, and the regions where $\sigma_{\hbar}\left(\Pi_{H}\right)$, and $\sigma_{\hbar}\left(\Pi_{L}\right)$ equal one. These locations/regions are obtained using (4.15) and (4.10) respectively. For example, since $1-\chi_{\mu}\left(|\xi|^{2}\right)=1$ for $|\xi|^{2} \geq 2 \mu$ and $=0$ for $|\xi|^{2} \leq \mu$, (4.10) and (4.15) imply that
$\sigma_{\hbar}\left(\Pi_{H}\right)=1$ on $\left\{\xi:|\xi|^{2} \geq 2 \mu\right\} \quad$ and $\quad \mathrm{WF}_{\hbar}\left(\Pi_{H}\right) \subset\left\{\xi:|\xi|^{2} \geq \mu\right\}$.

We also record the following key consequence of the results summarised in Fig. 5.1.

Lemma 5.3. If $\mu \geq \mu_{0}$, then $P_{\hbar}$ is elliptic on $\mathrm{WF}_{\hbar}\left(\Pi_{H}\right)$.
This property is central to our proof of the bound (5.1) on $u_{H^{2}}$, i.e., the high-frequency component. It is a consequence of (5.8), and the reason why we choose $\mu_{0}$ as in (5.7) is for this ellipticity result to hold.

The definitions of $u_{\mathcal{A}}$ and $u_{H^{2}}$ As described in $\S 1.2$, we choose $\varphi \in$ $C_{\text {comp }}^{\infty}\left(\mathbb{R}^{d}\right)$ be equal to one on $B_{R+1}$ and vanish outside $B_{R+2}$. We then let
$w:=\varphi u$
and we define
$u_{\mathcal{A}}:=\left.\left(\Pi_{L} w\right)\right|_{B_{R}} \quad$ and $\quad u_{H^{2}}:=\left.\left(\Pi_{H} w\right)\right|_{B_{R}}$.

### 5.3. Step 2: proof of the bound (5.2) on $u_{\mathcal{A}}$ (the low-frequency component)

Since $\Pi_{L} \in \Psi_{\hbar}^{-\infty}$, Part (iii) of Theorem 4.1, together with Sobolev embedding, gives $\Pi_{L} w \in C^{\infty}$.

The definition of $\Pi_{L}$ (1.9) and Plancherel's identity (4.4) for the standard (i.e. non semiclassical) Fourier transform imply that

$$
\begin{align*}
\left\|\partial^{\beta}\left(\Pi_{L} w\right)\right\|_{L^{2}} & =\frac{1}{(2 \pi)^{d / 2}}\left\|(\cdot)^{\beta} \mathcal{F}\left(\Pi_{L} w\right)(\cdot)\right\|_{L^{2}} \\
& =\frac{1}{(2 \pi)^{d / 2}}\left\|(\cdot)^{\beta} \chi_{\mu}\left(\hbar^{2}|\cdot|^{2}\right) \mathcal{F} w(\cdot)\right\|_{L^{2}} \tag{5.12}
\end{align*}
$$

The definitions of $\chi$ (5.5) and $\chi_{\mu}$ (5.6) imply that $\chi_{\mu}(\xi)=0$ for $|\xi| \geq 2 \mu$, so
$\chi_{\mu}\left(\hbar^{2}|\zeta|^{2}\right)=0 \quad$ for $|\zeta| \geq \sqrt{2 \mu} \hbar^{-1}$.
Using this fact, and then (in this order) the fact that $\left|\chi_{\mu}\right| \leq 1$, Plancherel's identity for the standard Fourier transform, the fact that $\varphi=0$ outside $B_{R+2}$, and the definition of $C_{\text {sol }}$ (2.12), we find from (5.12) that

$$
\begin{aligned}
\left\|\partial^{\beta}\left(\Pi_{L} \varphi u\right)\right\|_{L^{2}} & \leq \frac{(2 \mu)^{|\beta| / 2}}{(2 \pi)^{d / 2}} \hbar^{-|\beta|}\left\|\chi_{\mu}\left(\hbar^{2}|\cdot|^{2}\right) \mathcal{F}(\varphi u)(\cdot)\right\|_{L^{2}} \\
& \leq \frac{(2 \mu)^{|\beta| / 2}}{(2 \pi)^{d / 2}} \hbar^{-|\beta|}\|\mathcal{F}(\varphi u)\|_{L^{2}} \\
& \leq(2 \mu)^{|\beta| / 2} \hbar^{-|\beta|}\|\varphi u\|_{L^{2}} \\
& \leq(2 \mu)^{|\beta| / 2} \hbar^{-|\beta|} \hbar C_{\mathrm{sol}}\left(\hbar^{-1} ; R+2\right)\|f\|_{L^{2}\left(B_{R}\right)}
\end{aligned}
$$

Since
$\left\|\partial^{\beta} u_{\mathcal{A}}\right\|_{L^{2}\left(B_{R}\right)}=\left\|\partial^{\beta}\left(\Pi_{L} w\right)\right\|_{L^{2}\left(B_{R}\right)} \leq\left\|\partial^{\beta}\left(\Pi_{L} w\right)\right\|_{L^{2}}$,
the bound (5.2) then follows with $C_{4}:=\sqrt{2 \mu}$ and $C_{5}:=\sqrt{2 \mu}$.

### 5.4. Step 3: proof of the bound (5.1) on $u_{H^{2}}$ (the high-frequency component)

By the inequality (4.6), it is sufficient to prove that

$$
\begin{equation*}
\left\|\Pi_{H} w\right\|_{H_{\hbar}^{2}} \leq C_{3} \hbar^{2}\|f\|_{L^{2}\left(B_{R}\right)} \quad \text { for all } \hbar \in H \subset\left(0, \hbar_{0}\right) \tag{5.13}
\end{equation*}
$$

It is instructive to first prove (5.13) under the assumption that $C_{\text {sol }}(k) \lesssim 1$ (which, by Theorem 2.7 is ensured if A and $n$ are nontrapping). Indeed, as discussed in $\S 1.2$, this proof only requires that $P_{\hbar}$ is elliptic on $\mathrm{WF}_{\hbar}\left(\Pi_{H}\right)$; i.e., Lemma 5.3. Throughout the rest of this section, therefore, we assume that $\mu \geq \mu_{0}$, so that the result of Lemma 5.3 holds.

### 5.4.1. Proof of $(5.13)$ under the assumption that $C_{\text {sol }}(k) \lesssim 1$

We seek to apply Part (i) of Theorem 4.3 with $A=\Pi_{H}$ (so $m_{1}=0$ ), $B_{1}=1$ (so $m_{2}=0$ ), and $P=P_{\hbar}$ (so $\ell=2$ ). By Lemma 5.3, $B_{1} P$ is elliptic on $\mathrm{WF}_{\hbar}(A)$. We can therefore apply Theorem 4.3 and obtain that, given $N, N^{\prime}>0$,
$\left\|\Pi_{H} w\right\|_{H_{\hbar}^{2}} \lesssim\left\|P_{\hbar} w\right\|_{L^{2}}+\hbar^{N^{\prime}}\|w\|_{H_{\hbar}^{-N}}$,
where the omitted constant in $\lesssim$ depends on $N$ and $N^{\prime}$. Since $P_{\hbar} u=\hbar^{2} f$,
$P_{\hbar} w=\left[P_{\hbar}, \varphi\right] u+\hbar^{2} \varphi f$,
where $[\cdot, \cdot]$ is the standard commutator defined by $\left[A_{1}, A_{2}\right]:=A_{1} A_{2}-$ $A_{2} A_{1}$, so that (5.14) becomes
$\left\|\Pi_{H} w\right\|_{H_{\hbar}^{2}} \lesssim\left\|\left[P_{\hbar}, \varphi\right] u\right\|_{L^{2}}+\hbar^{2}\|f\|_{L^{2}}+\hbar^{N^{\prime}}\|w\|_{H_{\hbar}^{-N}}$.
Direct calculation, using the fact that $\operatorname{supp} \varphi \subset B_{R+2}$, implies that
$\left\|\left[P_{\hbar}, \varphi\right] u\right\|_{L^{2}} \lesssim \hbar\|u\|_{H_{\hbar}^{1}\left(B_{R+2}\right)}$,
where the omitted constant depends on $\varphi$, and hence on $R$.

Combining (5.15) and (5.16), and recalling that $\operatorname{supp} \varphi \subset B_{R+2}$, we have
$\left\|\Pi_{H} w\right\|_{H_{\hbar}^{2}} \lesssim \hbar\|u\|_{H_{\hbar}^{1}\left(B_{R+2}\right)}+\hbar^{2}\|f\|_{L^{2}\left(B_{R}\right)}+\hbar^{N^{\prime}}\|u\|_{H_{\hbar}^{-N}\left(B_{R+2}\right)}$.
Choosing $N=0$ and $N^{\prime}=1$, and then using (5.3), we obtain
$\left\|\Pi_{H} w\right\|_{H_{\hbar}^{2}} \lesssim \hbar^{2}\left(1+C_{\text {sol }}\left(\hbar^{-1}\right)\right)\|f\|_{L^{2}\left(B_{R}\right)}$.
If $C_{\text {sol }}\left(\hbar^{-1}\right) \lesssim 1$, then this implies (5.13). However, if $C_{\text {sol }}\left(\hbar^{-1}\right) \gg 1$ (as occurs when $C_{\text {sol }}$ is polynomially bounded in the sense of Definition 2.6 with $M>0$ ) then (5.17) is a weaker bound than (5.13).

### 5.4.2. Proof of $(5.13)$ under the assumption that $C_{\text {sol }}(k)$ is polynomially bounded

Inspecting the argument in §5.4.1, we see that the assumption that $C_{\text {sol }}(k) \lesssim 1$ is needed to get a good bound on the commutator term [ $\left.P_{\hbar}, \varphi\right] u$. To remove this commutator term, one idea is to use the elliptic estimate in Part (i) of Theorem 4.3, using the fact that $P_{\hbar}$ is elliptic on $\mathrm{WF}_{\hbar}\left(\Pi_{H} \varphi\right)$, and apply the estimate with $v:=u$. However, the error term would not be compactly supported and we would be unable to control it using the polynomial bound on the solution operator (5.4). We therefore introduce additional spatial cut-offs on the left of $\Pi_{H} \varphi$ and $P_{\hbar}$ to create compactly-supported operators and have a compactly-supported error term thanks to Part (ii) of Theorem 4.3.

To this end, let $\varphi_{1}, \varphi_{2} \in C_{\text {comp }}^{\infty}\left(\mathbb{R}^{d}\right)$ be such that $\varphi_{1}=1$ on $\operatorname{supp} \varphi$ and $\varphi_{2}=1$ on $\operatorname{supp} \varphi_{1}$; we then write
$\Pi_{H} \varphi u=\left(1-\varphi_{1}\right) \Pi_{H} \varphi u+\varphi_{1} \Pi_{H} \varphi u$.
Since $1-\varphi_{1}=0$ on $\operatorname{supp} \varphi$, using (4.14) and (4.15), we obtain that
$\mathrm{WF}_{\hbar}\left(\left(1-\varphi_{1}\right) \Pi_{H} \varphi\right) \subset \bar{T}^{*}\left(\operatorname{supp}\left(1-\varphi_{1}\right)\right) \cap \bar{T}^{*}(\operatorname{supp} \varphi)=\emptyset$.
Hence, by (4.16), $\left(1-\varphi_{1}\right) \Pi_{H} \varphi=O\left(\hbar^{\infty}\right)_{\Psi^{-\infty}}$, and, by the definition of the residual class (4.9), for any $N \geq 1$ there exists $C_{N}>0$ so that
$\left\|\left(1-\varphi_{1}\right) \Pi_{H} \varphi u\right\|_{H_{\hbar}^{2}}=\left\|\left(1-\varphi_{1}\right) \Pi_{H} \varphi \varphi_{1} u\right\|_{H_{\hbar}^{2}} \leq C_{N} \hbar^{N}\left\|\varphi_{1} u\right\|_{L^{2}}$,
were we used the fact that $\varphi_{1}=1$ on $\operatorname{supp} \varphi$ in the first equality.
It therefore remains to control $\varphi_{1} \Pi_{H} \varphi u$; to do this, we use the elliptic estimate of Theorem 4.3.

Lemma 5.4. $\varphi_{2} P_{\hbar}$ is elliptic on $\mathrm{WF}_{\hbar}\left(\varphi_{1} \Pi_{H} \varphi\right)$.

Proof. By (4.14) and (4.15), $\mathrm{WF}_{\hbar}\left(\varphi_{1} \Pi_{H} \varphi\right) \subset \bar{T}^{*}\left(\operatorname{supp} \varphi_{1}\right) \cap \mathrm{WF}_{\hbar} \Pi_{H}$. Since $\varphi_{2}=1$ on $\operatorname{supp} \varphi_{1}$, the result is a direct consequence of Lemma 5.3.

By the facts about compactly-supported operators recalled in $\S 4$, $\varphi_{1} \Pi_{H} \varphi$ and $\varphi_{2} P_{\hbar}$ are compactly supported. Therefore, by Lemma 5.4, we can apply Part (ii) of Theorem 4.3 with $A=\varphi_{1} \Pi_{H} \varphi, B_{1}=\varphi_{2}, P=P_{\hbar}$, $m_{1}=0, m_{2}=0, \ell=2$. This result implies that there exists $\tilde{\chi} \in C_{\mathrm{comp}}^{\infty}$, and, for any $N^{\prime} \geq 1$, there exists $C_{N^{\prime}}>0$ such that

$$
\begin{align*}
\left\|\varphi_{1} \Pi_{H} \varphi u\right\|_{H_{\hbar}^{2}} & \lesssim\left\|\varphi_{2} P_{\hbar} u\right\|_{L^{2}}+C_{N^{\prime}} \hbar^{N^{\prime}}\|\tilde{\chi} u\|_{L^{2}} \\
& =\hbar^{2}\left\|\varphi_{2} f\right\|_{L^{2}}+C_{N^{\prime}} \hbar^{N^{\prime}}\|\tilde{\chi} u\|_{L^{2}} . \tag{5.20}
\end{align*}
$$

Collecting (5.18), (5.19), (5.20), using (5.4), and choosing $N=N^{\prime}=$ $M+1$, we obtain (5.13).

## 6. Proof of Theorem 3.4

The two ingredients for the proof of Theorem 3.4 are

- Lemma 6.4 , which is the standard duality argument giving a condition for quasi-optimality to hold in terms of how well the solution of the adjoint problem is approximated by the finite-element space (measured by the quantity $\eta\left(V_{N}\right)$ defined by (6.3)), and
- Lemma 6.5 that bounds $\eta\left(V_{N}\right)$ using the decomposition from Theorem 3.1.

Regarding Lemma 6.4: we recall that this argument came out of ideas introduced in [43], was then formalised in [42], and has been used extensively in the analysis of the Helmholtz FEM; see, e.g., [1,26,33,42, $35,36,51,50,12,9,31,10,17,21,16]$.

Before stating Lemma 6.4 we need to introduce some notation.
Definition 6.1 (The adjoint sesquilinear form $a^{*}(\cdot, \cdot)$ ). The adjoint sesquilinear form, $a^{*}(u, v)$, to the sesquilinear form $a(\cdot, \cdot)$ defined in (2.4) is given by
$a^{*}(u, v):=\overline{a(v, u)}=\int_{B_{R}}\left((\mathrm{~A} \nabla u) \cdot \overline{\nabla v}-k^{2} n u \bar{v}\right)-\left\langle\gamma u, \mathrm{DtN}_{k}(\gamma v)\right\rangle_{\partial B_{R}}$.
A key role is played by the solution operator of the adjoint variational problem with data in $L^{2}\left(B_{R}\right)$; we therefore introduce the following notation.

Definition 6.2 (Adjoint solution operator $S^{*}$ ). Given $f \in L^{2}\left(B_{R}\right)$, let $S^{*} f$ be defined as the solution of the variational problem
find $S^{*} f \in H^{1}\left(B_{R}\right) \quad$ such that
$a^{*}\left(S^{*} f, v\right)=\int_{B_{R}} f \bar{v} \quad$ for all $v \in H^{1}\left(B_{R}\right)$
Green's second identity applied to solutions of the Helmholtz equation satisfying the Sommerfeld radiation condition (1.2) implies that $\left\langle\operatorname{DtN}_{k} \psi, \bar{\phi}\right\rangle_{\partial B_{R}}=\left\langle\operatorname{DtN}_{k} \phi, \bar{\psi}\right\rangle_{\partial B_{R}}$ (see, e.g., [46, Lemma 6.13]); thus $a(\bar{v}, u)=a(\bar{u}, v)$ and so the definition (6.1) implies that
$a\left(\overline{S^{*} f}, v\right)=(\bar{f}, v)_{L^{2}\left(B_{R}\right)} \quad$ for all $v \in H^{1}\left(B_{R}\right)$.
Definition $6.3\left(\eta\left(V_{N}\right)\right)$. Given a sequence of finite-dimensional spaces $\left(V_{N}\right)_{N=0}^{\infty}$ (as described in §2.3), let
$\eta\left(V_{N}\right):=\sup _{0 \neq f \in L^{2}\left(B_{R}\right)^{v_{N}}} \min _{N} \in V_{N} \frac{\left\|S^{*} f-v_{N}\right\|_{H_{k}^{1}\left(B_{R}\right)}}{\|f\|_{L^{2}\left(B_{R}\right)}}$.
Lemma 6.4 (Conditions for quasi-optimality). If
$k \eta\left(V_{N}\right) \leq \frac{1}{C_{\text {cont }}} \sqrt{\frac{A_{\min }}{2\left(n_{\max }+A_{\min }\right)}}$,
then the Galerkin equations (2.15) have a unique solution which satisfies
$\left\|u-u_{h}\right\|_{H_{k}^{1}\left(B_{R}\right)} \leq \frac{2 C_{\mathrm{cont}}}{A_{\min }}\left(\min _{v_{N} \in V_{N}}\left\|u-v_{N}\right\|_{H_{k}^{1}\left(B_{R}\right)}\right)$.
Proof. Using the inequality (2.10), we see that $a(\cdot, \cdot)$ satisfies the Gårding inequality
$\Re(a(v, v)) \geq A_{\min }\|v\|_{H_{k}^{1}\left(B_{R}\right)}^{2}-2 k^{2}\left(n_{\max }+A_{\min }\right)\|v\|_{L^{2}\left(B_{R}\right)}^{2}$
and the result follows from, e.g., the account [46, Theorem 6.32] of the standard duality argument with (in the notation of [46]) $\alpha=A_{\text {min }}$ and $C_{v}=2 k^{2}\left(n_{\max }+A_{\min }\right)$.

Lemma 6.5 (Bound on $\eta\left(V_{N}\right)$ using the decomposition from Theorem 3.1). Let A and $n$ satisfy Assumption 2.1 and let $R>0$ be such that $\operatorname{supp}(\mathrm{I}-\mathrm{A}) \cup$ $\operatorname{supp}(1-n) \Subset B_{R}$. Let $\left(V_{N}\right)_{N=0}^{\infty}$ be the piecewise-polynomial approximation spaces described in [35, §5]. There exists $C_{6}, C_{7}, \sigma>0$, all independent of $k, h$, and $p$, such that

$$
\begin{align*}
k \eta\left(V_{N}\right) \leq & C_{6} C_{3} \frac{h k}{p}\left(1+\frac{k h}{p}\right) \\
& +C_{7} C_{\mathrm{sol}}(k)\left[\left(\frac{h}{h+\sigma}\right)^{p}\left(1+\frac{h k}{h+\sigma}\right)+k\left(\frac{k h}{\sigma p}\right)^{p}\left(\frac{1}{p}+\frac{k h}{\sigma p}\right)\right] \tag{6.5}
\end{align*}
$$

The constants $C_{6}$ and $\sigma$ only depend on the constants in [35, Assumption 5.2] defining the element maps from the reference element; $C_{7}$ depends on these constants, and additionally on $C_{5}$.

Proof. This proof is very similar to the proof of [35, Theorem 5.5]. Indeed, [35, Theorem 5.5] proves a bound very similar to (6.5) starting from bounds almost identical to the bounds (3.1) and (3.2) (recalling Remark 3.3 about notation). The only difference is that the bound (3.2) contains $C_{\text {sol }}$, which depends on $k$ (whereas in [35] $C_{\text {sol }} \sim 1$ ), and so we now need to keep track of how $C_{\text {sol }}$ enters the proof of [35, Theorem 5.5].

From the definition (6.3), it is sufficient to show that, given $f \in$ $L^{2}\left(B_{R}\right)$, there exists $w_{N} \in V_{N}$ such that
$\left\|S^{*} f-w_{N}\right\|_{H_{k}^{1}\left(B_{R}\right)} \leq C\|f\|_{L^{2}\left(B_{R}\right)}$,
where $C$ is the right-hand side of (6.5) divided by $k$. Let $v:=S^{*} f$; by (6.2) and Part (i) of Lemma 2.3, $\bar{v}$ satisfies the assumptions of Theorem 3.1 with $f$ replaced by $\bar{f}$, and so the bounds (3.1) and (3.2) hold with $u$ replaced by $v$.

By [35, First equation on Page 1896] (which uses [35, Theorem B.4]), the bound (3.6) holds, and thus there exists $w_{N}^{(1)} \in V_{N}$ such that $\left\|v_{H^{2}}-w_{N}^{(1)}\right\|_{H_{k}^{1}\left(B_{R}\right)} \leq C_{6} \frac{h}{p}\left(1+\frac{k h}{p}\right)|v|_{H^{2}\left(B_{R}\right)}$
and so
$\left\|v_{H^{2}}-w_{N}^{(1)}\right\|_{H_{k}^{1}\left(B_{R}\right)} \leq C_{6} \frac{h}{p}\left(1+\frac{k h}{p}\right) C_{3}\|f\|_{L^{2}\left(B_{R}\right)}$
by (3.1).
For the approximation of $v_{\mathcal{A}}$, the only change to the argument in [35] is that a multiplicative factor of $\left(C_{\text {sol }}\right)^{2}$ must be included on the right-hand side of [35, Equation 5.8]. Then [35, Equations 5.8 and 5.9] implies that there exists $C_{7}$ and $w_{N}^{(2)} \in V_{N}$ such that

$$
\begin{align*}
& k\left\|v_{\mathcal{A}}-w_{N}^{(2)}\right\|_{H_{k}^{1}\left(B_{R}\right)} \\
& \quad \leq C_{7} C_{\mathrm{sol}}(k)\left[\left(\frac{h}{h+\sigma}\right)^{p}\left(1+\frac{h k}{h+\sigma}\right)+k\left(\frac{k h}{\sigma p}\right)^{p}\left(\frac{1}{p}+\frac{k h}{\sigma p}\right)\right]\|f\|_{L^{2}\left(B_{R}\right)} \tag{6.8}
\end{align*}
$$

(observe that this equation is identical to [35, Last equation on Page 1896] except for the factor $C_{\text {sol }}$ on the right-hand side).

Let $w_{N}:=w_{N}^{(1)}+w_{N}^{(2)}$. By the triangle inequality, the decomposition $v=v_{H^{2}}+v_{\mathcal{A}}$ on $B_{R}$, and the inequalities (6.7) and (6.8), the inequality (6.6) holds with $C$ the right-hand side of (6.5) and the proof is complete.

Corollary 6.6 (Conditions under which $k \eta\left(V_{N}\right)$ is arbitrarily small). Let the assumptions of Lemma 6.5 hold. Given $\varepsilon>0$ and $k_{0}>0$, there exists $C_{1}, C_{2}>0$, depending only on $\varepsilon, C_{3}, C_{6}, C_{7}, \sigma$, and $k_{0}$, such that if
$\frac{h k}{p} \leq C_{1} \quad$ and $\quad p \geq C_{2}\left(1+\log k+\log \left(C_{\text {sol }}(k)\right)\right)$,
then
$k \eta\left(V_{N}\right) \leq \varepsilon \quad$ for all $k \geq k_{0}$.
Proof. This proof is essentially identical to the proofs of [35, Corollary 5.6] and [36, Theorem 5.8]. First choose $C_{1}$ sufficiently small such that $C_{1}<\sigma$ and

$$
C_{6} C_{3} C_{1}\left(1+C_{1}\right) \leq \frac{\varepsilon}{2}
$$

From the bound on $k \eta\left(V_{N}\right)$ (6.5), it is then sufficient to show that
$C_{7} C_{\mathrm{sol}}(k)\left[\left(\frac{h}{h+\sigma}\right)^{p}\left(1+\frac{h k}{h+\sigma}\right)+k\left(\frac{k h}{\sigma p}\right)^{p}\left(\frac{1}{p}+\frac{k h}{\sigma p}\right)\right]$
can be made $\leq \varepsilon / 2$. Let
$\theta_{1}:=\frac{h}{h+\sigma} \quad$ and $\quad \theta_{2}:=\frac{C_{1}}{\sigma}$,
so that (6.9) is bounded by
$C_{7} C_{\mathrm{sol}}(k)\left[\left(\theta_{1}\right)^{p}\left(1+\frac{C_{1} p}{\sigma}\right)+k\left(\theta_{2}\right)^{p}\left(\frac{1}{p}+\frac{C_{1}}{\sigma}\right)\right] ;$
the result then follows since $\theta_{1}, \theta_{2}<1$.

Proof of Theorem 3.4. This follows by combining Lemma 6.4 and Corollary 6.6.

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[^1]:    ${ }^{1}$ The counterpart of "semiclassical" involving differential/pseudodifferential operators without a small parameter is usually called "homogeneous" (owing to the homogeneity of the principal symbol) rather than "classical." "Classical" describes the behaviour in either calculus in the small- $\hbar$ or high-frequency limit respectively, where commutators of operators become Poisson brackets of symbols, hence classical particle dynamics replaces wave motion.

[^2]:    ${ }^{2}$ We highlight that working in $\mathbb{R}^{d}$ (as opposed to on a general manifold defined by coordinate charts) allows us to remove the proper-support assumption appearing in [14, Proposition E.32, Theorem E.33].

