# Functional calculus of Dirac operators and complex perturbations of Neumann and Dirichlet problems 

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#### Abstract

We prove that Neumann, Dirichlet and regularity problems for divergence form elliptic equations in the half-space are well posed in $L_{2}$ for small complex $L_{\infty}$ perturbations of a coefficient matrix which is either real symmetric, of block form or constant. All matrices are assumed to be independent of the transversal coordinate. We solve the Neumann, Dirichlet and regularity problems through a new boundary operator method which makes use of operators in the functional calculus of an underlaying first order Dirac type operator. We establish quadratic estimates for this Dirac operator, which implies that the associated Hardy projection operators are bounded and depend continuously on the coefficient matrix. We also prove that certain transmission problems for $k$-forms are well posed for small perturbations of block matrices.


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## 1. Introduction

In this paper we prove that the Neumann, Dirichlet and regularity problems are well posed in $L_{2}\left(\mathbf{R}^{n}\right)$ for divergence form second order elliptic equations

$$
\begin{equation*}
\operatorname{div}_{t, x} A(x) \nabla_{t, x} U(t, x)=0 \tag{1.1}
\end{equation*}
$$

on the half-space $\mathbf{R}_{+}^{n+1}:=\left\{(t, x) \in \mathbf{R} \times \mathbf{R}^{n} ; t>0\right\}, n \geqslant 1$, when $A$ is a small complex $L_{\infty}$ perturbation of either a block matrix, a constant matrix or a real symmetric matrix. Furthermore, the matrix $A=\left(a_{i j}(x)\right)_{i, j=0}^{n} \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\mathbf{C}^{n+1}\right)\right)$ is assumed to be $t$-independent with complex coefficients and accretive, with quantitative bounds $\|A\|_{\infty}$ and $\kappa_{A}$, where $\kappa_{A}>0$ is the largest constant such that

$$
\operatorname{Re}(A(x) v, v) \geqslant \kappa_{A}|v|^{2}, \quad \text { for all } v \in \mathbf{C}^{n+1}, x \in \mathbf{R}^{n}
$$

We shall approach Eq. (1.1) from a first order point of view, rewriting it as the first order system

$$
\left\{\begin{array}{l}
\operatorname{div}_{t, x} A(x) F(t, x)=0  \tag{1.2}\\
\operatorname{curl}_{t, x} F(t, x)=0
\end{array}\right.
$$

where $F(t, x)=\nabla_{t, x} U(t, x)$. Recall that a vector field $F=F_{0} e_{0}+F_{1} e_{1}+\cdots+F_{n} e_{n}$ can be written in this way as a gradient if and only if $\operatorname{curl}_{t, x} F=0$, by which we understand that $\partial_{j} F_{i}=\partial_{i} F_{j}$, for all $i, j=0, \ldots, n$. We write $\left\{e_{0}, e_{1}, \ldots, e_{n}\right\}$ for the standard basis for $\mathbf{R}^{n+1}$ with $e_{0}$ upward pointing into $\mathbf{R}_{+}^{n+1}$, and write $t=x_{0}$ for the vertical coordinate. For the vertical derivative, we write $\partial_{0}=\partial_{t}$. Denote also by $F_{\|}:=F_{1} e_{1}+\cdots+F_{n} e_{n}$, the tangential part of $F$, and write $\operatorname{curl}_{x} F_{\|}=0$ if $\partial_{j} F_{i}=\partial_{i} F_{j}$, for all $i, j=1, \ldots, n$.

In the formulation of the boundary value problems below, we assume that $A=A(x)$ is a given coefficient matrix with properties as above. Furthermore, by saying that $F_{t}(x)=F(t, x)$
satisfies (1.2) we shall mean that $F_{t} \in C^{1}\left(\mathbf{R}_{+} ; L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}^{n+1}\right)\right)$ and that for each fixed $t>0$, we have $\operatorname{div}_{x}(A F)_{\|}=-\left(A(x) \partial_{0} F\right)_{0}, \nabla_{x} F_{0}=\partial_{0} F_{\|}$and $\operatorname{curl}_{x} F_{\|}=0$, where the derivatives on the left-hand sides are taken in the sense of distributions. If this holds for $F$, then in particular we can write $F=\nabla_{t, x} U$, with $U \in W_{2, \text { loc }}^{1}\left(\mathbf{R}_{+}^{n+1}\right)$, and we see that $U$ satisfies (1.1) in the sense that

$$
\iint_{\mathbf{R}_{+}^{n+1}}\left(A(x) \nabla_{t, x} U(t, x), \nabla_{t, x} \varphi(t, x)\right) d t d x=0, \quad \text { for all } \varphi \in C_{0}^{\infty}\left(\mathbf{R}_{+}^{n+1}\right)
$$

Neumann problem (Neu-A). Given a function $\phi(x) \in L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}\right)$, find a vector field $F_{t}(x)=$ $F(t, x)$ in $\mathbf{R}_{+}^{n+1}$ such that $F_{t} \in C^{1}\left(\mathbf{R}_{+} ; L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}^{n+1}\right)\right)$ and $F$ satisfies (1.2) for $t>0$, and furthermore $\lim _{t \rightarrow \infty} F_{t}=0$ and $\lim _{t \rightarrow 0} F_{t}=f$ in $L_{2}$ norm, where the conormal part of $f$ satisfies the boundary condition

$$
e_{0} \cdot(A f)=\sum_{j=0}^{n} A_{0 j} f_{j}=\phi, \quad \text { on } \mathbf{R}^{n}=\partial \mathbf{R}_{+}^{n+1}
$$

For $U$, this means that the conormal derivative $\frac{\partial U}{\partial v_{A}}(0, x)=\phi(x)$ in $L_{2}\left(\mathbf{R}^{n}\right)$.
Regularity problem (Reg- $\boldsymbol{A}$ ). Given a function $\psi: \mathbf{R}^{n} \rightarrow \mathbf{C}$ with tangential gradient $\nabla_{x} \psi \in$ $L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}^{n}\right)$, find a vector field $F_{t}(x)=F(t, x)$ in $\mathbf{R}_{+}^{n+1}$ such that $F_{t} \in C^{1}\left(\mathbf{R}_{+} ; L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}^{n+1}\right)\right)$ and $F$ satisfies (1.2) for $t>0$, and furthermore $\lim _{t \rightarrow \infty} F_{t}=0$ and $\lim _{t \rightarrow 0} F_{t}=f$ in $L_{2}$ norm, where the tangential part of $f$ satisfies the boundary condition

$$
f_{\|}=f_{1} e_{1}+\cdots+f_{n} e_{n}=\nabla_{x} \psi, \quad \text { on } \mathbf{R}^{n}=\partial \mathbf{R}_{+}^{n+1}
$$

For $U$, this means that $\nabla_{x} U(x, 0)=\nabla_{x} \psi(x)$, i.e. $U(x, 0)=\psi(x)$ in $\dot{W}_{2}^{1}\left(\mathbf{R}^{n}\right)$.
Dirichlet problem (Dir- $\boldsymbol{A}$ ). Given a function $u(x) \in L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}\right)$, find a function $U_{t}(x)=U(t, x)$ in $\mathbf{R}_{+}^{n+1}$ such that $U_{t} \in C^{2}\left(\mathbf{R}_{+} ; L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}\right)\right), \nabla_{t, x} U_{t} \in C^{1}\left(\mathbf{R}_{+} ; L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}^{n+1}\right)\right)$ and $\nabla_{t, x} U_{t}$ satisfies (1.2) for $t>0$, and furthermore $\lim _{t \rightarrow \infty} U_{t}=0, \lim _{t \rightarrow \infty} \nabla_{t, x} U_{t}=0$ and $\lim _{t \rightarrow 0} U_{t}=u$ in $L_{2}$ norm.

We shall also use first order methods based on (1.2) to solve the Dirichlet problem. However, here we use a different relation between (1.1) and (1.2) which we now describe. Assume $F(t, x)=\sum_{k=0}^{n} F_{i}(t, x) e_{i}$ is a vector field satisfying (1.2). Applying $\partial_{t}$ to the first equation and using that $\operatorname{curl}_{t, x} F=0$ yields

$$
0=\partial_{t}\left(\operatorname{div}_{t, x} A(x) F\right)=\operatorname{div}_{t, x} A(x)\left(\partial_{t} F\right)=\operatorname{div}_{t, x} A(x)\left(\nabla_{t, x} F_{0}\right),
$$

since the coefficients are assumed to be $t$-independent. Thus the normal component $U:=F_{0}$ satisfies (1.1). Note that when $A=I$, the functions $F_{1}, \ldots, F_{n}$ are conjugates to $U$ in the sense of Stein and Weiss [25]. From this we see that solvability of (Dir- $A$ ) is a direct consequence of solvability of the following auxiliary Neumann problem.

Neumann problem ( $\left.\mathbf{N e u}^{\perp}-\boldsymbol{A}\right)$. Given a function $\phi(x) \in L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}\right)$, find a vector field $F_{t}(x)=$ $F(t, x)$ in $\mathbf{R}_{+}^{n+1}$ such that $F_{t} \in C^{1}\left(\mathbf{R}_{+} ; L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}^{n+1}\right)\right)$ and $F$ satisfies (1.2) for $t>0$, and
furthermore $\lim _{t \rightarrow \infty} F_{t}=0$ and $\lim _{t \rightarrow 0} F_{t}=f$ in $L_{2}$ norm, where the normal part of $f$ is $e_{0} \cdot f=f_{0}=\phi$.

The main result of this paper is the following $L_{\infty}$ perturbation result for the boundary value problems.

Theorem 1.1. Let $A_{0}(x)=\left(\left(a_{0}\right)_{i j}(x)\right)_{i, j=0}^{n}$ be a $t$-independent, complex, accretive coefficient matrix function. Furthermore assume that $A_{0}$ has one of the following extra properties.
(b) $A_{0}$ is a block matrix, i.e. $\left(a_{0}\right)_{0 i}(x)=\left(a_{0}\right)_{i 0}(x)=0$ for all $1 \leqslant i \leqslant n$ and all $x \in \mathbf{R}^{n}$.
(c) $A_{0}$ is a constant coefficient matrix, i.e. $A_{0}(x)=A_{0}(y)$ for all $x, y \in \mathbf{R}^{n}$.
(s) $A_{0}$ is a real symmetric matrix, i.e. $\left(a_{0}\right)_{i j}(x)=\left(a_{0}\right)_{j i}(x) \in \mathbf{R}$ for all $0 \leqslant i, j \leqslant n$ and all $x \in \mathbf{R}^{n}$.

Then there exists $\varepsilon>0$ depending only on $\left\|A_{0}\right\|_{\infty}$, the accretivity constant $\kappa_{A_{0}}$ and the dimension $n$, such that if $A \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\mathbf{C}^{n+1}\right)\right)$ is $t$-independent and satisfies $\left\|A-A_{0}\right\|_{\infty}<\varepsilon$, then Neumann and Regularity problems $(\mathrm{Neu}-A),\left(\mathrm{Neu}^{\perp}-A\right)$ and $(\operatorname{Reg}-A)$ above have a unique solution $F(t, x)$ with the required properties for every boundary function $g(x)$, being $\phi(x)$ and $\nabla_{x} \psi(x)$, respectively. Furthermore Dirichlet problem (Dir- $A$ ) above has a unique solution $U(t, x)$ with the required properties for every boundary function $u(x)$.

The solutions depend continuously on the data with the following equivalences of norms. If we define the triple bar norm $\left\|G_{t}\right\|^{2}:=\int_{0}^{\infty}\left\|G_{t}\right\|_{2}^{2} t^{-1} d t$ and the non-tangential maximal function

$$
\widetilde{N}_{*}(F)(x):=\sup _{t>0}\left(f_{|s-t|<c_{0} t} f_{|y-x|<c_{1} t}|F(s, y)|^{2} d s d y\right)^{1 / 2}
$$

where $f_{E}:=|E|^{-1} \int_{E}$ and $c_{0} \in(0,1), c_{1}>0$ are constants, then for Neumann and Regularity problems we have

$$
\|g\|_{2} \approx\|f\|_{2} \approx \sup _{t>0}\left\|F_{t}\right\|_{2} \approx\left\|t \partial_{t} F_{t}\right\|\left\|\approx \tilde{N}_{*}(F)\right\|_{2}
$$

and for Dirichlet problem we have

$$
\|u\|_{2} \approx \sup _{t>0}\left\|U_{t}\right\|_{2} \approx\left\|t \partial_{t} U_{t}\right\|\|\approx\| t \nabla_{x} U_{t}\| \| \approx\left\|\tilde{N}_{*}(U)\right\|_{2}
$$

Moreover, the solution operators $S_{A}$, being $S_{A}(g)=F$ or $S_{A}(u)=U$, respectively, depend Lipschitz continuously on $A$, i.e. there exists $C<\infty$ such that

$$
\left\|S_{A_{2}}-S_{A_{1}}\right\|_{L_{2}\left(\mathbf{R}^{n}\right) \rightarrow \mathcal{X}} \leqslant C\left\|A_{2}-A_{1}\right\|_{L_{\infty}\left(\mathbf{R}^{n}\right)}
$$

when $\left\|A_{i}-A_{0}\right\|_{\infty}<\varepsilon, i=1,2$, where $\|F\|_{\mathcal{X}}$ or $\|U\|_{\mathcal{X}}$ denotes any of the norms above.
Throughout this paper, we use the notation $X \approx Y$ and $X \lesssim Y$ to mean that there exists a constant $C>0$ so that $X / C \leqslant Y \leqslant C X$ and $X \leqslant C Y$, respectively. The value of $C$ varies from one usage to the next, but then is always fixed.

Let us now review the history of works on these boundary value problems, starting with the case of matrices of the form $\tilde{A}$ which we now describe. By standard arguments, Theorem 1.1 also shows well-posedness of the corresponding boundary value problems on the region $\Omega$ above a Lipschitz graph $\Sigma=\{(t, x) ; t=g(x)\}$, where $g: \mathbf{R}^{n} \rightarrow \mathbf{R}$ is a Lipschitz function. Indeed, if the function $U(t, x)$ satisfies $\operatorname{div}_{t, x} A(x) \nabla_{t, x} U=0$ in $\Omega$ then $\widetilde{U}(t, x):=U(t+g(x), x)$ satisfies $\operatorname{div}_{t, x} \tilde{A}(x) \nabla_{t, x} \tilde{U}=0$ in $\mathbf{R}_{+}^{n+1}$, where

$$
\tilde{A}(x):=\left[\begin{array}{cc}
1 & -\left(\nabla_{x} g(x)\right)^{t} \\
0 & I
\end{array}\right] A(x)\left[\begin{array}{cc}
1 & 0 \\
-\nabla_{x} g(x) & I
\end{array}\right] .
$$

Thus Theorem 1.1 gives conditions on $A$ for which Neumann problem (conormal derivative $e_{0} \cdot \tilde{A} \nabla_{t, x} \widetilde{U}=\left(e_{0}-\nabla_{x} g\right) \cdot A \nabla_{t, x} U$ given), Regularity problem (tangential gradient $\nabla_{x} \widetilde{U}=$ $\left(\partial_{t} U\right) \nabla_{x} g+\nabla_{x} U$ given $)$, and Dirichlet problem ( $\widetilde{U}=U$ given) are well posed. Note that $A$ is real symmetric if and only if $\tilde{A}$ is, but that $\tilde{A}$ being constant or of block form does not imply the same for $A$. For the Laplace equation $A=I$ in $\Omega$, solvability of (Neu- $\tilde{I}$ ) and (Reg- $\tilde{I}$ ) was first proved by Jerison and Kenig [16], and solvability of (Dir- $\tilde{I}$ ) was first proved by Dahlberg [12]. Later Verchota [26] showed that these boundary value problems are solvable with the layer potential integral equation method.

For general real symmetric matrices $A$, not being of the "Jacobian type" $\tilde{I}$ above, the wellposedness of (Dir- $A$ ) was first proved by Jerison and Kenig [17], and (Neu-A) and (Reg-A) by Kenig and Pipher [20]. These results make use of the Rellich estimate technique. For Neumann and Regularity problems, this integration by parts technique yields an equivalence

$$
\begin{equation*}
\left\|\frac{\partial U}{\partial v_{A}}\right\|_{L_{2}\left(\mathbf{R}^{n}\right)} \approx\left\|\nabla_{x} U\right\|_{L_{2}\left(\mathbf{R}^{n}\right)}, \tag{1.3}
\end{equation*}
$$

which is seen to be equivalent with the first estimate $\|f\| \approx\|g\|$ in the theorem above, and shows that the boundary trace $f$ splits into two parts of comparable size.

Turning to the unperturbed case where $A=A_{0}$ satisfies (b), then (1.3) is still valid, but the proof is far deeper than Rellich estimates. In fact, it is equivalent with the Kato square root estimate proved by Auscher, Hofmann, Lacey, M${ }^{\mathrm{c}}$ Intosh and Tchamitchian in [4]. (For the nondivergence form case $a_{00} \neq 1$, see [10].) For details concerning this equivalence between the Kato problem and the boundedness and invertibility of the Dirichlet-to-Neumann map $\nabla_{x} U \mapsto \frac{\partial U}{\partial \nu}$ we refer to Kenig [19, Remark 2.5.6], where also many further references in the field can be found.

We now consider what is previously known in the case when $A$ does not satisfy (b), (c) or (s). Here ( $\operatorname{Dir}-A$ ) has been showed to be well posed by Fabes, Jerison and Kenig [14] for small perturbations of (c), using the method of multilinear expansions. More recently, the boundary value problems have been studied in the $L_{p}$ setting and for real but non-symmetric matrices in the plane, i.e. $n=1$. Here Kenig, Koch, Pipher and Toro [22] have obtained solvability of the Dirichlet problem for sufficiently large $p$, and Kenig and Rule [21] have shown solvability of the Neumann and regularity problems for sufficiently small dual exponent $p^{\prime}$.

In the perturbed case $A \approx A_{0}$ when $A_{0}$ satisfies (c) or (s), the well-posedness of (Neu- $A$ ), $(\operatorname{Reg}-A)$ and $(\operatorname{Dir}-A)$ is also proved in [2] by Alfonseca, Auscher, Axelsson, Hofmann and Kim. With the further assumption of pointwise resolvent kernel bounds, perturbation of case (b) is also implicit in [2]. It is worth comparing the present methods to those of [2]. In [2], due to the presence of kernel bounds, the solvability of the boundary value problems is meant in the sense of non-tangential maximal estimates at the boundary and this follows from the use of layer
potentials. The first main result in [2] in the unperturbed case (s), is the proof via singular integral operator theory of boundedness and invertibility of layer potentials. The second main result in [2] is the stability of the simultaneous occurrence of both boundedness and invertibility, which hold in the unperturbed cases (c), (s) and (b). Solvability then follows.

Here, we setup a different resolution algorithm (forcing us to introduce some substantial material), which consists in solving the first order system (1.2) instead of (1.1), also by a boundary operator method, but acting on the gradient of solutions involving a generalised Cauchy operator $E_{A}$, the goal being to establish boundedness of $E_{A}$ and invertibility of related operators $E_{A} \pm N_{A}$. Boundedness of $E_{A}=\operatorname{sgn}\left(T_{A}\right)$ is obtained via quadratic estimates of an underlaying first order differential operator $T_{A}$, and the deep fact is here that those quadratic estimates alone are stable under perturbations. Stability of invertibility is then easy. The perturbation argument requires sophisticated harmonic analysis techniques inspired by the strategy of [2]. In particular, the latter uses extensively the technology of the solution of the Kato problem for second order operators in [4], whereas we utilise here the work of Axelsson, Keith and $\mathrm{M}^{\mathrm{c}}$ Intosh [10], which adapts and extends this technology to first order operators of Dirac type. Indeed, we note that our Dirac type operators $T_{A}$ are of the form $\Pi_{A}$ of [10] in the case of block matrices (b) of Theorem 1.1. But $T_{A}$ has a more complicated structure when $A$ is not a block matrix and we understand how to prove boundedness of $E_{A}$ at the moment only in the cases specified by Theorem 1.1. We also note that the present paper, like [10], makes no use of kernel bounds and only needs $L_{2}$ off-diagonal bounds for the operators, which always holds.

The boundary operator method for first order Dirac type operators, used here to solve second order boundary value problems, was developed in the thesis of Axelsson [6], which has been published as the four papers [5,7-9]. It covers operators on Lipschitz domains as described above and in Example 1.5. The result in [10] pursued the program initiated by Auscher, $\mathrm{M}^{\mathrm{c}}$ Intosh and Nahmod in [3], consisting of connecting the Kato problem and the functional calculus of first order differential operators of Dirac type. As said, it thus applies to the boundary value problems for operators of case (b). What is new here is the setup for full matrices encompassing the above. We prove also a sort of meta-theorem (see Theorem 1.3) which roughly says that the set of matrices for which the needed quadratic estimates on $T_{A}$ hold, is open.

We also show that non-tangential maximal estimates hold for our solutions. By uniqueness in the class of solutions of (1.1) with non-tangential maximal estimates, this implies that our solutions are the same as those in [2] for perturbations of the real symmetric and constant cases. The non-tangential maximal estimate here also yields an indirect proof of non-tangential limits of solutions of (1.2) which hold for the solutions of (1.1) in [2]. We do not know how to prove this fact directly in the framework of this article. Note also that we prove here that the nontangential maximal functions have comparable $L_{2}$-norms for different values of the parameters $c_{0}$ and $c_{1}$, and that the slightly different non-tangential maximal function used in [2] therefore has comparable norm.

Before turning to the method of proof for Theorem 1.1, we would like to stress the importance of the final result that the solution operators $g \mapsto F$ and $u \mapsto U$ depend Lipschitz continuously on $L_{\infty}$ changes of the matrix $A$ around $A_{0}$. This is an important motivation for considering complex $A$, as the authors do not know any proof of this perturbation result which does not make use of boundedness of the operators in a complex neighbourhood of $A_{0}$. We also remark that we in fact prove that $A \mapsto S_{A}$ is holomorphic, from which we deduce Lipschitz continuity as a corollary.

In this paper, we shall mainly focus on the boundary value problems (Neu-A) and (Reg-A). The reason is that, assuming the Cauchy operator $E_{A}$ is bounded, we prove in Section 2.5 that well-posedness follows as

$$
\left(\operatorname{Reg}-A^{*}\right) \quad \Longleftrightarrow \quad\left(\mathrm{Neu}^{\perp}-A\right) \quad \Longrightarrow \quad(\operatorname{Dir}-A)
$$

That (Reg- $A^{*}$ ) implies (Dir- $A$ ) has been proved by Kenig and Pipher [20, Theorem 5.4] in the case of real matrices $A$.

### 1.1. Operators and vector fields

We now explain the basic ideas of the method we use for the proof of Theorem 1.1. The appropriate Hilbert space on the boundary $\mathbf{R}^{n}$ is

$$
\hat{\mathcal{H}}^{1}:=\left\{f \in L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}^{n+1}\right) ; \operatorname{curl}_{x}\left(f_{\|}\right)=0\right\} .
$$

The condition on $f$ means that its tangential part is curl-free. Indeed, the trace $f(x)$ of a vector field $F(t, x)$ solving (1.2) belongs to $\hat{\mathcal{H}}^{1}$ due to the second equation in (1.2). The basic picture, building on ideas from [7], is that the Hilbert space splits into two different pairs of complementary subspaces as

$$
\begin{equation*}
\hat{\mathcal{H}}^{1}=E_{A}^{+} \hat{\mathcal{H}}^{1} \oplus E_{A}^{-} \hat{\mathcal{H}}^{1}=N_{A}^{+} \hat{\mathcal{H}}^{1} \oplus N_{A}^{-} \hat{\mathcal{H}}^{1} . \tag{1.4}
\end{equation*}
$$

We first discuss the splitting into the Hardy type subspaces $E_{A}^{ \pm} \hat{\mathcal{H}}^{1}$, consisting of $L_{2}$ boundary traces of vector fields $F^{ \pm}$solving (1.2) in $\mathbf{R}_{ \pm}^{n+1}$, respectively. Our main work in this paper is to establish boundedness of the projection operators $E_{A}^{ \pm}$for certain $A$. These projections can be written $E_{A}^{ \pm}=\frac{1}{2}\left(I \pm E_{A}\right)$, where $E_{A}$ for simple $A$ is a singular integral operator of Cauchy type. However, in the general case $E_{A}$ may fail to be a singular integral operator. To handle the projections $E_{A}^{ \pm}$we make use of functional calculus of closed Hilbert space operators, and show that $E_{A}^{ \pm}=\chi_{ \pm}\left(T_{A}\right)$ are the spectral projections of an underlaying bisectorial operator $T_{A}$ in $\hat{\mathcal{H}}^{1}$. The functions $\chi_{ \pm}(z)$ are the characteristic functions for the right and left complex half-planes. To find $T_{A}$, assume $F(t, x)$ satisfies (1.2) in $\mathbf{R}_{+}^{n+1}$ and solve for the vertical derivative

$$
\begin{aligned}
\partial_{t} F_{0} & =-a_{00}^{-1}\left(\sum_{i=1}^{n} a_{0 i} \partial_{i} F_{0}+\partial_{i}(A F)_{i}\right) \\
\partial_{t} F_{i} & =\partial_{i} F_{0}, \quad i=1, \ldots, n
\end{aligned}
$$

The right-hand side defines an operator $-T_{A}$ in $\hat{\mathcal{H}}^{1}$ which on $F(t, x)$, for fixed $t>0$, satisfies

$$
\partial_{t} F+T_{A} F=0
$$

Concretely, if we identify $f=f_{0} e_{0}+f_{\|}$with $\left(f_{0}, f_{\|}\right)^{t}$, where $f_{\|}$is a tangential curl-free vector field, then

$$
T_{A} f=\left[\begin{array}{cc}
A_{00}^{-1}\left(\left(A_{0 \|}, \nabla_{x}\right)+\operatorname{div}_{x} A_{\| 0}\right) & A_{00}^{-1} \operatorname{div}_{x} A_{\| \|}  \tag{1.5}\\
-\nabla_{x} & 0
\end{array}\right]\left[\begin{array}{l}
f_{0} \\
f_{\|}
\end{array}\right]
$$

where

$$
\begin{gathered}
\mathrm{D}\left(T_{A}\right)=\left\{f=\left(f_{0}, f_{\|}\right)^{t} \in \hat{\mathcal{H}}^{1} ; \nabla_{x} f_{0} \in L_{2}, \operatorname{div}_{x}(A f)_{\|} \in L_{2}\right\} \quad \text { and } \\
A=\left[\begin{array}{ll}
A_{00} & A_{0 \|} \\
A_{\| 0} & A_{\| \|}
\end{array}\right] .
\end{gathered}
$$

If $F(t, x)$ is a vector field in $\mathbf{R}_{+}^{n+1}$ satisfying (1.2), then using this operator $T_{A}$, we can reproduce $F$ provided we know the full trace $f=\left.F\right|_{\mathbf{R}^{n}}$, through a Cauchy type reproducing formula $F(t, x)=\left(e^{-t\left|T_{A}\right|} f\right)(x)$. However, in (Neu- $A$ ) and (Reg- $A$ ) only "half" of the trace $f$ is known since the boundary conditions for $f$ are $e_{0} \cdot A f=\phi$ and $e_{0} \wedge f=e_{0} \wedge \nabla \psi$, respectively.

We now turn to the second splitting in (1.4), which is used to split the boundary trace $f$ into the regularity and Neumann data. We define the $A$-tangential and normal subspaces of $\hat{\mathcal{H}}^{1}$ to be the null spaces of these two operators:

$$
\begin{aligned}
& N_{A}^{+} \hat{\mathcal{H}}^{1}:=\left\{f \in \hat{\mathcal{H}}^{1} ; e_{0} \cdot A f=0\right\}, \\
& N_{A}^{-} \hat{\mathcal{H}}^{1}:=\left\{f \in \hat{\mathcal{H}}^{1} ; e_{0} \wedge f=0\right\} .
\end{aligned}
$$

In contrast with the Hardy subspaces, it is straightforward to show that we have a topological splitting $\hat{\mathcal{H}}^{1}=N_{A}^{+} \hat{\mathcal{H}}^{1} \oplus N_{A}^{-} \hat{\mathcal{H}}^{1}$, and therefore that the corresponding pair of projections $N_{A}^{ \pm}$are bounded. We can now reformulate (Neu- $A$ ) and (Reg- $A$ ) as follows. Neumann problem (Neu- $A$ ) being well posed means that the restricted projection

$$
N_{A}^{-}: E_{A}^{+} \hat{\mathcal{H}}^{1} \rightarrow N_{A}^{-} \hat{\mathcal{H}}^{1}
$$

is an isomorphism, since $N_{A}^{-} \hat{\mathcal{H}}^{1}$ is a complement of the null space of $e_{0} \cdot A(\cdot)$. Similarly Regularity problem $(\operatorname{Reg}-A)$ being well posed means that the restricted projection

$$
N_{A}^{+}: E_{A}^{+} \hat{\mathcal{H}}^{1} \rightarrow N_{A}^{+} \hat{\mathcal{H}}^{1}
$$

is an isomorphism, since $N_{A}^{+} \hat{\mathcal{H}}^{1}$ is a complement of the null space of $e_{0} \wedge(\cdot)$. Note that what is important here is which subspace $N_{A}^{ \pm}$projects along, not what subspace they project onto.

We shall also find it convenient to use the operators $E_{A}:=E_{A}^{+}-E_{A}^{-}=\operatorname{sgn}\left(T_{A}\right)$ and $N_{A}:=$ $N_{A}^{+}-N_{A}^{-}$. These are reflection operators, i.e. $E_{A}^{2}=I$ and $N_{A}^{2}=I$, and we have $E_{A}^{ \pm}=\frac{1}{2}\left(I \pm E_{A}\right)$ and $N_{A}^{ \pm}=\frac{1}{2}\left(I \pm N_{A}\right)$.

Example 1.2. Let $n=1$ and $A=I$. Then the space $\hat{\mathcal{H}}^{1}$ is simply $L_{2}\left(\mathbf{R} ; \mathbf{C}^{2}\right)$ and the fundamental operator $T_{A}$ becomes

$$
T:=T_{I}=\left[\begin{array}{cc}
0 & \frac{d}{d x} \\
-\frac{d}{d x} & 0
\end{array}\right] \approx\left[\begin{array}{cc}
0 & i \xi \\
-i \xi & 0
\end{array}\right]
$$

if $\approx$ denotes conjugation with Fourier transform. Furthermore

$$
E:=E_{I}=\operatorname{sgn}(T)=\left[\begin{array}{cc}
0 & i H \\
-i H & 0
\end{array}\right], \quad N:=N_{I}=\left[\begin{array}{cc}
-1 & 0 \\
0 & 1
\end{array}\right],
$$

where

$$
H f(x):=\frac{i}{\pi} \text { p.v. } \int \frac{f(y)}{x-y} d y .
$$

Note that the operator $E$ is contained in the Borel functional calculus of the self-adjoint operator $T$, and it follows that $\|E\|=1$. On the other hand the operator $N$ is outside the Borel functional calculus of $T$. Indeed, the operators $b(T)$ in the Borel functional calculus of $T$ all commute, but we have the anticommutation relation

$$
\begin{equation*}
E N+N E=0 \tag{1.6}
\end{equation*}
$$

A consequence of this equation is that the boundary value problems (Neu- $I$ ) and (Reg- $I$ ) are well posed, or equivalently that $N^{-}: E^{+} L_{2} \rightarrow N^{-} L_{2}$ and $N^{+}: E^{+} L_{2} \rightarrow N^{+} L_{2}$ are isomorphisms. Consider for example (Reg-I) and assume we want to solve $N^{+} f=g$ for $f \in E^{+} L_{2}$. Applying $4 E^{+}$to the equation gives

$$
4 E^{+} g=(I+E)(I+N) f=f+E f+N f+E N f=f+f+N f-N f=2 f
$$

so $f=2 E^{+} g$ and it follows that $N^{+}: E^{+} L_{2} \rightarrow N^{+} L_{2}$ is an isomorphism. Having solved for $f \in E^{+} L_{2}$, we can find the solution $F(t, x)=C_{t}^{+} f(x)$ in $\mathbf{R}_{+}^{2}$ by using the Cauchy extension $C_{t}^{+}:=e^{-t|T|} E^{+}$for $t>0$. As a convolution operator, the Fourier multiplier $C_{t}^{+}$has the expression

$$
\begin{aligned}
C_{t}^{+}\left(u_{0} e_{0}+u_{1} e_{1}\right)= & \left(\frac{1}{2 \pi} \int_{\mathbf{R}} \frac{t u_{0}(y)-(x-y) u_{1}(y)}{t^{2}+(x-y)^{2}} d y\right) e_{0} \\
& +\left(\frac{1}{2 \pi} \int_{\mathbf{R}} \frac{(x-y) u_{0}(y)+t u_{1}(y)}{t^{2}+(x-y)^{2}} d y\right) e_{1},
\end{aligned}
$$

and in particular

$$
F(t, x)=C_{t}^{+} f(x)=2 C_{t}^{+}\left(g_{1} e_{1}\right)=\frac{1}{\pi} \int_{\mathbf{R}} \frac{\left(-(x-y) e_{0}+t e_{1}\right) g_{1}(y)}{t^{2}+(x-y)^{2}} d y .
$$

For a more general $A$, even if $A$ is real symmetric, the operator $T_{A}$ is not self-adjoint and proving boundedness of $E_{A}^{ \pm}$is a highly non-trivial problem. Also Eq. (1.6) fails for general $A$, and therefore such explicit formulae for the solution $F(t, x)$ as in Example 1.2 are not available. However, to show well-posedness of (Neu- $A$ ) and (Reg- $A$ ) it suffices to show that $I \pm \frac{1}{2}\left(E_{A} N_{A}+\right.$ $\left.N_{A} E_{A}\right)=\frac{1}{2}\left(E_{A} \pm N_{A}\right)^{2}$ are invertible, as explained in [7]. To summarise, in order to solve (Neu- $A$ ) and (Reg- $A$ ) we need that:
(i) the Hardy projections $E_{A}^{ \pm}$are bounded, so that we have a topological splitting $\hat{\mathcal{H}}^{1}=E_{A}^{+} \hat{\mathcal{H}}^{1} \oplus$ $E_{A}^{-} \hat{\mathcal{H}}^{1}$, and
(ii) the restricted projections $N_{A}^{+}: E_{A}^{+} \hat{\mathcal{H}}^{1} \rightarrow N_{A}^{+} \hat{\mathcal{H}}^{1}$ and $N_{A}^{-}: E_{A}^{+} \hat{\mathcal{H}}^{1} \rightarrow N_{A}^{-} \hat{\mathcal{H}}^{1}$ are isomorphisms.

To prove Theorem 1.1, we shall prove the following.
(i') That $T_{A}$ satisfies quadratic estimates for all $A$ such that $\left\|A-A_{0}\right\|_{\infty}<\varepsilon$. From this it will follow that $E_{A}^{ \pm}$are bounded for all such $A$ and that $A \mapsto E_{A}^{ \pm}$is continuous (in fact holomorphic).
(ii') That (ii) holds for the unperturbed $A_{0}$. From (i') and since clearly $N_{A}^{ \pm}$are bounded and depend continuously on $A$, it then follows from a continuity argument that (ii) holds for all $A$ in a neighbourhood of $A_{0}$.

We emphasise that boundedness of the Hardy projections $E_{A}^{ \pm}$alone does not show that (Neu- $A$ ) and ( $\operatorname{Reg}-A$ ) are well posed. In our framework, the boundary value problems being well posed means that the Hardy space $E_{A}^{+} \hat{\mathcal{H}}^{1}$, which exists as a closed subspace when the Hardy projections are bounded, is transversal to the $A$-tangential and normal subspaces $N_{A}^{ \pm} \hat{\mathcal{H}}^{1}$. A concrete example showing that $(\operatorname{Neu}-A)$ and ( $\operatorname{Reg}-A$ ) may fail to be well posed, even though $E_{A}$ is bounded, is furnished by the matrix

$$
A(x)=\left[\begin{array}{cc}
1 & k \operatorname{sgn}(x) \\
-k \operatorname{sgn}(x) & 1
\end{array}\right],
$$

with parameter $k \in \mathbf{R}$. The corresponding elliptic equation (1.1) in $\mathbf{R}_{+}^{2}$ was studied by Kenig, Koch, Pipher and Toro in [22], where they showed that (Dir-A) fails to be well posed for certain values of $k$. Moreover, that (Neu-A) and (Reg-A) also fails for some $k$, is shown by Kenig and Rule [21]. On the other hand, $\left\|E_{A}\right\|=1$ for all $k \in \mathbf{R}$ since according to (1.5)

$$
T_{A}=\left[\begin{array}{cc}
k\left(\operatorname{sgn}(x) \frac{d}{d x}-\frac{d}{d x} \operatorname{sgn}(x)\right) & \frac{d}{d x} \\
-\frac{d}{d x} & 0
\end{array}\right]
$$

is self-adjoint, and therefore $E_{A}^{ \pm}$are orthogonal projections.

### 1.2. Embedding in a Dirac equation

Unfortunately there is a technical problem in applying harmonic analysis to the operator $T_{A}$ in order to prove ( $\mathrm{i}^{\prime}$ ): the space $\hat{\mathcal{H}}^{1}$ is defined through the non-local condition $\operatorname{curl}_{x}\left(f_{\|}\right)=0$. This prevents us from using multiplication operators, for example when localising with a cutoff $f \mapsto \eta f$, as these does not preserve $\hat{\mathcal{H}}^{1}$. To avoid this problem we embed $\hat{\mathcal{H}}^{1} \subset \mathcal{H}:=$ $L_{2}\left(\mathbf{R}^{n} ; \bigwedge_{\mathbf{C}} \mathbf{R}^{n+1}\right)$, where

$$
\bigwedge_{\mathbf{C}} \mathbf{R}^{n+1}=\bigwedge^{0} \oplus \bigwedge^{1} \oplus \bigwedge^{2} \oplus \cdots \oplus \bigwedge^{n+1}
$$

is the full complex exterior algebra of $\mathbf{R}^{n+1}$, which in particular contains the vectors $\bigwedge^{1}=\mathbf{C}^{n+1}$ and the scalars $\bigwedge^{0}=\mathbf{C}$ along with all $k$-vectors $\bigwedge^{k}$. (We identify $k$-vectors with the dual $k$-forms in Euclidean space.) In this way we obtain closure, i.e. all operators, including multiplication operators preserve $\mathcal{H}$. Furthermore we embed Eq. (1.2) in a Dirac type equation

$$
\begin{equation*}
\left(d_{t, x}+\widetilde{B}(x)^{-1} d_{t, x}^{*} B(x)\right) F(t, x)=0 \tag{1.7}
\end{equation*}
$$

where $B, \widetilde{B} \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\bigwedge_{\mathbf{C}} \mathbf{R}^{n+1}\right)\right)$ are accretive and $\left.B\right|_{\bigwedge^{1}}=A$. Here $d_{t, x}$ and $d_{t, x}^{*}$ are the exterior and interior derivative operators, as defined in Section 2. In particular, if $F(t, x): \mathbf{R}_{+}^{n+1} \rightarrow$ $\bigwedge^{1}$ is a vector field satisfying (1.2) then

$$
\begin{gathered}
d_{t, x}^{*}(B F)=-\operatorname{div}_{t, x}(A F)=0 \\
d_{t, x} F=\operatorname{curl}_{t, x} F=0
\end{gathered}
$$

Thus (1.7) follows from (1.2) for any choice of the auxiliary function $\widetilde{B}$. In the same way as in Section 1.1 we shall solve for the vertical derivative in (1.7) and obtain an operator $T_{B}$ acting in $\mathcal{H}$, such that $T_{A}=\left.T_{B}\right|_{\hat{\mathcal{H}}^{1}}$. For applications to the Neumann and regularity problems it suffices to consider $B=I \oplus A \oplus I \oplus \cdots \oplus I$. The stability result for the operator $T_{A}$ we prove in this paper is the following.

Theorem 1.3. Let $A_{0}$ be a t-independent, complex, accretive coefficient matrix function such that $T_{B_{0}}$ has quadratic estimates in $\mathcal{H}$, where $B_{0}=I \oplus A_{0} \oplus I \oplus \cdots \oplus I$. Then there exists $\varepsilon_{0}>0$ depending only on $\left\|A_{0}\right\|_{\infty}, \kappa_{A_{0}}$ and $n$, such that if $A \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\mathbf{C}^{n+1}\right)\right)$ is $t$-independent and satisfies $\left\|A-A_{0}\right\|_{\infty}<\varepsilon_{0}$, then $T_{B}$ has quadratic estimates in $\mathcal{H}$, where $B=I \oplus A \oplus I \oplus \cdots \oplus I$. In particular $T_{A}$ has quadratic estimates in $\hat{\mathcal{H}}^{1}$.

Thus $\hat{\mathcal{H}}^{1}$ splits into Hardy subspaces, the spectral subspaces of $T_{A}$, i.e. each $f \in \hat{\mathcal{H}}^{1}$ can be uniquely written $f=f^{+}+f^{-}$, where $f^{ \pm}=\left.F^{ \pm}\right|_{\mathbf{R}^{n}}$ and $F^{ \pm}(t, x)=e^{\mp t\left|T_{A}\right|} E_{A}^{ \pm} f(x)$ satisfy (1.2) in $\mathbf{R}_{ \pm}^{n+1}$. Moreover, we have equivalence of the norms

$$
\|f\|_{2} \approx\left\|f^{+}\right\|_{2}+\left\|f^{-}\right\|_{2}
$$

and $\left\|f^{ \pm}\right\|_{2} \approx \sup _{t>0}\left\|F_{ \pm t}^{ \pm}\right\|_{2} \approx\left\|t \partial_{t} F_{ \pm t}^{ \pm}\right\|\| \| \tilde{N}_{*}\left(F^{ \pm}\right) \|_{2}$.
With a more general choice

$$
B=B^{0} \oplus B^{1} \oplus B^{2} \oplus \cdots \oplus B^{n+1}
$$

where $B^{k} \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\bigwedge^{k}\right)\right)$, we also obtain new perturbation results concerning boundary value problems and more generally transmission problems for $k$-vector fields. For $k$-vector fields we consider the function space

$$
\hat{\mathcal{H}}_{B}^{k}:=\left\{f \in L_{2}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\bigwedge^{k}\right)\right) ; d_{x}\left(f_{\|}\right)=0=d_{x}^{*}\left(\left(B^{k} f\right)_{\perp}\right)\right\} \subset \mathcal{H}
$$

where $f_{\|}$and $f_{\perp}$ denote the tangential and normal parts of $f$. This is the appropriate function space since traces of $k$-vector fields $F(t, x)$ in $\mathbf{R}_{ \pm}^{n+1}$ satisfying

$$
\left\{\begin{array}{l}
d_{t, x}^{*}\left(B^{k}(x) F(t, x)\right)=0  \tag{1.8}\\
d_{t, x} F(t, x)=0
\end{array}\right.
$$

belong to $\hat{\mathcal{H}}_{B}^{k}$. Note that for 1 -vectors, i.e. vectors, the condition $d_{x}^{*}\left(\left(B^{1} f\right)_{\perp}\right)=0$ is void and $\hat{\mathcal{H}}_{B}^{1}=\hat{\mathcal{H}}^{1}$. For $k$-vector fields, we consider the following.

Transmission problem $\left(\operatorname{Tr}-\boldsymbol{B}^{k} \boldsymbol{\alpha}^{ \pm}\right)$. Let $B^{k}=B^{k}(x) \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\bigwedge^{k}\right)\right)$ be accretive and let $\alpha^{ \pm} \in \mathbf{C}$ be given jump parameters. Given a $k$-vector field $g \in \hat{\mathcal{H}}_{B}^{k}$, find $k$-vector fields $F_{t}^{+}(x)=$ $F^{+}(t, x)$ in $\mathbf{R}_{+}^{n+1}$ and $F_{t}^{-}(x)=F^{-}(t, x)$ in $\mathbf{R}_{-}^{n+1}$ such that $F_{t}^{ \pm} \in C^{1}\left(\mathbf{R}_{+} ; L_{2}\left(\mathbf{R}^{n} ; \bigwedge^{k}\right)\right)$ and $F^{ \pm}$ satisfies (1.8) for $\pm t>0$, and furthermore $\lim _{t \rightarrow \pm \infty} F_{t}^{ \pm}=0$ and $\lim _{t \rightarrow 0^{ \pm}} F_{t}^{ \pm}=f^{ \pm}$in $L_{2}$ norm, where the traces $f^{ \pm}$satisfy the jump conditions

$$
\left\{\begin{array}{l}
e_{0} \wedge\left(\alpha^{-} f^{+}-\alpha^{+} f^{-}\right)=e_{0} \wedge g \\
\left.\left.e_{0}\right\lrcorner\left(B^{k}\left(\alpha^{+} f^{+}-\alpha^{-} f^{-}\right)\right)=e_{0}\right\lrcorner\left(B^{k} g\right) .
\end{array}\right.
$$

The second main result of this paper is the following $L_{\infty}$ perturbation result of Transmission problem $\left(\operatorname{Tr}-B^{k} \alpha^{ \pm}\right)$.

Theorem 1.4. Let $B_{0}^{k}=B_{0}^{k}(x) \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\bigwedge^{k}\right)\right)$ be t-independent, accretive and possibly complex, and assume that $B_{0}^{k}$ is a block matrix, i.e.

$$
\left(B_{0}^{k}\right)_{s t}=0, \quad \text { whenever } 0 \in(s \backslash t) \cup(t \backslash s)
$$

and $s, t \subset\{0,1, \ldots, n\}$ has lengths $|s|=|t|=k$. Then there exist $\varepsilon>0$ and $C<\infty$ depending only on $\left\|B_{0}^{k}\right\|_{\infty}$, the accretivity constant $\kappa_{B_{0}^{k}}$ and dimension $n$, such that if $B^{k}=B^{k}(x) \in$ $L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\bigwedge^{k}\right)\right)$ is $t$-independent and satisfies

$$
\begin{equation*}
\left\|B^{k}-B_{0}^{k}\right\|_{L_{\infty}\left(\mathbf{R}^{n}\right)}<\min \left(\varepsilon, C\left|\left(\alpha^{+} / \alpha^{-}\right)^{2}+1\right|\right) \tag{1.9}
\end{equation*}
$$

then Transmission problem $\left(\operatorname{Tr}-B^{k} \alpha^{ \pm}\right)$above is well posed in the sense that for every boundary function $g \in \hat{\mathcal{H}}_{B}^{k}$, there exist unique $k$-vector fields $F^{ \pm}(t, x)$ with properties as in $\left(\operatorname{Tr}-B^{k} \alpha^{ \pm}\right)$. The solution $F_{t}^{ \pm}$depends continuously on $g$ with equivalences of norms $\|g\|_{2} \approx\left\|f^{+}+f^{-}\right\|_{2} \approx$ $\left\|f^{+}\right\|_{2}+\left\|f^{-}\right\|_{2}$ and

$$
\left\|f^{ \pm}\right\|_{2} \approx \sup _{t>0}\left\|F_{ \pm t}^{ \pm}\right\|_{2} \approx\left\|t \partial_{t} F_{ \pm t}^{ \pm}\right\| \|
$$

This perturbation theorem for transmission problems for $k$-vectors has two important corollaries. On one hand it specialises when $k=1$ to a generalisation of Theorem 1.1(b), giving perturbation results for transmission problems across $\mathbf{R}^{n}$ for the divergence form equation (1.1). The details of this Neumann-regularity transmission problem is stated as ( $\operatorname{Tr}-A \alpha^{ \pm}$) in Section 4.3.

On the other hand it specialises when either $\alpha^{+}=0$ or $\alpha^{-}=0$ to a generalisation of Theorem 1.1(b), giving perturbation results for boundary-value problems for $k$-vectors. Our result for these boundary-value problems (Nor- $B^{k}$ ) and (Tan- $B^{k}$ ) is given as Corollary 4.17 in Section 4.3.

Example 1.5. In the case $\widetilde{B}=B$, operators of the form $d_{t, x}+B^{-1} d_{t, x}^{*} B$ appear naturally when pulling back the unperturbed Hodge-Dirac operator $d+d^{*}$ with a change of variables. As above, consider the region $\Omega$ above a Lipschitz graph $\Sigma=\{(t, x) ; t=g(x)\}$. We define the pullback of the field $F: \Omega \rightarrow \bigwedge \mathbf{R}^{n+1}$ to be the field

$$
\left(\rho^{*} F\right)(t, x):=\underline{\rho}^{\mathrm{T}}(x) F(\rho(t, x))
$$

in $\mathbf{R}_{+}^{n+1}$, where $\rho(t, x)=(t+g(x), x)$ is the parametrisation of $\Omega$, having differential

$$
\left.\underline{\rho}(x)\right|_{\Lambda^{1}}=\left[\begin{array}{cc}
1 & \nabla g(x) \\
0 & I
\end{array}\right]
$$

acting on vector fields and extended naturally to $\bigwedge \mathbf{R}^{n+1}$, and $\underline{\rho}^{\mathrm{T}}(x)$ denotes the transposed matrix. From the well-known fact that $d_{t, x}$ commutes with $\rho^{*}$, we get the intertwining relation

$$
\begin{equation*}
\left(d_{t, x}+G d_{t, x}^{*} G^{-1}\right)\left(\rho^{*} F\right)=\rho^{*}\left(\left(d_{t, x}+d_{t, x}^{*}\right) F\right), \tag{1.10}
\end{equation*}
$$

where $G(x)=\left(g_{i j}(x)\right)=\underline{\rho}^{\mathrm{T}}(x) \underline{\rho}(x)$ is the metric for the parametrisation, being real symmetric. Solving for the vertical derivative in the equation $\left(d_{t, x}+d_{t, x}^{*}\right) F=0$ in $\Omega$, gives us an operator $\mathbf{D}_{\Sigma}$ in $L_{2}\left(\Sigma ; \bigwedge \mathbf{R}^{n+1}\right)$, which is similar to the operator $\left(e_{0}-\nabla g(x)\right)^{-1}\left(d_{x}+d_{x}^{*}\right)$ in $\mathcal{H}$. From (1.10) it follows that

$$
T_{G^{-1}}\left(\underline{\rho}(x)^{\mathrm{T}} f(x)\right)=\underline{\rho}(x)^{\mathrm{T}}\left(\mathbf{D}_{\Sigma} f\right)(x) .
$$

It is known that the operator $\mathbf{D}_{\Sigma}$ satisfies quadratic estimates, and therefore so does $T_{G^{-1}}$. For references and further discussion of $\mathbf{D}_{\Sigma}$, see [10, Consequence 3.6]. From this we get the bounded Clifford-Cauchy singular integral operator

$$
\begin{aligned}
& E_{G^{-1}}\left(\underline{\rho}(x)^{\mathrm{T}} f(x)\right) \\
& \quad=\underline{\rho}(x)^{\mathrm{T}} \frac{2}{\sigma_{n}} \text { p.v. } \int_{\mathbf{R}^{n}} \frac{\left(g(x) e_{0}+x\right)-\left(g(y) e_{0}+y\right)}{\left(|y-x|^{2}+(g(y)-g(x))^{2}\right)^{(n+1) / 2}}\left(e_{0}-\nabla g(y)\right) f(y) d y,
\end{aligned}
$$

where $\sigma_{n}$ is the area of the unit $n$-sphere in $\mathbf{R}^{n+1}$ and $E_{G^{-1}}=\operatorname{sgn}\left(T_{G^{-1}}\right)$. For this reason, we shall refer to $E_{B}=\operatorname{sgn}\left(T_{B}\right)$ as generalised Cauchy integral operators and $E_{B}^{ \pm}=\chi_{ \pm}\left(T_{B}\right)$ as generalised Hardy projection operators, also when $B \neq G^{-1}$ and $\widetilde{B} \neq B$.

### 1.3. Outline of the paper

In Section 2 we explain how we use the exterior algebra $\bigwedge_{\mathbf{C}} \mathbf{R}^{n+1}$ and the exterior and interior derivative operators $d$ and $d^{*}$. In Section 2.1 we introduce the Dirac type operator $T_{B}$ which extends $T_{A}$ to the full exterior algebra as well as projection operators $N_{B}^{ \pm}$which extend the $A$ tangential and normal projections $N_{A}^{ \pm}$from above. Section 2.2 is concerned with the spectral properties of $T_{B}$, where we prove that $T_{B}$ is a bisectorial operator and that the resolvents ( $\lambda I-$ $\left.T_{B}\right)^{-1}$ has $L_{2}$ off-diagonal estimates. Section 2.3 surveys the theory of functional calculus of bisectorial operators like $T_{B}$. In Section 2.5, Lemmas 2.49 and 2.55 characterise the classes of solutions $F_{t}$ and $U_{t}$, respectively, to the boundary value problems, and are used in particular to prove uniqueness.

In Section 3 we prove ( $\mathrm{i}^{\prime}$ ) quadratic estimates and (ii') invertibility in the unperturbed case $B=B_{0}$. The most involved case is when $B_{0}$ is real symmetric. In order to prove the quadratic estimates we use the results from Section 2.4 that $T_{B_{0}}$ leaves certain subspaces $\hat{\mathcal{H}}_{B_{0}}^{k}$ and $\check{\mathcal{H}}_{B_{0}}$ invariant and that therefore it suffices to establish quadratic estimates in each subspace separately.

Section 4 contains the core harmonic analysis results of the paper, where we establish quadratic estimates for $T_{B}$ when $B \approx B_{0}$ through a perturbation argument based on Eqs. (4.3)(4.8). Section 4.1 treats the case when $B_{0}$ is a block matrix and makes use of techniques from the solution of the Kato square root problem [4]. In Lemma 4.10 we construct a new set of test functions $f_{Q}^{w}$ which also can be used in the Carleson measure estimate in [10] to simplify the proof there. Section 4.2 treats matrices of the form $B=I \oplus A \oplus I \oplus \cdots \oplus I$. The key technique is Lemma 4.14 where we compare $B_{0}$ with a corresponding block matrix $\hat{B}_{0}$ for which the results from Section 4.1 applies. The reason why this approximation $B_{0} \approx \hat{B}_{0}$ works is that the normal vector component $F^{1,0}$ has additional regularity by Lemma 2.12 when $F$ satisfy the Dirac type equation (1.7).

The paper ends with Section 4.3, where we bring the results together and prove Theorems 1.4, 1.1 and 1.3.

## 2. Operator theory and algebra

In this section we develop the operator theoretic framework we use to prove the perturbation theorems stated in the introduction. In particular we introduce our basic operator $T_{B}$, along with perturbations of the normal and tangential projection operators $N^{-}$and $N^{+}$. These all act in the Hilbert space $\mathcal{H}:=L_{2}\left(\mathbf{R}^{n} ; \bigwedge_{\mathbf{C}} \mathbf{R}^{n+1}\right)$ on the boundary $\mathbf{R}^{n}=\partial \mathbf{R}_{+}^{n+1}=\partial \mathbf{R}_{-}^{n+1}$. In $\mathbf{R}^{n+1}$ we write the standard ON basis as $\left\{e_{0}, e_{1}, \ldots, e_{n}\right\}$, where $e_{0}$ denote the vertical direction and $e_{1}, \ldots, e_{n}$ span the horizontal hyperplane $\mathbf{R}^{n}$. We write the corresponding coordinates as $x_{0}, x_{1}, \ldots, x_{n}$ and we also use the notation $t=x_{0}$. The corresponding partial derivatives we write as $\partial_{i}=\frac{\partial}{\partial x_{i}}$ and $\partial_{t}=\partial_{0}=\frac{\partial}{\partial t}$. Our functions $f \in \mathcal{H}$ take values in the full complex exterior algebra over $\mathbf{R}^{n+1}$

$$
\Lambda=\bigwedge_{\mathbf{C}} \mathbf{R}^{n+1}=\bigwedge^{0} \oplus \bigwedge^{1} \oplus \cdots \oplus \bigwedge^{n+1}
$$

This is a $2^{n+1}$-dimensional linear space with $n+2$ pairwise orthogonal subspaces $\bigwedge^{k}$ of dimensions $\binom{n+1}{k}$. With the notation $e_{s}:=e_{s_{1}} \wedge \ldots \wedge e_{s_{k}}$ if $s=\left\{s_{1}, \ldots, s_{k}\right\} \subset\{0,1, \ldots, n\}$ and $s_{1}<s_{2}<\cdots<s_{k}$, the space $\bigwedge^{k}$ of homogeneous $k$-vectors is the linear span of $\left\{e_{s} ;|s|=k\right\}$. In particular, identifying $e_{\emptyset}$ with 1 and the singleton set $\{j\}$ with $j$, we have $\bigwedge^{0}=\mathbf{C}$ and $\bigwedge^{1}=\mathbf{C}^{n+1}$. A general element in $\Lambda$ is called a multivector and is a direct sum of $k$-vectors of different degrees $k$.

Definition 2.1. Introduce the sesqui-linear scalar product

$$
(f, g)=\left(\sum_{s} f_{s} e_{s}, \sum_{t} g_{t} e_{t}\right)=\sum_{s} f_{s} \bar{g}_{s}
$$

on $\bigwedge$ and the bilinear scalar product $f \cdot g=\sum_{s} f_{s} g_{s}$. Define the counting function $\sigma(s, t):=$ $\#\left\{\left(s_{i}, t_{j}\right) ; s_{i}>t_{j}\right\}$, where $s=\left\{s_{i}\right\}, t=\left\{t_{j}\right\} \subset\{0,1, \ldots, n\}$.
(i) The exterior product $f \wedge g$ is the complex bilinear product for which

$$
e_{s} \wedge e_{t}=(-1)^{\sigma(s, t)} e_{s \cup t} \text { if } s \cap t=\emptyset \quad \text { and } \quad \text { zero otherwise. }
$$

(ii) The (left) interior product $f\lrcorner g$ is the complex bilinear product for which $\left.\left(e_{s}\right\lrcorner e_{t}, e_{u}\right)=$ $\left(e_{t}, e_{s} \wedge e_{u}\right)$ for all $s, t, u \subset\{0,1, \ldots, n\}$. Explicitly we have

$$
\left.e_{s}\right\lrcorner e_{t}=(-1)^{\sigma(s, t \backslash s)} e_{t \backslash s} \text { if } s \subset t \quad \text { and } \quad \text { zero otherwise. }
$$

Example 2.2. The most common situation is when forming a product between a vector $a=$ $\sum_{j=0}^{n} a_{j} e_{j} \in \bigwedge^{1}$ and a $k$-vector

$$
f=\sum_{0 \leqslant s_{1}<\cdots<s_{k} \leqslant n} f_{s_{1}, \ldots, s_{k}} e_{s_{1}} \wedge \ldots \wedge e_{s_{k}} .
$$

In this case the interior product $a\lrcorner f$ is a $(k-1)$-vector, whereas the exterior product $a \wedge f$ yields a $(k+1)$-vector. Note also that we embed all different spaces $\bigwedge^{k}$ of homogeneous $k$-vectors as pairwise orthogonal subspaces in the $2^{n+1}$-dimensional linear space $\bigwedge$. Thus we may add the two products to obtain (the Clifford product)

$$
a\lrcorner f+a \wedge f \in \bigwedge^{k-1} \oplus \bigwedge^{k+1} \subset \bigwedge
$$

In the special case when $f=b$ is also a vector, i.e. $k=1$, we have

$$
\begin{gathered}
a\lrcorner b=\sum_{j=0}^{n} a_{j} b_{j} \in \bigwedge^{0}, \\
a \wedge b=\sum_{0 \leqslant i<j \leqslant n}\left(a_{i} b_{j}-a_{j} b_{i}\right) e_{i} \wedge e_{j} \in \bigwedge^{2},
\end{gathered}
$$

where we write $b_{\{j\}}=b_{j}$. We see that $\left.a\right\lrcorner b$ coincide with the bilinear scalar product $a \cdot b$. Furthermore, in three dimensions $n+1=3$, the exterior product $a \wedge b \in \bigwedge^{2}$ can be identified with the vector product $a \times b \in \bigwedge^{1}$ by using the Hodge star identifications $e_{\{1,2\}} \approx e_{0},-e_{\{0,2\}} \approx e_{1}$ and $e_{\{0,1\}} \approx e_{2}$.

The following anticommutativity, associativity and derivation properties of these products summarise the fundamental algebra we shall need in this paper.

Lemma 2.3. If $a, b \in \bigwedge^{1}$ are vectors and $f, g, h \in \bigwedge$, then

$$
\begin{gathered}
a \wedge b=-b \wedge a, \quad a \wedge a=0, \\
f \wedge(g \wedge h)=(f \wedge g) \wedge h, \quad f\lrcorner(g\lrcorner h)=(g \wedge f)\lrcorner h, \\
a\lrcorner(b \wedge f)=(a \cdot b) f-b \wedge(a\lrcorner f) .
\end{gathered}
$$

We shall also frequently use that if $a \in \bigwedge^{1}$ is a real vector, then $\left.(a\lrcorner f, g\right)=(f, a \wedge g)$.
Proof. That $a \wedge b=-b \wedge a$ and $a \wedge a=0$ is readily seen from Example 2.2. These and the associativity $f \wedge(g \wedge h)=(f \wedge g) \wedge h$ are well-known properties of the exterior product. To see how $f\lrcorner(g\lrcorner h)=(g \wedge f)\lrcorner h$ follows, note first that by linearity it suffices to consider the case
when $f, g$ and $h$ all are real, i.e. have real coefficients in the standard basis $\left\{e_{s}\right\}$. Under this assumption we pair with an arbitrary $w \in \bigwedge$ and use that left interior and exterior multiplication are adjoint operations. We get

$$
\begin{aligned}
(f\lrcorner(g\lrcorner h), w) & =(g\lrcorner h, f \wedge w)=(h, g \wedge(f \wedge w)) \\
& =(h,(g \wedge f) \wedge w)=((g \wedge f)\lrcorner h, w)
\end{aligned}
$$

For the derivation identity, by linearity it suffices to prove

$$
\left.\left.e_{i}\right\lrcorner\left(e_{j} \wedge e_{s}\right)+e_{j} \wedge\left(e_{i}\right\lrcorner e_{s}\right)= \begin{cases}e_{s}, & i=j, \\ 0, & i \neq j,\end{cases}
$$

for all $s \subset\{0,1, \ldots, n\}$. This is straightforward to verify from Definition 2.1.

Definition 2.4. For a multivector $f \in \Lambda$, we write $\left.\mu f:=e_{0} \wedge f, \mu^{*} f:=e_{0}\right\lrcorner f$ and $m f:=$ $\left.e_{0} \wedge f+e_{0}\right\lrcorner f=\left(\mu+\mu^{*}\right) f$. We call

$$
\begin{aligned}
f_{\perp} & =N^{-} f:=\mu \mu^{*} f=m \mu^{*} f=\mu m f \\
f_{\|} & =N^{+} f:=\mu^{*} \mu f=m \mu f=\mu^{*} m f
\end{aligned}
$$

the normal and tangential parts of $f$, respectively.
Concretely, if $f=\sum_{\left\{s_{1}, \ldots, s_{k}\right\}} f_{s_{1}, \ldots, s_{k}} e_{s_{1}} \wedge \ldots \wedge e_{s_{k}}$, then its normal part is

$$
f_{\perp}=\sum_{\left\{s_{1}, \ldots, s_{k}\right\} \ni 0} f_{s_{1}, \ldots, s_{k}} e_{s_{1}} \wedge \ldots \wedge e_{s_{k}}
$$

and its tangential part is

$$
f_{\|}=\sum_{\left\{s_{1}, \ldots, s_{k}\right\} \ngtr 0} f_{s_{1}, \ldots, s_{k}} e_{s_{1}} \wedge \ldots \wedge e_{s_{k}} .
$$

In particular $e_{s}$ is normal if $0 \in s$ and tangential if $0 \notin s$. Note that both the subspace of normal multivectors $N^{-} \bigwedge$ and the subspace of tangential multivectors $N^{+} \bigwedge$ have dimension $2^{n}$, although the subspaces $N^{ \pm} \bigwedge^{k}$ have different dimensions in general.

Throughout this paper, upper case letters denote $\bigwedge$-valued functions $F(t, x)=F_{t}(x)$ in the domain $\mathbf{R}_{+}^{n+1}$ or $\mathbf{R}_{-}^{n+1}$ whereas lower case letters will denote $\Lambda$-valued functions $f$ on the boundary $\mathbf{R}^{n}$. We use the sesqui-linear scalar product $(f, g):=\int_{\mathbf{R}^{n}}(f(x), g(x)) d x$ on the Hilbert space $\mathcal{H}$.

Definition 2.5. Using the nabla symbols $\nabla_{x}=\nabla=\sum_{j=1}^{n} e_{j} \partial_{j}$ in $\mathbf{R}^{n}$ and $\nabla_{t, x}=\sum_{j=0}^{n} e_{j} \partial_{j}$ in $\mathbf{R}^{n+1}$, we define the operators of exterior and interior derivation as

$$
\begin{aligned}
d_{x} f=d f=\nabla \wedge f:=\sum_{j=1}^{n} e_{j} \wedge \partial_{j} f, & \left.\left.d_{x}^{*} f=d^{*} f=-\nabla\right\lrcorner f:=-\sum_{j=1}^{n} e_{j}\right\lrcorner \partial_{j} f \\
d_{t, x} F=\nabla_{t, x} \wedge F:=\sum_{j=0}^{n} e_{j} \wedge \partial_{j} F, & \left.\left.d_{t, x}^{*} F=-\nabla_{t, x}\right\lrcorner F:=-\sum_{j=0}^{n} e_{j}\right\lrcorner \partial_{j} F .
\end{aligned}
$$

We shall also find it convenient to use the operators $\underline{d}:=i m d$ and $\underline{d}^{*}=-i d^{*} m=i m d^{*}$.
Remark 2.6. Clearly, in the splitting $\Lambda=\bigwedge^{0} \oplus \bigwedge^{1} \oplus \cdots \oplus \bigwedge^{n+1}$, the operators $d$ and $d^{*}$ map $d: L_{2}\left(\mathbf{R}^{n} ; \bigwedge^{k}\right) \rightarrow L_{2}\left(\mathbf{R}^{n} ; \bigwedge^{k+1}\right)$ and $d^{*}: L_{2}\left(\mathbf{R}^{n} ; \bigwedge^{k}\right) \rightarrow L_{2}\left(\mathbf{R}^{n} ; \bigwedge^{k-1}\right)$ as unbounded operators. Moreover, if we further decompose the space of homogeneous $k$-vectors into its normal and tangential subspaces as

$$
\Lambda=\bigwedge_{\|}^{0} \oplus\left(\bigwedge_{\perp}^{1} \oplus \bigwedge_{\|}^{1}\right) \oplus\left(\bigwedge_{\perp}^{2} \oplus \bigwedge_{\|}^{2}\right) \oplus \cdots \oplus\left(\bigwedge_{\perp}^{n} \oplus \bigwedge_{\|}^{n}\right) \oplus \bigwedge_{\perp}^{n+1}
$$

where $\bigwedge_{\perp}^{k}:=N^{-} \bigwedge^{k}$ and $\bigwedge_{\|}^{k}:=N^{+} \bigwedge^{k}$, then we see that the operators $\underline{d}$ and $\underline{d}^{*}$ map

$$
\begin{aligned}
\underline{d}: L_{2}\left(\mathbf{R}^{n} ; \bigwedge_{\perp}^{k}\right) \rightarrow L_{2}\left(\mathbf{R}^{n} ; \bigwedge_{\|}^{k}\right), & \underline{d}: L_{2}\left(\mathbf{R}^{n} ; \bigwedge_{\|}^{k}\right) \rightarrow L_{2}\left(\mathbf{R}^{n} ; \bigwedge_{\perp}^{k+2}\right) \\
\underline{d}^{*}: L_{2}\left(\mathbf{R}^{n} ; \bigwedge_{\|}^{k}\right) \rightarrow L_{2}\left(\mathbf{R}^{n} ; \bigwedge_{\perp}^{k}\right), & \underline{d}^{*}: L_{2}\left(\mathbf{R}^{n} ; \bigwedge_{\perp}^{k}\right) \rightarrow L_{2}\left(\mathbf{R}^{n} ; \bigwedge_{\|}^{k-2}\right) .
\end{aligned}
$$

Lemma 2.7. We have, on appropriate domains, $\left(d_{t, x}\right)^{2}=\left(d_{t, x}^{*}\right)^{2}=d^{2}=\left(d^{*}\right)^{2}=\mu^{2}=$ $\left(\mu^{*}\right)^{2}=0$ and the anti-commutation relations

$$
\begin{array}{ll}
\left\{m, d_{t, x}\right\}=\partial_{t}=\left\{m,-d_{t, x}^{*}\right\}, & m^{2}=\left\{\mu, \mu^{*}\right\}=I, \\
\left\{d, d^{*}\right\}=-\Delta=-\sum_{1}^{n} \partial_{k}^{2}, & \{d, \mu\}=\left\{d, \mu^{*}\right\}=0,
\end{array}
$$

where $\{A, B\}=A B+B A$ denotes the anticommutator.

Proof. The proofs are straightforward, using Lemma 2.3. Let us prove that $\partial_{t}=\left\{m,-d_{t, x}^{*}\right\}$. For a function $F(t, x)$ we have

$$
\left.\left.\left.\left.\left.\nabla_{x, t}\right\lrcorner(m F)=\sum_{0}^{n} e_{i}\right\lrcorner\left(m \partial_{i} F\right)=\sum_{0}^{n} e_{i}\right\lrcorner\left(e_{0}\right\lrcorner \partial_{i} F\right)+\sum_{0}^{n} e_{i}\right\lrcorner\left(e_{0} \wedge \partial_{i} F\right) .
$$

Using the anticommutativity and associative properties, the first sum is

$$
\left.\left.\left.\left.\sum_{0}^{n}\left(e_{0} \wedge e_{i}\right)\right\lrcorner \partial_{i} F=-\sum_{0}^{n}\left(e_{i} \wedge e_{0}\right)\right\lrcorner \partial_{i} F=-e_{0}\right\lrcorner\left(\sum_{0}^{n} e_{i}\right\lrcorner \partial_{i} F\right),
$$

whereas using the derivation property, the second sum is

$$
\left.\left.\sum_{0}^{n} e_{i}\right\lrcorner\left(e_{0} \wedge \partial_{i} F\right)=\partial_{0} F-e_{0} \wedge\left(\sum_{0}^{n} e_{i}\right\lrcorner \partial_{i} F\right)
$$

Adding up we obtain $\left.\left.\nabla_{x, t}\right\lrcorner(m F)=\partial_{t} F-m\left(\nabla_{x, t}\right\lrcorner F\right)$.

### 2.1. The basic operators

Definition 2.8. Throughout this paper we denote by $B \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}(\bigwedge)\right)$ a bounded, accretive and complex matrix function acting on $f \in \mathcal{H}$ as $f(x) \mapsto B(x) f(x)$, with quantitative bounds $\|B\|_{\infty}$ and $\kappa_{B}>0$, where $\kappa_{B}$ is the largest constant such that

$$
\operatorname{Re}(B(x) w, w) \geqslant \kappa_{B}|w|^{2}, \quad \text { for all } w \in \Lambda, x \in \mathbf{R}^{n}
$$

We shall also assume that $B$ is of the form $B=B^{0} \oplus B^{1} \oplus B^{2} \oplus \cdots \oplus B^{n+1}$, where $B^{k} \in$ $L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\bigwedge^{k}\right)\right)$, so that $B$ preserve the space of $k$-vectors. For the matrix part $B^{1}$ acting on vectors, we use the alternative notation $A=B^{1}$. We also define $\widetilde{B}:=m B m$. It is not true in general that $\widetilde{B}$ preserve $k$-vectors.

Remark 2.9. It would be more optimal to replace the quantitative bounds $\|B\|_{\infty}$ and $\kappa_{B}>0$ with $K_{B}$ and $k_{B}$, where $K_{B}$ and $k_{B}$ are the optimal constants such that

$$
k_{B}\|f\|^{2} \leqslant|(B f, f)| \leqslant K_{B}\|f\|^{2}, \quad \text { for all } f \in \mathcal{H} .
$$

We consider $F(t, x)$ satisfying the Dirac type equation

$$
\begin{equation*}
\left(m d_{t, x}+B^{-1} m d_{t, x}^{*} B\right) F(t, x)=0 \tag{2.1}
\end{equation*}
$$

or equivalently

$$
\begin{equation*}
\left(d_{t, x}+\widetilde{B}^{-1} d_{t, x}^{*} B\right) F(t, x)=0 \tag{2.2}
\end{equation*}
$$

In order to solve for the vertical derivative $\partial_{t} F$, we note that

$$
d_{t, x} F=d F+\mu \partial_{t} F, \quad d_{t, x}^{*} F=d^{*} F-\mu^{*} \partial_{t} F
$$

Inserted in (2.1) this yields

$$
\partial_{t} F+\frac{1}{N^{+}-B^{-1} N^{-B}}\left(m d+B^{-1} m d^{*} B\right) F=0 .
$$

Definition 2.10. Write $M_{B}:=N^{+}-B^{-1} N^{-} B$ and define the unbounded operator

$$
T_{B}:=M_{B}^{-1}\left(m d+B^{-1} m d^{*} B\right)=-i M_{B}^{-1}\left(\underline{d}+B^{-1} \underline{d}^{*} B\right)=\frac{1}{\mu-\widetilde{B}^{-1} \mu^{*} B}\left(d+\widetilde{B}^{-1} d^{*} B\right)
$$

in $\mathcal{H}$ with domain $\mathrm{D}\left(T_{B}\right):=\mathrm{D}(d) \cap B^{-1} \mathrm{D}\left(d^{*}\right)$.

A rather surprising fact is that the most obvious choice for $\widetilde{B}$, namely $\widetilde{B}=B$ is not the best, but rather $\widetilde{B}=m B m$. For example, this is the only choice for which a Rellich type formula, as in Proposition 3.8, holds on all $\mathcal{H}$, in the case of Hermitean coefficients $B$.

Note that $T_{B}$ is closely related to operators of the form $\Pi_{B}=\Gamma+B^{-1} \Gamma^{*} B$, where $\Gamma$ denotes a first order, homogeneous partial differential operator with constant coefficients such that $\Gamma^{2}=0$, which were studied in [10]. Unfortunately, the factor $M_{B}^{-1}$ does not commute with $\Pi_{B}$ for general $B$. However, it has other useful commutation properties.

Lemma 2.11. The operator $M_{B}$ is an isomorphism and

$$
\begin{array}{ll}
\left(B^{-1} N^{-} B\right) M_{B}^{-1}=M_{B}^{-1} N^{-}, & N^{-} M_{B}^{-1}=M_{B}^{-1}\left(B^{-1} N^{-} B\right) \\
\left(B^{-1} N^{+} B\right) M_{B}^{-1}=M_{B}^{-1} N^{+}, & N^{+} M_{B}^{-1}=M_{B}^{-1}\left(B^{-1} N^{+} B\right)
\end{array}
$$

Proof. Note that $M_{B}=B^{-1}\left(B N^{+}-N^{-} B\right)$, where the last factor is the diagonal matrix $\left(-B_{\perp \perp}\right) \oplus B_{\| \|}$in the splitting $\mathcal{H}=N^{-} \mathcal{H} \oplus N^{+} \mathcal{H}$, if $B_{\perp \perp}$ and $B_{\| \| \|}$denote the diagonal blocks of $B$, and thus $B N^{+}-N^{-} B$ commutes with $N^{ \pm}$. Furthermore, the diagonal blocks are accretive, so $B N^{+}-N^{-} B$ and thus $M_{B}$ is invertible. This proves the two equations to the right. To obtain the equations to the left, replace $B$ by $B^{-1}$ and note that

$$
-N^{-}+B^{-1} N^{+} B=N^{+}-B^{-1} N^{-} B
$$

Let us comment on the terminology "Dirac type equation" for (2.1). Normally this denotes a first-order differential operator, like $d_{t, x}+d_{t, x}^{*}$, whose square acts componentwise as the Laplace operator. In our situation, the following holds.

Lemma 2.12. If $F(t, x)$ satisfies (2.1), then $d_{t, x}^{*} B d_{t, x}(m F)=0$. In particular it follows that

$$
\operatorname{div}_{t, x} A \nabla_{t, x} F^{1,0}=0
$$

if $F=F^{0}+\left(F^{1,0} e_{0}+F^{1, \|}\right)+F^{2}+\cdots+F^{n+1}$, i.e. $F^{1,0}$ is the normal component of the vector part of $F$.

Proof. We use the anticommutation relations $m d_{t, x}=\partial_{t}-d_{t, x} m$ and $m d_{t, x}^{*}=-\partial_{t}-d_{t, x}^{*} m$ from Lemma 2.7, which shows that (2.1) is equivalent with

$$
\left(d_{t, x} m+B^{-1} d_{t, x}^{*} m B\right) F(t, x)=0
$$

since $\partial_{t}-B^{-1} \partial_{t} B=0$. Applying $d_{t, x}^{*} B$ to this equation shows that $d_{t, x}^{*} B d_{t, x}(m F)=0$, and evaluating the scalar $\bigwedge^{0}$ part shows that $\operatorname{div}_{t, x} A \nabla_{t, x} F^{1,0}=0$ since $d_{t, x}^{*} B d_{t, x}$ preserves $k$ vectors.

The following notions are central in our operator theoretic framework.

Definition 2.13. Let $\mathcal{H}$ be a Hilbert space.
(i) A (topological) splitting of $\mathcal{H}$ is a decomposition $\mathcal{H}=\mathcal{H}_{1} \oplus \mathcal{H}_{2}$ into closed subspaces $\mathcal{H}_{1}$ and $\mathcal{H}_{2}$. In particular, we have $\left\|f_{1}+f_{2}\right\| \approx\left\|f_{1}\right\|+\left\|f_{2}\right\|$ if $f_{i} \in \mathcal{H}_{i}$.
(ii) Two bounded operators $R^{+}$and $R^{-}$in $\mathcal{H}$ are called complementary projections if $R^{+}+$ $R^{-}=I,\left(R^{ \pm}\right)^{2}=R^{ \pm}$and $R^{ \pm} R^{\mp}=0$.
(iii) A bounded operator $R$ in $\mathcal{H}$ is called a reflection operator if $R^{2}=I$.

We note the following connection between these concepts.
Lemma 2.14. There is a one-to-one correspondence

$$
R=R^{+}-R^{-} \longleftrightarrow R^{ \pm}=\frac{1}{2}(I \pm R) \longleftrightarrow \mathcal{H}=R^{+} \mathcal{H} \oplus R^{-} \mathcal{H}
$$

between reflection operators in $\mathcal{H}$, complementary projections in $\mathcal{H}$ and topological splittings of $\mathcal{H}$. We write $R^{ \pm} \mathcal{H}=\mathrm{R}\left(R^{ \pm}\right)$for the range of the projection $R^{ \pm}$.

In Definition 2.4 we introduced the complementary projections $N^{ \pm}$associated with the splitting of $\mathcal{H}$ into the subspaces of tangential and normal multivector fields. The corresponding reflection operator is

$$
N:=N^{+}-N^{-}=\mu^{*} \mu-\mu \mu^{*}=\left(\mu^{*}-\mu\right) m=m\left(\mu-\mu^{*}\right)
$$

We also introduce $B$-perturbed versions of the tangential and normal subspaces $N^{-} \mathcal{H}$ and $N^{+} \mathcal{H}$ as

$$
\begin{aligned}
B^{-1} N^{+} \mathcal{H} & :=\left\{B^{-1} f ; f \in N^{+} \mathcal{H}\right\} \\
B^{-1} N^{-} \mathcal{H} & :=\left\{B^{-1} f ; f \in N^{-} \mathcal{H}\right\}
\end{aligned}
$$

In Definition 2.10 we encountered one of the complementary projections $B^{-1} N^{+} B$ and $B^{-1} N^{-} B$ associated with the splitting

$$
\mathcal{H}=B^{-1} N^{+} \mathcal{H} \oplus B^{-1} N^{-} \mathcal{H}
$$

However, more important will be the following complementary projections.
Definition 2.15. Let $\hat{N}_{B}^{+}$and $\hat{N}_{B}^{-}$be the complementary projections associated with the splitting

$$
\mathcal{H}=B^{-1} N^{+} \mathcal{H} \oplus N^{-} \mathcal{H}
$$

We sometimes use the shorter notation $N_{B}^{+}:=\hat{N}_{B}^{+}$and $N_{B}^{-}:=\hat{N}_{B}^{-}$. Also let $\check{N}_{B}^{+}$and $\check{N}_{B}^{-}$be the complementary projections associated with the splitting

$$
\mathcal{H}=N^{+} \mathcal{H} \oplus B^{-1} N^{-} \mathcal{H}
$$

Let $N_{B}=\hat{N}_{B}:=\hat{N}_{B}^{+}-\hat{N}_{B}^{-}$and $\check{N}_{B}:=\check{N}_{B}^{+}-\check{N}_{B}^{-}$be the associated reflection operators.

With the notation $\mu_{B}^{*}:=B^{-1} \mu^{*} B$, these operators are

$$
\begin{aligned}
& \hat{N}_{B}^{+}=\mu_{B}^{*}\left(\mu+\mu_{B}^{*}\right)^{-1}=\left(\mu+\mu_{B}^{*}\right)^{-1} \mu \\
& \hat{N}_{B}^{-}=\mu\left(\mu+\mu_{B}^{*}\right)^{-1}=\left(\mu+\mu_{B}^{*}\right)^{-1} \mu_{B}^{*} \\
& \hat{N}_{B}=\left(\mu_{B}^{*}-\mu\right)\left(\mu+\mu_{B}^{*}\right)^{-1}=\left(\mu+\mu_{B}^{*}\right)^{-1}\left(\mu-\mu_{B}^{*}\right),
\end{aligned}
$$

and similarly for $\check{N}_{B}^{ \pm}$and $\check{N}_{B}$. We shall prove in Section 2.2 that all these operators are bounded. It is mainly the operators $\hat{N}_{B}^{+}, \hat{N}_{B}^{-}$and $\hat{N}_{B}$ that we shall use.

Definition 2.16. Let $\langle f, g\rangle_{B}:=\left(\left(B N^{+}-N^{-} B\right) f, g\right)$, for $f, g \in \mathcal{H}$. As $B N^{+}-N^{-} B=B M_{B}$ is invertible, $\langle\cdot, \cdot\rangle_{B}$ is a duality, i.e. there exists $C<\infty$ such that

$$
\begin{gathered}
\left|\langle f, g\rangle_{B}\right| \leqslant C\|f\|\|g\|, \\
\|f\| \leqslant C \sup _{g \neq 0}\left|\langle f, g\rangle_{B}\right| /\|g\|, \quad\|g\| \leqslant C \sup _{f \neq 0}\left|\langle f, g\rangle_{B}\right| /\|f\| .
\end{gathered}
$$

We write $T^{\prime}$ for the adjoint of an operator $T$ with respect to this duality, i.e. if $\langle T f, g\rangle_{B}=$ $\left\langle f, T^{\prime} g\right\rangle_{B}$ for all $f, g \in \mathcal{H}$.

Proposition 2.17. We have adjoint operators $\left(\hat{N}_{B}\right)^{\prime}=\check{N}_{B^{*}},\left(\check{N}_{B}\right)^{\prime}=\hat{N}_{B^{*}}, N^{\prime}=N$ and $\left(T_{B}\right)^{\prime}=$ $-T_{B^{*}}$.

Proof. To prove that $\left(\hat{N}_{B}\right)^{\prime}=\check{N}_{B^{*}}$, recall that $\hat{N}_{B}$ is the reflection operator for the splitting $\mathcal{H}=$
 Thus we need to prove that

$$
\begin{aligned}
& \left\langle B^{-1} N^{+} f_{1}-N^{-} f_{2}, N^{+} g_{1}+\left(B^{*}\right)^{-1} N^{-} g_{2}\right\rangle_{B} \\
& \quad=\left\langle B^{-1} N^{+} f_{1}+N^{-} f_{2}, N^{+} g_{1}-\left(B^{*}\right)^{-1} N^{-} g_{2}\right\rangle_{B},
\end{aligned}
$$

for all $f_{i}, g_{i}$, i.e. that $\left\langle B^{-1} N^{+} f_{1},\left(B^{*}\right)^{-1} N^{-} g_{2}\right\rangle_{B}=0=\left\langle N^{-} f_{2}, N^{+} g_{1}\right\rangle_{B}$. Use Lemma 2.11 to obtain

$$
\begin{aligned}
\left\langle B^{-1} N^{+} f_{1},\left(B^{*}\right)^{-1} N^{-} g_{2}\right\rangle_{B} & =\left(\left(B N^{+}-N^{-} B\right) B^{-1} N^{+} f_{1},\left(B^{*}\right)^{-1} N^{-} g_{2}\right) \\
& =\left(N^{-} M_{B} B^{-1} N^{+} f_{1}, g_{2}\right)=\left(M_{B} B^{-1} N^{-} N^{+} f_{1}, g_{2}\right)=0 .
\end{aligned}
$$

A similar calculation shows that $\left\langle N^{-} f_{2}, N^{+} g_{1}\right\rangle_{B}=0$. The proof of $\left(\check{N}_{B}\right)^{\prime}=\hat{N}_{B^{*}}$ is similar. For the unperturbed reflection operator $N$, the duality $N^{\prime}=N$ follows directly from the fact the $B N^{+}-N^{-} B$ is diagonal in the splitting $\mathcal{H}=N^{+} \mathcal{H} \oplus N^{-} \mathcal{H}$.

To prove $\left(T_{B}\right)^{\prime}=-T_{B^{*}}$, we calculate

$$
\begin{aligned}
\left\langle T_{B} f, g\right\rangle_{B} & =\left(\left(B N^{+}-N^{-} B\right) M_{B}^{-1}\left(m d+B^{-1} m d^{*} B\right) f, g\right) \\
& =\left(\left(B m d+m d^{*} B\right) f, g\right)=\left(f,\left(B^{*} d m+d^{*} m B^{*}\right) g\right)
\end{aligned}
$$

$$
\begin{aligned}
& =-\left(f,\left(N^{+} B^{*}-B^{*} N^{-}\right)\left(N^{+} B^{*}-B^{*} N^{-}\right)^{-1}\left(B^{*} m d+m d^{*} B^{*}\right) g\right) \\
& =-\left\langle f,\left(N^{+}-\left(B^{*}\right)^{-1} N^{-} B^{*}\right)^{-1}\left(m d+\left(B^{*}\right)^{-1} m d^{*} B^{*}\right) g\right\rangle_{B}=\left\langle f,-T_{B^{*}} g\right\rangle_{B}
\end{aligned}
$$

where we have used that $\left(B^{*}\right)^{-1} N^{+} B^{*}-N^{-}=N^{+}-\left(B^{*}\right)^{-1} N^{-} B^{*}$.

Remark 2.18. Our main use of these dualities is for proving surjectivity. Recall the following standard technique for proving invertibility of a bounded operator $T: \mathcal{H}_{1} \rightarrow \mathcal{H}_{2}$. Assume we have at our disposal two pairs of dual spaces $\left\langle\mathcal{H}_{1}, \mathcal{K}_{1}\right\rangle_{1}$ and $\left\langle\mathcal{H}_{2}, \mathcal{K}_{2}\right\rangle_{2}$ and that the adjoint operator is $T^{\prime}: \mathcal{K}_{2} \rightarrow \mathcal{K}_{1}$. If we can prove a priori estimates

$$
\|T f\| \gtrsim\|f\|, \quad \text { for all } f \in \mathcal{H}_{1}
$$

then $T$ is injective and has closed range. If furthermore $T^{\prime}$ is injective, in particular if it satisfies a priori estimates, then $T$ is surjective and therefore an isomorphism.

Remark 2.19. We shall also need to restrict dualities to subspaces. Let $\langle\mathcal{H}, \mathcal{K}\rangle$ be a duality, i.e. let $\langle\cdot, \cdot\rangle: \mathcal{H} \times \mathcal{K} \rightarrow \mathbf{C}$ satisfy the estimates in Definition 2.16. If $\mathcal{H}_{1} \subset \mathcal{H}$ is a subspace, then a subspace $\mathcal{K}_{1} \subset \mathcal{K}$ is such that $\left\langle\mathcal{H}_{1}, \mathcal{K}_{1}\right\rangle$ satisfies estimates as in Definition 2.16 if and only if $\mathcal{K}_{1}$ is a complementary subspace to the annihilator $\left\{g \in \mathcal{K} ;\langle f, g\rangle=0\right.$, for all $\left.f \in \mathcal{H}_{1}\right\}$, in the sense of Definition 2.13(i).

In particular, if $R^{ \pm}$are complementary projections in $\mathcal{H}$, the adjoint operators $\left(R^{ \pm}\right)^{\prime}$ are also complementary projections and the duality $\langle\mathcal{H}, \mathcal{H}\rangle$ restricts to a duality $\left\langle R^{+} \mathcal{H},\left(R^{+}\right)^{\prime} \mathcal{H}\right\rangle$, since the annihilator of $R^{+} \mathcal{H}$ is $\mathrm{N}\left(\left(R^{+}\right)^{\prime}\right)$, which is complementary to $\left(R^{+}\right)^{\prime} \mathcal{H}$.

### 2.2. Hodge decompositions and resolvent estimates

In this section, we estimate the spectrum of the operator $T_{B}$. For this we make use of Hodge type decompositions of $\mathcal{H}$ as explained below.

Definition 2.20. By a nilpotent operator $\Gamma$ in a Hilbert space $\mathcal{H}$, we mean a closed, densely defined operator such that $\overline{\mathrm{R}(\Gamma)} \subset \mathrm{N}(\Gamma)$. In particular $\Gamma^{2} f=0$ if $f \in \mathrm{D}(\Gamma)$. We say that a nilpotent operator is exact if $\overline{\mathrm{R}(\Gamma)}=\mathrm{N}(\Gamma)$.

If $\widetilde{\Gamma}$ is another nilpotent operator, then we say that $\Gamma$ and $\widetilde{\Gamma}$ are transversal if there is a constant $c=c(\Gamma, \widetilde{\Gamma})<1$ such that

$$
|(f, g)| \leqslant c\|f\|\|g\|, \quad f \in \overline{\mathrm{R}(\Gamma)}, g \in \overline{\mathrm{R}(\widetilde{\Gamma})}
$$

or equivalently if $\|f+g\| \approx\|f\|+\|g\|$ for all $f \in \overline{\mathrm{R}(\Gamma)}$ and $g \in \overline{\mathrm{R}(\widetilde{\Gamma})}$. Note that any nilpotent operator $\Gamma$ is transversal to its adjoint $\Gamma^{*}$ with $c=0$, since $\overline{\mathrm{R}(\Gamma)} \subset \mathrm{N}(\Gamma)=\overline{\mathrm{R}\left(\Gamma^{*}\right)^{\perp}}$.

Below we collect the properties of Hodge type splittings which we need in this paper. This generalises results from [10, Proposition 2.2] and [9, Proposition 3.11].

Lemma 2.21. Let $\Gamma$ and $\widetilde{\Gamma}$ be two nilpotent operators in a Hilbert space $\mathcal{H}$ which are exact and transversal with constant $c(\Gamma, \widetilde{\Gamma})$, and assume also that the adjoints $\Gamma^{*}$ and $\widetilde{\Gamma}^{*}$ are transversal. Let $B$ be a bounded, accretive multiplication operator and assume that

$$
\max \left(c(\Gamma, \widetilde{\Gamma}), c\left(\Gamma^{*}, \widetilde{\Gamma}^{*}\right)\right)<\kappa_{B} /\|B\|_{\infty}
$$

Define operators $\Pi:=\Gamma+\widetilde{\Gamma}, \widetilde{\Gamma}_{B}:=B^{-1} \widetilde{\Gamma} B$ and $\Pi_{B}:={\underset{\sim}{\tilde{C}}}_{B}+\widetilde{\Gamma}_{B}$ with domains $\mathrm{D}(\Pi):=$ $\mathrm{D}(\Gamma) \cap \mathrm{D}(\widetilde{\Gamma}), \mathrm{D}\left(\widetilde{\Gamma}_{B}\right):=B^{-1} \mathrm{D}(\widetilde{\Gamma})$ and $\mathrm{D}\left(\Pi_{B}\right):=\mathrm{D}(\Gamma) \cap \mathrm{D}\left(\widetilde{\Gamma}_{B}\right)$, respectively.
(i) We have a topological splitting

$$
\mathcal{H}=\mathrm{N}(\Gamma) \oplus \mathrm{N}\left(\widetilde{\Gamma}_{B}\right),
$$

so that $\left\|f_{1}+f_{2}\right\| \approx\left\|f_{1}\right\|+\left\|f_{2}\right\|, f_{1} \in \mathrm{~N}(\Gamma), f_{2} \in \mathrm{~N}\left(\widetilde{\Gamma}_{B}\right)$. The operator $\Pi_{B}$ is a closed operator with dense domain and range. Furthermore $\widetilde{\Gamma}_{B}$ is an exact nilpotent operator. The complementary Hodge type projections associated with the splitting are

$$
\begin{aligned}
\mathbf{P}_{B}^{1} & :=\Gamma \Pi_{B}^{-1}=\Pi_{B}^{-1} \widetilde{\Gamma}_{B}=\Gamma \Pi_{B}^{-2} \widetilde{\Gamma}_{B} \\
\mathbf{P}_{B}^{2} & :=\widetilde{\Gamma}_{B} \Pi_{B}^{-1}=\Pi_{B}^{-1} \Gamma=\widetilde{\Gamma}_{B} \Pi_{B}^{-2} \Gamma
\end{aligned}
$$

where we identify a bounded, densely defined operator with its bounded extension to $\mathcal{H}$. In the special case when $B=I$ we write $\mathbf{P}^{1}$ and $\mathbf{P}^{2}$ for these projections.
(ii) Let $B_{1}$ and $B_{2}$ be two bounded, accretive operators in $\mathcal{H}$, with $\kappa_{B_{i}} /\left\|B_{i}\right\|_{\infty}>\max (c(\Gamma, \widetilde{\Gamma})$, $\left.c\left(\Gamma^{*}, \widetilde{\Gamma}^{*}\right)\right), i=1,2$. Then there exists $C<\infty$, depending only on $\left\|B_{i}\right\|_{\infty}, \kappa_{B_{i}}, c(\Gamma, \widetilde{\Gamma})$ and $c\left(\Gamma^{*}, \widetilde{\Gamma}^{*}\right)$, such that

$$
\left\|\left(\Gamma+B_{1}^{-1} \widetilde{\Gamma} B_{2}\right) f\right\| \geqslant C^{-1} \lambda\|f\|, \quad \text { for all } f \in \mathrm{D}(\Gamma) \cap B_{2}^{-1} \mathrm{D}(\widetilde{\Gamma}),
$$

if $\|\Pi f\| \geqslant \lambda\|f\|$ for all $f \in \mathrm{D}(\Pi)$.
Proof. Since $\widetilde{\Gamma}_{B}$ is conjugated to $\widetilde{\Gamma}$, it is an exact nilpotent operator. To prove (i), it suffices to prove the estimate

$$
\begin{equation*}
\|f\|+\|g\| \lesssim\|f+g\|, \quad f \in \mathrm{~N}(\Gamma), g \in \mathrm{~N}\left(\widetilde{\Gamma}_{B}\right) . \tag{2.3}
\end{equation*}
$$

Indeed, this shows that $\mathrm{N}(\Gamma) \oplus \mathrm{N}\left(\widetilde{\Gamma}_{B}\right) \subset \mathcal{H}$. Furthermore, replacing $(\Gamma, \widetilde{\Gamma}, B)$ with $\left(\widetilde{\Gamma}^{*}, \Gamma^{*}, B^{*}\right)$ shows that $\mathrm{N}\left(\widetilde{\Gamma}^{*}\right) \oplus \mathrm{N}\left(\Gamma_{B^{*}}^{*}\right) \subset \mathcal{H}$. Conjugating with $B^{*}$ then shows that $\mathrm{N}\left(\left(\widetilde{\Gamma}_{B}\right)^{*}\right) \oplus \mathrm{N}\left(\Gamma^{*}\right) \subset \mathcal{H}$. Therefore a duality argument proves the splitting $\mathcal{H}=\mathrm{N}(\Gamma) \oplus \mathrm{N}\left(\widetilde{\Gamma}_{B}\right)$. From this it follows that

$$
\mathrm{D}\left(\Pi_{B}\right)=\left(\mathrm{D}\left(\widetilde{\Gamma}_{B}\right) \cap \mathrm{N}(\Gamma)\right) \oplus\left(\mathrm{D}(\Gamma) \cap \mathrm{N}\left(\widetilde{\Gamma}_{B}\right)\right)
$$

and $\mathrm{R}\left(\Pi_{B}\right)=\mathrm{R}(\Gamma) \oplus \mathrm{R}\left(\widetilde{\Gamma}_{B}\right)$ are dense and that $\Pi_{B}$ is closed, as well as the boundedness of the associated projections.

To prove (2.3), we use that $B g \in \mathrm{~N}(\widetilde{\Gamma})$ and thus $|(f, B g)| \leqslant c(\Gamma, \widetilde{\Gamma})\|f\|\|B g\|$, and estimate:

$$
\begin{aligned}
\|f\|^{2} \leqslant \kappa_{B}^{-1} \operatorname{Re}(B f, f) & \leqslant \kappa_{B}^{-1} \operatorname{Re}((B(f+g), f)-(B g, f)) \\
& \leqslant \kappa_{B}^{-1}\|B\|_{\infty}\|f+g\|\|f\|+\kappa_{B}^{-1} c(\Gamma, \widetilde{\Gamma})\|B g\|\|f\| \\
& \leqslant \kappa_{B}^{-1}\|B\|_{\infty}\|f+g\|\|f\|+\kappa_{B}^{-1} c(\Gamma, \widetilde{\Gamma})\|B\|_{\infty}(\|f+g\|+\|f\|)\|f\| \\
& \leqslant \kappa_{B}^{-1}\|B\|_{\infty}(1+c(\Gamma, \widetilde{\Gamma}))\|f+g\|\|f\|+\kappa_{B}^{-1}\|B\|_{\infty} c(\Gamma, \widetilde{\Gamma})\|f\|^{2}
\end{aligned}
$$

Solving for $\|f\|^{2}$, this shows that $\|f\| \lesssim\|f+g\|$ provided $\kappa_{B}^{-1}\|B\|_{\infty} c(\Gamma, \widetilde{\Gamma})<1$, which proves (2.3).

To prove (ii) we factorise

$$
\Gamma+B_{1}^{-1} \widetilde{\Gamma} B_{2}=\left(\mathbf{P}^{1}+B_{1}^{-1} \mathbf{P}^{2}\right) \Pi\left(\mathbf{P}^{2}+\mathbf{P}^{1} B_{2}\right)
$$

Here $\left(\mathbf{P}^{1}+B_{1}^{-1} \mathbf{P}^{2}\right)^{-1}=\mathbf{P}_{B_{1}}^{1}+B_{1} \mathbf{P}_{B_{1}}^{2}$ and $\left(\mathbf{P}^{2}+\mathbf{P}^{1} B_{2}\right)^{-1}=\mathbf{P}_{B_{2}}^{1} B_{2}^{-1}+\mathbf{P}_{B_{2}}^{2}$. Indeed, we have

$$
\begin{aligned}
\left(\mathbf{P}^{1}+B_{1}^{-1} \mathbf{P}^{2}\right)\left(\mathbf{P}_{B_{1}}^{1}+B_{1} \mathbf{P}_{B_{1}}^{2}\right) & =\mathbf{P}^{1} \mathbf{P}_{B_{1}}^{1}+\mathbf{P}^{1} B_{1} \mathbf{P}_{B_{1}}^{2}+B_{1}^{-1} \mathbf{P}^{2} \mathbf{P}_{B_{1}}^{1}+B_{1}^{-1} \mathbf{P}^{2} B_{1} \mathbf{P}_{B_{1}}^{2} \\
& =\mathbf{P}^{1} \mathbf{P}_{B_{1}}^{1}+B_{1}^{-1} \mathbf{P}^{2} B_{1} \mathbf{P}_{B_{1}}^{2} \\
& =\left(\mathbf{P}^{1}+\mathbf{P}^{2}\right) \mathbf{P}_{B_{1}}^{1}+B_{1}^{-1}\left(\mathbf{P}^{1}+\mathbf{P}^{2}\right) B_{1} \mathbf{P}_{B_{1}}^{2} \\
& =\mathbf{P}_{B_{1}}^{1}+\mathbf{P}_{B_{1}}^{2}=I, \\
\left(\mathbf{P}^{2}+\mathbf{P}^{1} B_{2}\right)\left(\mathbf{P}_{B_{2}}^{1} B_{2}^{-1}+\mathbf{P}_{B_{2}}^{2}\right)= & \mathbf{P}^{2} \mathbf{P}_{B_{2}}^{1} B_{2}^{-1}+\mathbf{P}^{2} \mathbf{P}_{B_{2}}^{2}+\mathbf{P}^{1} B_{2} \mathbf{P}_{B_{2}}^{1} B_{2}^{-1}+\mathbf{P}^{1} B_{2} \mathbf{P}_{B_{2}}^{2} \\
& =\mathbf{P}^{2} \mathbf{P}_{B_{2}}^{2}+\mathbf{P}^{1} B_{2} \mathbf{P}_{B_{2}}^{1} B_{2}^{-1} \\
& =\mathbf{P}^{2}\left(\mathbf{P}_{B_{2}}^{1}+\mathbf{P}_{B_{2}}^{2}\right)+\mathbf{P}^{1} B_{2}\left(\mathbf{P}_{B_{2}}^{1}+\mathbf{P}_{B_{2}}^{2}\right) B_{2}^{-1}=\mathbf{P}^{2}+\mathbf{P}^{1}=I .
\end{aligned}
$$

A similar calculation shows that $\mathbf{P}_{B_{1}}^{1}+B_{1} \mathbf{P}_{B_{1}}^{2}$ and $\mathbf{P}_{B_{2}}^{1} B_{2}^{-1}+\mathbf{P}_{B_{2}}^{2}$ also are left inverses. Thus $\mathbf{P}^{1}+B_{1}^{-1} \mathbf{P}^{2}$ and $\mathbf{P}^{2}+\mathbf{P}^{1} B_{2}$ are isomorphisms and (ii) follows.

We can now prove the following perturbed normal and tangential splittings of $\mathcal{H}$.

$$
\begin{aligned}
& \mathcal{H}=\hat{N}_{B}^{+} \mathcal{H} \oplus \hat{N}_{B}^{-} \mathcal{H}=B^{-1} N^{+} \mathcal{H} \oplus N^{-} \mathcal{H} \\
& \mathcal{H}=\check{N}_{B}^{+} \mathcal{H} \oplus \check{N}_{B}^{-} \mathcal{H}=N^{+} \mathcal{H} \oplus B^{-1} N^{-} \mathcal{H}
\end{aligned}
$$

To see this, let first $\Gamma=\mu$ and $\widetilde{\Gamma}=\Gamma^{*}=\mu^{*}$ in Lemma 2.21(i). It follows that $\mathrm{N}(\Gamma)=N^{-} \mathcal{H}$, $\mathrm{N}\left(\Gamma_{B}^{*}\right)=B^{-1} N^{+} \mathcal{H}, \mathbf{P}_{B}^{1}=\hat{N}_{B}^{-}$and $\mathbf{P}_{B}^{2}=\hat{N}_{B}^{+}$. On the other hand, choosing $\Gamma=\mu^{*}$ and $\widetilde{\Gamma}=\Gamma^{*}=\mu$, we see that $\mathrm{N}(\Gamma)=N^{+} \mathcal{H}, \mathrm{N}\left(\Gamma_{B}^{*}\right)=B^{-1} N^{-} \mathcal{H}, \mathbf{P}_{B}^{1}=\check{N}_{B}^{+}$and $\mathbf{P}_{B}^{2}=\check{N}_{B}^{-}$. Lemma 2.21(i) thus shows that all the oblique normal and tangential projections $\hat{N}_{B}^{ \pm}$and $\check{N}_{B}^{ \pm}$ from Definition 2.15 are bounded operators on $\mathcal{H}$, i.e. we have the stated splittings.

Next we apply Lemma 2.21 (ii) to prove resolvent bounds for the operator $T_{B}$ from Definition 2.10. Define closed and open sectors and double sectors in the complex plane by

$$
\begin{aligned}
S_{\omega+} & :=\{z \in \mathbf{C} ;|\arg z| \leqslant \omega\} \cup\{0\}, & & S_{\omega}:=S_{\omega+} \cup\left(-S_{\omega+}\right), \\
S_{v+}^{o} & :=\{z \in \mathbf{C} ; z \neq 0,|\arg z|<v\}, & & S_{v}^{o}:=S_{v+}^{o} \cup\left(-S_{v+}^{o}\right) .
\end{aligned}
$$

Proposition 2.22. The operator $T_{B}$ is a closed operator in $\mathcal{H}$ with dense domain and range, for any accretive, complex matrix function $B \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}(\bigwedge)\right)$. Furthermore, $T_{B}$ is a bisectorial operator with $\sigma\left(T_{B}\right) \subset S_{\omega}$, where

$$
\omega:=\arccos \left(\kappa_{B} /\left(2\|B\|_{\infty}\right)\right) \in[\pi / 3, \pi / 2)
$$

and if $\omega<\nu<\pi / 2$, then there exists $C<\infty$ depending only on $\nu, \kappa_{B}$ and $\|B\|_{\infty}$ such that

$$
\left\|\left(\lambda-T_{B}\right)^{-1}\right\| \leqslant C /|\lambda|, \quad \text { for all } \lambda \notin S_{\nu} .
$$

Proof. It is a consequence of Lemma 2.21(i) that the operator $\Pi_{B}=\underline{d}+B^{-1} \underline{d}^{*} B$ is a closed operator with dense domain and range. Since $T_{B}=-i M_{B}^{-1} \Pi_{B}$ with $M_{B}$ an isomorphism, it follows that $T_{B}$ also is closed with dense domain and range.

To prove the resolvent estimate, write $\lambda=1 /(i \tau)$ where $\tau \in S_{\pi / 2-v}^{o}$. We first prove that $\|u\| \leqslant$ $C\|f\|$ if $\left(I-i \tau T_{B}\right) u=f$, uniformly for $\tau \in S_{\pi / 2-v}^{o}$. Multiply the equation with $i\left(\mu-\widetilde{B}^{-1} \mu^{*} B\right)$ to obtain

$$
\left(\Gamma+\widetilde{B}^{-1} \widetilde{\Gamma} B\right) u=i\left(\mu-\widetilde{B}^{-1} \mu^{*} B\right) f
$$

where $\Gamma:=i \mu+\tau d$ and $\widetilde{\Gamma}:=-i \mu^{*}+\tau d^{*}$ are nilpotent by Lemma 2.7. It suffices to prove $\|u\| \lesssim\left\|\left(\Gamma+\widetilde{B}^{-1} \widetilde{\Gamma} B\right) u\right\|$.
(i) By orthogonality we have

$$
\begin{aligned}
\left\|\left(\Gamma+\Gamma^{*}\right) u\right\|^{2}=\|\Gamma u\|^{2}+\left\|\Gamma^{*} u\right\|^{2}= & \|\mu u\|^{2}+|\tau|^{2}\|d u\|^{2}+2 \operatorname{Re}(i \mu u, \tau d u) \\
& +\left\|\mu^{*} u\right\|^{2}+|\tau|^{2}\left\|d^{*} u\right\|^{2}+2 \operatorname{Re}\left(-i \mu^{*} u, \bar{\tau} d^{*} u\right),
\end{aligned}
$$

where

$$
\begin{aligned}
\operatorname{Re}(i \mu u, \tau d u)+\operatorname{Re}\left(-i \mu^{*} u, \bar{\tau} d^{*} u\right) & =\operatorname{Re}\left(i \bar{\tau} d^{*} \mu u, u\right)+\operatorname{Re}\left(u, i \bar{\tau} \mu d^{*} u\right) \\
& =\operatorname{Re}\left(i \bar{\tau}\left\{d^{*}, \mu\right\} u, u\right)=0,
\end{aligned}
$$

by Lemma 2.7. Thus $\left\|\left(\Gamma+\Gamma^{*}\right) u\right\|^{2}=\left\|\left(\mu+\mu^{*}\right) u\right\|^{2}+|\tau|^{2}\left\|\left(d+d^{*}\right) u\right\|^{2} \geqslant\|m u\|^{2}=\|u\|^{2}$. In particular $\Gamma$ is exact.
(ii) Next we prove that $\Gamma$ and $\widetilde{\Gamma}$ are transversal, with a bound $c<1$ uniformly for all $\tau \in$ $S_{\pi / 2-v}^{o}$. By exactness, it suffices to bound $(f, g)$ for $f=\Gamma u \in \mathrm{R}(\Gamma)$ and $\widetilde{\Gamma} g=0$. Furthermore, using the orthogonal Hodge splitting $\mathcal{H}=\mathrm{N}(\Gamma) \oplus \mathrm{N}\left(\Gamma^{*}\right)$, we may assume that $\Gamma^{*} u=0$. We get

$$
\begin{aligned}
(f, g) & =(\Gamma u, g)=\left(u, \Gamma^{*} g\right)=\left(u,\left(-i \mu^{*}+\bar{\tau} d^{*}\right) g\right) \\
& =\left(u,-i \mu^{*} g+i(\bar{\tau} / \tau) \mu^{*} g\right)=2(\tau-\bar{\tau}) /(2 i \bar{\tau})\left(u, \mu^{*} g\right),
\end{aligned}
$$

and thus $|(f, g)| \leqslant 2|\sin (\arg \tau)|\|u\|\|g\| \leqslant c\|f\|\|g\|$ by (i), where $c<\kappa_{B} /\|B\|_{\infty}$ since $\pi / 2-v<$ $\arcsin \left(\kappa_{B} /\left(2\|B\|_{\infty}\right)\right)$. A similar argument shows that $\Gamma^{*}$ and $\widetilde{\Gamma}^{*}$ are transversal with the same constant $c<\kappa_{B} /\|B\|_{\infty}$ uniformly for $\tau \in S_{\pi / 2-v}^{o}$.
(iii) To apply Lemma 2.21 (ii) it now suffices to prove that $\|(\Gamma+\widetilde{\Gamma}) u\| \geqslant C^{-1}\|u\|$ uniformly for all $\tau \in S_{\pi / 2-v}^{o}$. From Lemma 2.21(i) with $B=I$ we have $\|(\Gamma+\widetilde{\Gamma}) u\| \approx\|\Gamma u\|+\|\widetilde{\Gamma} u\|$.

Using the Hodge splitting $\mathcal{H}=\mathrm{N}(\Gamma) \oplus \mathrm{N}(\widetilde{\Gamma})$ it suffices to prove $\|\Gamma u\| \gtrsim\|u\|$ for $u \in \mathrm{~N}(\widetilde{\Gamma})$, and $\|\widetilde{\Gamma} v\| \gtrsim\|v\|$ for $v \in \mathrm{~N}(\Gamma)$. To prove for example the first estimate, write $\mathrm{N}(\widetilde{\Gamma}) \ni u=u_{1}+u_{2} \in$ $\mathrm{N}(\Gamma) \oplus \mathrm{N}\left(\Gamma^{*}\right)$. Then

$$
\begin{aligned}
\|\Gamma u\|^{2} & =\left\|\Gamma u_{2}\right\|^{2} \geqslant\left\|u_{2}\right\|^{2}=\left\|u-u_{1}\right\|^{2} \\
& \geqslant\|u\|^{2}+\left\|u_{1}\right\|^{2}-2 c\|u\|\left\|u_{1}\right\| \geqslant\left(1-c^{2}\right)\|u\|^{2}
\end{aligned}
$$

where we have used (i) in the second step and (ii) in the fourth step.
This proves that $\|u\| \leqslant C\|f\|$ if $\left(I-i \tau T_{B}\right) u=f$. Since $\left(I-i \tau T_{B}\right)^{\prime}=I-i \bar{\tau} T_{B^{*}}$ by Proposition 2.17 , a duality argument shows that $I-i \tau T_{B}$ is onto, and the proof is complete.

From the uniform boundedness of the resolvents $R_{t}^{B}:=\left(I+i t T_{B}\right)^{-1}$ for $t \in \mathbf{R}$ and the boundedness of the Hodge projections $\mathbf{P}_{B}^{1}:=\underline{d} \Pi_{B}^{-1}$ and $\mathbf{P}_{B}^{2}:=\underline{d}_{B}^{*} \Pi_{B}^{-1}$, where $\Pi_{B}=\underline{d}+\underline{d}_{B}^{*}, \underline{d}=$ imd and $\underline{d}_{B}^{*}=B^{-1} \underline{d}^{*} B$, we can now deduce boundedness of operators related to $P_{t}^{B}=\frac{1}{2}\left(R_{-t}^{B}+R_{t}^{B}\right)$ and $Q_{t}^{B}=\frac{1}{2 i}\left(R_{-t}^{B}-R_{t}^{B}\right)$.

Corollary 2.23. The following families of operators are all uniformly bounded for $t>0$.

$$
\begin{aligned}
& R_{t}^{B}:=\left(I+i t T_{B}\right)^{-1}, \\
& t \underline{d} P_{t}^{B}=i \mathbf{P}_{B}^{1} M_{B} Q_{t}^{B}, \\
& P_{t}^{B}:=\left(I+t^{2} T_{B}^{2}\right)^{-1}, \\
& t \underline{d}_{B}^{*} P_{t}^{B}=i \mathbf{P}_{B}^{2} M_{B} Q_{t}^{B}, \\
& Q_{t}^{B}:=t T_{B}\left(I+t^{2} T_{B}^{2}\right)^{-1}, \\
& P_{t}^{B} M_{B}^{-1} t \underline{d}=i Q_{t}^{B} \mathbf{P}_{B}^{2}, \\
& t^{2} T_{B}^{2} P_{t}^{B}=t T_{B} Q_{t}^{B}=I-P_{t}^{B}, \\
& P_{t}^{B} M_{B}^{-1} t \underline{d}_{B}^{*}=i Q_{t}^{B} \mathbf{P}_{B}^{1}, \\
& t \underline{d} Q_{t}^{B}=i \mathbf{P}_{B}^{1} M_{B}\left(I-P_{t}^{B}\right), \\
& t \underline{d} P_{t}^{B} M_{B}^{-1} t \underline{d}=\mathbf{P}_{B}^{1} M_{B}\left(P_{t}^{B}-I\right) \mathbf{P}_{B}^{2}, \\
& t \underline{d}_{B}^{*} Q_{t}^{B}=i \mathbf{P}_{B}^{2} M_{B}\left(I-P_{t}^{B}\right), \\
& t \underline{d} P_{t}^{B} M_{B}^{-1} t \underline{d}_{B}^{*}=\mathbf{P}_{B}^{1} M_{B}\left(P_{t}^{B}-I\right) \mathbf{P}_{B}^{1}, \\
& Q_{t}^{B} M_{B}^{-1} t \underline{d}=i\left(I-P_{t}^{B}\right) \mathbf{P}_{B}^{2}, \\
& t \underline{d}_{B}^{*} P_{t}^{B} M_{B}^{-1} t \underline{d}=\mathbf{P}_{B}^{2} M_{B}\left(P_{t}^{B}-I\right) \mathbf{P}_{B}^{2}, \\
& Q_{t}^{B} M_{B}^{-1} t \underline{d}_{B}^{*}=i\left(I-P_{t}^{B}\right) \mathbf{P}_{B}^{1}, \\
& t \underline{d}_{B}^{*} P_{t}^{B} M_{B}^{-1} t \underline{d}_{B}^{*}=\mathbf{P}_{B}^{2} M_{B}\left(P_{t}^{B}-I\right) \mathbf{P}_{B}^{1} .
\end{aligned}
$$

These families of operators are not only bounded, but have $L_{2}$ off-diagonal bounds in the following sense.

Definition 2.24. Let $\left(U_{t}\right)_{t>0}$ be a family of operators on $\mathcal{H}$, and let $M \geqslant 0$. We say that $\left(U_{t}\right)_{t>0}$ has $L_{2}$ off-diagonal bounds (with exponent $M$ ) if there exists $C_{M}<\infty$ such that

$$
\left\|U_{t} f\right\|_{L_{2}(E)} \leqslant C_{M}\langle\operatorname{dist}(E, F) / t\rangle^{-M}\|f\|
$$

whenever $E, F \subset \mathbf{R}^{n}$ and supp $f \subset F$. Here $\langle x\rangle:=1+|x|$, and $\operatorname{dist}(E, F):=\inf \{|x-y|$ : $x \in E, y \in F\}$. We write $\left\|U_{t}\right\|_{\text {off }, M}$ for the smallest constant $C_{M}$. The exact value of $M$ is normally not important and we write $\left\|U_{t}\right\|_{\text {off }}$, where it is understood that $M$ is chosen sufficiently large but fixed.

Proposition 2.25. All the operator families from Corollary 2.23 has $L_{2}$ off-diagonal bounds for all exponents $M \geqslant 0$.

Proof. First consider the resolvents $R_{t}^{B}=\left(I+i t T_{B}\right)^{-1}$. As we already have proved uniform bounds for $R_{t}^{B}$, it suffices to prove

$$
\left\|\left(I+i t T_{B}\right)^{-1} f\right\|_{L_{2}(E)} \leqslant C_{M}(|t| / \operatorname{dist}(E, F))^{M}\|f\|
$$

for $|t| \leqslant \operatorname{dist}(E, F)$. We prove this by induction on $M$ as in [10, Proposition 5.2]. Let $\eta: \mathbf{R}^{n} \rightarrow$ $[0,1]$ be a bump function such that $\left.\eta\right|_{E}=1, \operatorname{supp} \eta \subset \widetilde{E}:=\left\{x \in \mathbf{R}^{n} ; \operatorname{dist}(x, E) \leqslant \operatorname{dist}(x, F)\right\}$ and $\|\nabla \eta\|_{\infty} \lesssim 1 / \operatorname{dist}(E, F) \approx 1 / \operatorname{dist}(\widetilde{E}, F)$. Since the commutator is

$$
\left[\eta I, R_{t}^{B}\right]=t R_{t}^{B} M_{B}^{-1}\left([\eta I, \underline{d}]+B^{-1}\left[\eta I, \underline{d}^{*}\right] B\right) R_{t}^{B},
$$

where $\|[\eta I, \underline{d}]\|,\left\|\left[\eta I, \underline{d}^{*}\right]\right\| \lesssim\|\nabla \eta\|_{\infty}$, we get

$$
\left\|R_{t}^{B} f\right\|_{L_{2}(E)} \leqslant\left\|\eta R_{t}^{B} f\right\|=\left\|\left[\eta I, R_{t}^{B}\right] f\right\| \lesssim|t|\|\nabla \eta\|_{\infty}\left\|R_{t}^{B} f\right\|_{L_{2}(\tilde{E})},
$$

where we used that $\eta f=0$. By induction, this proves the off-diagonal bounds for $R_{t}^{B}$. From this, off-diagonal bounds for $P_{t}^{B}, Q_{t}^{B}$ and $I-P_{t}^{B}$ also follows immediately.

Next we consider $t \underline{d} P_{t}^{B}$ and use Lemma 2.21(i) to obtain

$$
\begin{align*}
\left\|t \underline{d} P_{t}^{B} f\right\|_{L_{2}(E)} & \leqslant\left\|\eta t \underline{d} P_{t}^{B} f\right\| \leqslant\left\|[\eta I, t \underline{d}] P_{t}^{B} f\right\|+\left\|t \underline{d} \eta P_{t}^{B} f\right\| \\
& \lesssim|t|\|\nabla \eta\|_{\infty}\left\|P_{t}^{B} f\right\|_{L_{2}(\widetilde{E})}+\left\|t T_{B} \eta P_{t}^{B} f\right\| \\
& \lesssim|t|\|\nabla \eta\|_{\infty}\left\|P_{t}^{B} f\right\|_{L_{2}(\widetilde{E})}+\left\|\left[\eta, t T_{B}\right] P_{t}^{B} f\right\|+\left\|\eta t T_{B} P_{t}^{B} f\right\| \\
& \lesssim|t|\|\nabla \eta\|_{\infty}\left\|P_{t}^{B} f\right\|_{L_{2}(\widetilde{E})}+\left\|Q_{t}^{B} f\right\|_{L_{2}(\widetilde{E})} . \tag{2.4}
\end{align*}
$$

This and the corresponding calculation with $\underline{d}$ replaced by $\underline{d}_{B}^{*}$ proves the off-diagonal bounds for $t \underline{d} P_{t}^{B}$ and $t \underline{d}_{B}^{*} P_{t}^{B}$. From this the result for $P_{t}^{B} M_{B}^{-1} t \underline{d}$ and $P_{t}^{B} M_{B}^{-1} t \underline{d}_{B}^{*}$ follows immediately with a duality argument. Indeed, $\left(M_{B}^{-1} \underline{d}\right)^{\prime}=M_{B^{*}}^{-1} \underline{d}_{B^{*}}^{*}$ and $\left(M_{B}^{-1} \underline{d}_{B}^{*}\right)^{\prime}=M_{B^{*}}^{-1} \underline{d}$ is proved similarly to $T_{B}^{\prime}=-T_{B^{*}}$ in Proposition 2.17.

The proof for $t \underline{d} Q_{t}^{B}, t \underline{d}_{B}^{*} Q_{t}^{B}, Q_{t}^{B} M_{B}^{-1} t \underline{d}$ and $Q_{t}^{B} M_{B}^{-1} t \underline{d}_{B}^{*}$ is similar, replacing $P_{t}$ and $Q_{t}$ with $Q_{t}$ and $I-P_{t}$. Finally, the last four estimates follows from a computation like (2.4), for example replacing $P_{t}$ and $Q_{t}$ with $P_{t}^{B} M_{B}^{-1} t \underline{d}$ and $Q_{t}^{B} M_{B}^{-1} t \underline{d}$ proves the estimate for $t \underline{d} P_{t}^{B} M_{B}^{-1} t \underline{d}$.

We finish this section with a lemma to be used in Section 4. This lemma is proved with an argument similar to that in [15, Lemma 2.3]. For completeness, we include a short proof.

Lemma 2.26. Assume that $\left(U_{t}\right)_{t>0}$ and $\left(V_{t}\right)_{t>0}$ both have $L_{2}$ off-diagonal bounds with exponent $M$. Then $\left(U_{t} V_{t}\right)_{t>0}$ has $L_{2}$ off-diagonal bounds with exponent $M$ and

$$
\left\|U_{t} V_{t}\right\|_{\text {off }, M} \leqslant 2^{M+1}\left\|U_{t}\right\|_{\text {off }, M}\left\|V_{t}\right\|_{\text {off }, M}
$$

Proof. By Definition 2.24 we need to prove that

$$
\left\|U_{t} V_{t} f\right\|_{L_{2}(E)} \lesssim\langle\operatorname{dist}(E, F) / t\rangle^{-M}\|f\|
$$

whenever supp $f \subset F$. To this end let $G:=\left\{x \in \mathbf{R}^{n} ; \operatorname{dist}(x, E) \leqslant \rho / 2\right\}$, where $\rho:=\operatorname{dist}(E, F)$. We get

$$
\begin{aligned}
\left\|U_{t} V_{t} f\right\|_{L_{2}(E)} & \leqslant\left\|U_{t}\left(\chi_{G} V_{t} f\right)\right\|_{L_{2}(E)}+\left\|U_{t}\left(\chi_{\mathbf{R}^{n} \backslash G} V_{t} f\right)\right\|_{L_{2}(E)} \\
& \leqslant C_{0}^{U_{t}}\left\|V_{t} f\right\|_{L_{2}(G)}+C_{M}^{U_{t}}\langle\rho / 2 t\rangle^{-M}\left\|V_{t} f\right\|_{L_{2}\left(\mathbf{R}^{n} \backslash G\right)} \\
& \leqslant C_{0}^{U_{t}} C_{M}^{V_{t}}\langle\rho / 2 t\rangle^{-M}\|f\|+C_{M}^{U_{t}}\langle\rho / 2 t\rangle^{-M} C_{0}^{V_{t}}\|f\| \leqslant 2^{M+1} C_{M}^{U_{t}} C_{M}^{V_{t}}\langle\rho / t\rangle^{-M}\|f\| .
\end{aligned}
$$

### 2.3. Quadratic estimates: generalities

In Proposition 2.22 we proved the spectral estimate $\sigma\left(T_{B}\right) \subset S_{\omega}$ for some angle $\omega<\pi / 2$ with bounds on the resolvent outside $S_{\omega}$. In this section we survey some general facts about the functional calculus of the operator $T_{B}$. For a further background and discussion of these matters we refer to $[1,10]$.

Definition 2.27. For $\omega<\nu<\pi / 2$, we define the following classes of holomorphic functions $f \in H\left(S_{v}^{o}\right)$ on the open double sector $S_{v}^{o}$ :

$$
\begin{aligned}
\Psi\left(S_{v}^{o}\right) & :=\left\{\psi \in H\left(S_{v}^{o}\right) ;|\psi(z)| \leqslant C \min \left(|z|^{s},|z|^{-s}\right), z \in S_{v}^{o}, \text { for some } s>0, C<\infty\right\} \\
H_{\infty}\left(S_{v}^{o}\right) & :=\left\{b \in H\left(S_{v}^{o}\right) ;|b(z)| \leqslant C, z \in S_{v}^{o}, \text { for some } C<\infty\right\} \\
F\left(S_{v}^{o}\right) & :=\left\{w \in H\left(S_{v}^{o}\right) ;|w(z)| \leqslant C \max \left(|z|^{s},|z|^{-s}\right), z \in S_{v}^{o}, \text { for some } s<\infty, C<\infty\right\} .
\end{aligned}
$$

Thus $\Psi\left(S_{v}^{o}\right) \subset H_{\infty}\left(S_{v}^{o}\right) \subset F\left(S_{v}^{o}\right) \subset H\left(S_{v}^{o}\right)$.
For $\psi \in \Psi\left(S_{v}^{o}\right)$, we define a bounded operator $\psi\left(T_{B}\right)$ through the Dunford functional calculus

$$
\begin{equation*}
\psi\left(T_{B}\right):=\frac{1}{2 \pi i} \int_{\gamma} \psi(\lambda)\left(\lambda I-T_{B}\right)^{-1} d \lambda \tag{2.5}
\end{equation*}
$$

where $\gamma$ is the unbounded contour $\left\{ \pm r e^{ \pm i \theta} ; r>0\right\}, \omega<\theta<\nu$, parametrised counterclockwise around $S_{\omega}$. The decay estimate on $\psi$ and the resolvent bounds of Proposition 2.22 guarantee that $\left\|\psi\left(T_{B}\right)\right\|<\infty$.

For general $w \in F\left(S_{v}^{o}\right)$ we define

$$
w\left(T_{B}\right):=\left(Q^{B}\right)^{-k}\left(q^{k} w\right)\left(T_{B}\right),
$$

where $k$ is an integer larger than $s$ if $|w(z)| \leqslant C \max \left(|z|^{s},|z|^{-s}\right)$, and $q(z):=z\left(1+z^{2}\right)^{-1}$ and $Q^{B}:=q\left(T_{B}\right)$. This yields a closed, densely defined operator $w\left(T_{B}\right)$ in $\mathcal{H}$. Furthermore, we have

$$
\begin{align*}
\overline{\lambda_{1} w_{1}\left(T_{B}\right)+\lambda_{2} w_{2}\left(T_{B}\right)} & =\left(\lambda_{1} w_{1}+\lambda_{2} w_{2}\right)\left(T_{B}\right), \\
\overline{w_{1}\left(T_{B}\right) w_{2}\left(T_{B}\right)} & =\left(w_{1} w_{2}\right)\left(T_{B}\right), \tag{2.6}
\end{align*}
$$

for all $w_{1}$ and $w_{2} \in F\left(S_{v}^{o}\right)$. Here $\bar{T}=S$ means that the graph $\mathrm{G}(T)$ is dense in the graph $\mathrm{G}(S)$.

The functional calculus $w \mapsto w\left(T_{B}\right)$ have the following convergence properties as proved in [1].

Lemma 2.28. If $b_{k} \in H_{\infty}\left(S_{v}^{o}\right)$ is a sequence uniformly bounded on $S_{v}^{o}$ which converges to $b$ uniformly on compact subsets, and if $b_{k}\left(T_{B}\right)$ are uniformly bounded operators, then

$$
b_{k}\left(T_{B}\right) f \rightarrow b\left(T_{B}\right) f, \quad \text { for all } f \in \mathcal{H},
$$

and $\left\|b\left(T_{B}\right)\right\| \leqslant \lim \sup _{k}\left\|b_{k}\left(T_{B}\right)\right\|$.
Definition 2.29. The following operators in the functional calculus are of special importance to us.
(1) $q_{t}(z)=q(t z):=t z\left(1+t^{2} z^{2}\right)^{-1} \in \Psi\left(S_{v}^{o}\right)$, which give the operator $Q_{t}^{B}$.
(2) $|z|^{s}:=\left(z^{2}\right)^{s / 2} \in F\left(S_{v}^{o}\right)$, which give the operator $\left|T_{B}\right|^{s}$. Note that $|z|$ does not denote absolute value here, but $z \mapsto|z|$ is holomorphic on $S_{v}^{o}$.
(3) $e^{-t|z|} \in H_{\infty}\left(S_{v}^{o}\right)$, which give the operator $e^{-t\left|T_{B}\right|}$.
(4) The characteristic functions

$$
\chi^{ \pm}(z)= \begin{cases}1 & \text { if } \pm \operatorname{Re} z>0 \\ 0 & \text { if } \pm \operatorname{Re} z<0\end{cases}
$$

which give the generalised Hardy projections $E_{B}^{ \pm}:=\chi^{ \pm}\left(T_{B}\right)$.
(5) The signum function

$$
\operatorname{sgn}(z)=\chi^{+}(z)-\chi^{-}(z)
$$

which give the generalised Cauchy integral $E_{B}:=\operatorname{sgn}\left(T_{B}\right)$.
The main work in this paper is to prove the boundedness the projections $E_{B}^{ \pm}$. As in Lemma 2.14, if these are bounded then they correspond to a splitting

$$
\mathcal{H}=E_{B}^{+} \mathcal{H} \oplus E_{B}^{-} \mathcal{H}
$$

of $\mathcal{H}$ into the Hardy subspaces $E_{B}^{ \pm} \mathcal{H}$ associated with Eq. (2.1). That the projections are bounded is also equivalent with having a bounded reflection operator $E_{B}$.

Definition 2.30. For a function $F(t, x)$ defined in $\mathbf{R}_{ \pm}^{n+1}$ we write

$$
\|F\|_{ \pm}:=\left(\int_{0}^{\infty}\|F( \pm t, x)\|^{2} \frac{d t}{t}\right)^{1 / 2}
$$

and for short $\|\mid F\|_{+}=:\|F\| \|$. When $F(t, x)=\left(\Theta_{t} f\right)(x)$ for some family of operators $\left(\Theta_{t}\right)_{t>0}$, we use the notation

$$
\left\|\Theta_{t}\right\|_{\text {op }}:=\sup _{\|f\|=1}\left\|\Theta_{t} f\right\| .
$$

Our main goal in this paper will be to prove quadratic estimates of the form

$$
\begin{equation*}
\int_{0}^{\infty}\left\|Q_{t}^{B} f\right\|^{2} \frac{d t}{t} \approx\|f\|^{2} \tag{2.7}
\end{equation*}
$$

for certain coefficients $B$. We recall the following two basic results concerning quadratic estimates which are proved by Schur estimates. For details we refer to [1].

Proposition 2.31. Let $\psi \in \Psi\left(S_{v}^{o}\right)$ be non-vanishing on both $S_{v+}^{o}$ and $S_{v-}^{o}$, and define $\psi_{t}(z):=$ $\psi(t z)$. Then there exists $0<C<\infty$ such that

$$
C^{-1}\| \| Q_{t}^{B} f\| \| \leqslant\left\|\psi_{t}\left(T_{B}\right) f\right\|\|\leqslant C\| Q_{t}^{B} f\| \|
$$

Proposition 2.32. If $T_{B}$ satisfies quadratic estimates (2.7), then $T_{B}$ has bounded $H_{\infty}\left(S_{v}^{o}\right)$ functional calculus, i.e.

$$
\left\|b\left(T_{B}\right)\right\| \lesssim\|b\|_{\infty}, \quad \text { for all } b \in H_{\infty}\left(S_{v}^{o}\right) .
$$

Thus $H_{\infty}\left(S_{v}^{o}\right) \ni b \mapsto b\left(T_{B}\right) \in \mathcal{L}(\mathcal{H})$ is a continuous homomorphism.
Before proving quadratic estimates for $T_{B}$ for certain $B$ in Sections 3 and 4, we introduce a dense subspace on which the operator $b\left(T_{B}\right)$ is defined for any $b \in H_{\infty}\left(S_{v}^{o}\right)$.

Definition 2.33. Let $\mathcal{V}_{B}$ be the dense linear subspace

$$
\mathcal{V}_{B}:=\bigcup_{s>0}\left(\mathrm{D}\left(\left|T_{B}\right|^{s}\right) \cap \mathrm{R}\left(\left|T_{B}\right|^{s}\right)\right) \subset \mathcal{H} .
$$

We see that $\mathrm{D}\left(\left|T_{B}\right|^{s}\right) \cap \mathrm{R}\left(\left|T_{B}\right|^{s}\right)$ increases when $s$ decreases. The density of $\mathrm{D}\left(\left|T_{B}\right|\right) \cap \mathrm{R}\left(\left|T_{B}\right|\right)$, and therefore of $\mathcal{V}_{B}$, follows from the fact that

$$
2 \int_{\alpha}^{\beta}\left(Q_{t}^{B}\right)^{2} f \frac{d t}{t}=\left(P_{\alpha}^{B}-P_{\beta}^{B}\right) f \rightarrow f
$$

as $(\alpha, \beta) \rightarrow(0, \infty)$, for all $f \in \mathcal{H}$.
Moreover, if $b \in H_{\infty}\left(S_{v}^{o}\right)$ and $f \in \mathcal{V}_{B}$ then $b\left(T_{B}\right) f \in \mathcal{V}_{B} \subset \mathcal{H}$. To see this, write

$$
b\left(T_{B}\right) f=(b \psi)\left(T_{B}\right)\left(\psi\left(T_{B}\right)^{-1} f\right)
$$

where $\psi(z)^{-1}:=\left(1+|z|^{s}\right) /|z|^{s / 2}$ if $f \in \mathrm{D}\left(\left|T_{B}\right|^{s}\right) \cap \mathrm{R}\left(\left|T_{B}\right|^{s}\right)$. Then $\psi\left(T_{B}\right)^{-1} f \in \mathcal{H}$ and $(b \psi)\left(T_{B}\right)$ is bounded since $b \psi \in \Psi\left(S_{v}^{o}\right)$. Furthermore, if $s^{\prime}<s / 2$ then $\left|T_{B}\right|^{s^{\prime}}\left(\psi\left(T_{B}\right)^{-1} f\right) \in \mathcal{H}$ and $\left|T_{B}\right|^{-s^{\prime}}\left(\psi\left(T_{B}\right)^{-1} f\right) \in \mathcal{H}$, so $b\left(T_{B}\right) f \in \mathrm{D}\left(\left|T_{B}\right|^{s^{\prime}}\right) \cap \mathrm{R}\left(\left|T_{B}\right|^{s^{\prime}}\right)$.

Lemma 2.34. We have an algebraic splitting

$$
\mathcal{V}_{B}=E_{B}^{+} \mathcal{V}_{B}+E_{B}^{-} \mathcal{V}_{B}
$$

$E_{B}^{+} \mathcal{V}_{B} \cap E_{B}^{-} \mathcal{V}_{B}=\{0\}$, and $f \in E_{B}^{ \pm} \mathcal{V}_{B}$ is in one-to-one correspondence with $F(t, x)=F_{t}(x):=$ $\left(e^{\mp t\left|T_{B}\right|} f\right)(x)$ in $\mathbf{R}_{ \pm}^{n+1}$, and

$$
\begin{gathered}
\lim _{t \rightarrow 0^{ \pm}} F_{t}=f, \quad \lim _{t \rightarrow \pm \infty} F_{t}=0, \\
F_{t} \in \mathrm{D}\left(T_{B}\right)=\mathrm{D}(d) \cap \mathrm{D}\left(d^{*} B\right), \quad \partial_{t} F_{t}=-T_{B} F_{t} \in L_{2}\left(\mathbf{R}^{n}\right), \quad \pm t>0 .
\end{gathered}
$$

Proof. That each $f \in \mathcal{V}_{B}$ can be uniquely written $f=f^{+}+f^{-}$, where $f^{ \pm} \in E_{B}^{ \pm} \mathcal{V}_{B}$, follows from (2.6). To verify the properties of $F_{t}$, it suffices to consider the case $f \in E_{B}^{+} \mathcal{V}_{B}$ as the case $f \in E_{B}^{-} \mathcal{V}_{B}$ is similar. Since $z e^{-t|z|} \in \Psi\left(S_{v}^{0}\right)$ it follows that $F_{t} \in \mathrm{D}\left(T_{B}\right)$. Moreover, since $\frac{1}{h}\left(e^{-(t+h)|z|}-e^{-t|z|}\right) \rightarrow-|z| e^{-t|z|}$ uniformly on $S_{v}^{0}$, it follows from Lemma 2.28 that $\partial_{t} F_{t}=$ $-\left|T_{B}\right| F_{t}$ and since $F_{t} \in E_{B}^{+} \mathcal{V}_{B}$ we have $\left|T_{B}\right| F_{t}=T_{B} F_{t}$.

To prove the limits, assume that $f \in \mathrm{D}\left(\left|T_{B}\right|^{s}\right) \cap \mathrm{R}\left(\left|T_{B}\right|^{s}\right)$ for some $0<s<1$. Writing $f=$ $\left|T_{B}\right|^{-s} u$, we see that

$$
F_{t}-f=t^{s} \psi\left(t T_{B}\right) u, \quad \text { where } \psi(z)=\left(e^{-|z|}-1\right) /|z|^{S} .
$$

Similarly with $f=\left|T_{B}\right|^{s} v$, we get

$$
F_{t}=t^{-s} \psi\left(t T_{B}\right) v, \quad \text { where } \psi(z)=|z|^{s} e^{-|z|}
$$

Since in both cases $\psi\left(t T_{B}\right)$ are uniformly bounded in $t$, using a direct norm estimate in (2.5), it follows that both limits are 0 as $t \rightarrow 0$ and $t \rightarrow \infty$, respectively. In particular, $f=\lim _{t \rightarrow 0} F_{t}$ is uniquely determined by $F$.

We now further discuss the quadratic estimates (2.7). First note the following consequence of the duality $T_{B}^{\prime}=-T_{B^{*}}$ from Proposition 2.17. Again we refer to [1] for further details.

Lemma 2.35. If $\left\|Q_{t}^{B} f\right\| \lesssim\|f\|$ for all $f \in \mathcal{H}$, then $\left\|Q_{t}^{B^{*}} f\right\| \gtrsim\|f\|$ for all $f \in \mathcal{H}$.
In Section 3.3 we shall use the following Hardy space reduction of the quadratic estimate. This is a technique due to Coifman, Jones and Semmes [11], and adapted to the setting of functional calculus by $\mathrm{M}^{\mathrm{c}}$ Intosh and Qian [23, Theorem 5.2].

Proposition 2.36. Assume that we have reverse quadratic estimates in $E_{B^{*}}^{+} \mathcal{V}_{B^{*}}$ and $E_{B^{*}}^{-} \mathcal{V}_{B^{*}}$, i.e.

$$
\|g\| \lesssim\left\|t t \partial_{t} G_{t}\right\|_{ \pm}, \quad g \in E_{B^{*}}^{ \pm} \mathcal{V}_{B^{*}},
$$

where $G_{t}=e^{\mp t\left|T_{B^{*}}\right|} g$. Then

$$
\left\|Q_{t}^{B} f\right\| \lesssim\|f\|, \quad f \in \mathcal{H}
$$

In particular, if we have reverse quadratic estimates in both Hardy spaces for both operators $T_{B}$ and $T_{B^{*}}$, then $\left\|Q_{t}^{B} f\right\| \approx\|f\| \approx\left\|Q_{t}^{B^{*}} f\right\|, f \in \mathcal{H}$, and thus $T_{B}$ and $T_{B^{*}}$ have bounded $H_{\infty}\left(S_{v}^{o}\right)$ functional calculus for all $\omega<\nu<\pi / 2$.

Proof. Assume first that

$$
\|g\| \lesssim\left\|t \partial_{t} G_{t}\right\|, \quad g \in E_{B^{*}}^{+} \mathcal{V}_{B^{*}}
$$

and write $\psi_{t}(z):=t z e^{-t|z|}$ so that $t \partial_{t} G_{t}=-\psi_{t}\left(T_{B^{*}}\right) g$ by Lemma 2.34. Let $f \in \mathcal{H}$ and define $q_{t}^{\mp}(z):=\chi^{\mp} q_{t}(z)$ so that $q_{t}^{\mp}\left(T_{B}\right)=E_{B}^{\mp} Q_{t}^{B}$. To prove that

$$
\left\|q_{t}^{-}\left(T_{B}\right) f\right\| \lesssim \lesssim f \|, \quad f \in \mathcal{H}
$$

it suffices to bound $\int_{\alpha}^{\beta}\left\|q_{t}^{-}\left(T_{B}\right) f\right\|^{2} \frac{d t}{t}$ uniformly for $\alpha>0$ and $\beta<\infty$. To this end, define the auxiliary functions

$$
\begin{aligned}
h_{t} & :=\left(\int_{\alpha}^{\beta}\left\|q_{t}^{-}\left(T_{B}\right) f\right\|^{2} \frac{d t}{t}\right)^{-1 / 2}\left(N^{+} B^{*}-B^{*} N^{-}\right)^{-1} q_{t}^{-}\left(T_{B}\right) f, \\
g & :=-\int_{\alpha}^{\beta} q_{t}^{+}\left(T_{B^{*}}\right) h_{t} \frac{d t}{t}
\end{aligned}
$$

so that $\int_{\alpha}^{\beta}\left\|h_{t}\right\|^{2} \frac{d t}{t} \leqslant C$ and $g \in E_{B^{*}}^{+} \mathcal{V}_{B^{*}}$, and calculate

$$
\begin{aligned}
\left(\int_{\alpha}^{\beta}\left\|q_{t}^{-}\left(T_{B}\right) f\right\|^{2} \frac{d t}{t}\right)^{1 / 2} & =\int_{\alpha}^{\beta}\left\langle q_{t}^{-}\left(T_{B}\right) f, h_{t}\right\rangle_{B} \frac{d t}{t}=\langle f, g\rangle_{B} \\
& \lesssim\|f\|\|g\| \lesssim\|f\|\left\|\left\|\psi_{s}\left(T_{B^{*}}\right) g\right\|\right\| \\
& \lesssim\|f\|\left(\int_{0}^{\infty}\left(\int_{\alpha}^{\beta}\left\|\left(\psi_{s} q_{t}^{+}\right)\left(T_{B^{*}}\right)\right\|\left\|h_{t}\right\| \frac{d t}{t}\right)^{2} \frac{d s}{s}\right)^{1 / 2} \\
& \lesssim\|f\|\left(\int_{\alpha}^{\beta}\left\|h_{t}\right\|^{2} \frac{d t}{t}\right)^{1 / 2} \lesssim\|f\|
\end{aligned}
$$

In the second equality we have used that $q_{t}^{-}\left(T_{B}\right)^{\prime}=q_{t}^{-}\left(-T_{B^{*}}\right)=-q_{t}^{+}\left(T_{B^{*}}\right)$, in the second estimate we used the hypothesis and the second last estimate is a Schur estimate. We here use that $\left\|\left(\psi_{s} q_{t}^{+}\right)\left(T_{B^{*}}\right)\right\| \lesssim \eta(t / s)$, where $\eta(x):=\min \left(x^{s}, x^{-s}\right)$ for some $s>0$.

With a similar argument $\left\|q_{t}^{+}\left(T_{B}\right) f\right\| \lesssim\|f\|, f \in \mathcal{H}$, follows from the reverse quadratic estimate for $g \in E_{B^{*}}^{-} \mathcal{V}_{B^{*}}$. If both reverse estimates holds for $B^{*}$, then

$$
\left\|\left\|Q_{t}^{B} f\right\|\right\| \leqslant\left\|q_{t}^{-}\left(T_{B}\right) f\right\|\|+\| q_{t}^{+}\left(T_{B}\right) f\| \| \lesssim\|f\|, \quad f \in \mathcal{H}
$$

and if the same holds for $B$ and $B^{*}$ interchanged, then Lemma 2.35 proves that $\left\|Q_{t}^{B} f\right\| \approx\|f\|$, $f \in \mathcal{H}$, and by Proposition 2.32 this proves that $T_{B}$ has bounded $H_{\infty}\left(S_{v}^{o}\right)$ functional calculus, and similarly for $T_{B^{*}}$.

We end this section with a discussion of the holomorphic perturbation theory for the functional calculus of $T_{B}$.

Definition 2.37. Let $z \mapsto U_{z} \in \mathcal{L}(\mathcal{X}, \mathcal{Y})$ be an operator-valued function defined on an open subset $D \subset \mathbf{C}$ of the complex plane. We say that $U_{z}$ is holomorphic if for all $z \in D$ there exists an operator $U_{z}^{\prime} \in \mathcal{L}(\mathcal{X}, \mathcal{Y})$ such that

$$
\left\|\frac{1}{w}\left(U_{z+w}-U_{z}\right)-U_{z}^{\prime}\right\|_{\mathcal{X} \rightarrow \mathcal{Y}} \rightarrow 0, \quad w \rightarrow 0
$$

Lemma 2.38. Let $z \mapsto U_{z} \in \mathcal{L}(\mathcal{X}, \mathcal{Y})$ be an operator-valued function defined on an open subset $D \subset \mathbf{C}$ of the complex plane. Then the following are equivalent.
(i) $z \mapsto U_{z}$ is holomorphic.
(ii) The scalar function $h(z)=\left(U_{z} f, g\right)$ is a holomorphic function for all $f \in \tilde{\mathcal{X}}$ and $g \in \tilde{\mathcal{Y}}^{*}$, where $\tilde{\mathcal{X}} \subset \mathcal{X}$ and $\tilde{\mathcal{Y}}^{*} \subset \mathcal{Y}^{*}$ are dense, and $\left\|U_{z}\right\|$ is locally bounded.

In particular, if $z \mapsto U_{z}^{k}$ are holomorphic on $D$ for $k=1,2, \ldots$ and

$$
\left(U_{z}^{k} f, g\right)=h_{k}(z) \rightarrow h(z)=\left(U_{z} f, g\right), \quad \text { for all } z \in D, f \in \tilde{\mathcal{X}}, g \in \tilde{\mathcal{Y}}^{*}
$$

and $\sup _{z \in K, k \geqslant 1}\left\|U_{z}^{k}\right\|<\infty$ for each compact subset $K \subset D$, then $z \mapsto U_{z}$ is holomorphic on $D$.
Proof. For the equivalence between (i) and (ii), see Kato [18, Theorem III 3.12]. To prove the convergence result, it suffices to show that $h(z)$ is holomorphic on $D$. That this is true follows from an application of the dominated convergence theorem in the Cauchy integral formula for $h_{k}(z)$.

Below, we shall assume that $z \mapsto B_{z}$ is a given holomorphic matrix-valued function defined on an open subset $D \subset \mathbf{C}$ such that $B_{z}$ is a multiplication operator as in Definition 2.8 for each $z \in D$, and that $\omega_{D}:=\sup _{z \in D} \arccos \left(\kappa_{B_{z}} /\left(2\left\|B_{z}\right\|_{\infty}\right)\right)<\pi / 2$. Let $\omega_{D}<\nu<\pi / 2$.

Lemma 2.39. For $\tau \in S_{\pi / 2-v}^{0}$, the operator-valued function $D \ni z \mapsto\left(I-i \tau T_{B_{z}}\right)^{-1}$ is holomorphic.

Proof. Similar to the proof of Proposition 2.22 we have

$$
\left(I-i \tau T_{B_{z}}\right)^{-1}=\left(\widetilde{B}_{z} \Gamma+\widetilde{\Gamma} B_{z}\right)^{-1} i\left(\widetilde{B}_{z} \mu-\mu^{*} B_{z}\right)
$$

where $\Gamma=i \mu+\tau d$ and $\widetilde{\Gamma}=-i \mu+\tau d^{*}$. It is clear that the multiplication operator $\widetilde{B}_{z} \mu-\mu^{*} B_{z}$ depends holomorphically on $z$, so it suffices to show that $z \mapsto\left(\widetilde{B}_{z} \Gamma+\widetilde{\Gamma} B_{z}\right)^{-1}$ is holomorphic. Let $z \in D$, write $B:=B_{z}$ and $B_{w}:=B_{z+w}$ and calculate

$$
\begin{align*}
& \left(\widetilde{B}_{w} \Gamma+\widetilde{\Gamma} B_{w}\right)^{-1}-(\widetilde{B} \Gamma+\widetilde{\Gamma} B)^{-1} \\
& =-\left(\widetilde{B}_{w} \Gamma+\widetilde{\Gamma} B_{w}\right)^{-1}\left(\widetilde{B}_{w}-\widetilde{B}\right)\left(\Gamma(\widetilde{B} \Gamma+\widetilde{\Gamma} B)^{-1}\right) \\
& \quad-\left(\left(\widetilde{B}_{w} \Gamma+\widetilde{\Gamma} B_{w}\right)^{-1} \widetilde{\Gamma}\right)\left(B_{w}-B\right)(\widetilde{B} \Gamma+\widetilde{\Gamma} B)^{-1} . \tag{2.8}
\end{align*}
$$

The first and last factors in both terms on the right-hand side are uniformly bounded by Lemma 2.21(i) so we deduce continuity of $w \mapsto\left(\widetilde{B}_{w} \Gamma+\widetilde{\Gamma} B_{w}\right)^{-1}$. Furthermore, multiplying Eq. (2.8) from the right with $\widetilde{\Gamma}$, the first term on the right vanishes as

$$
\Gamma(\widetilde{B} \Gamma+\widetilde{\Gamma} B)^{-1} \widetilde{\Gamma}=(\widetilde{B} \Gamma+\widetilde{\Gamma} B)^{-1} \widetilde{\Gamma}^{2}=0
$$

and we deduce continuity of $w \mapsto\left(\widetilde{B}_{w} \Gamma+\widetilde{\Gamma} B_{w}\right)^{-1} \widetilde{\Gamma}$. Thus, dividing Eq. (2.8) by $w$ and letting $w \rightarrow 0$ we see that the limit exists and equals

$$
-(\widetilde{B} \Gamma+\widetilde{\Gamma} B)^{-1} \widetilde{B}^{\prime}\left(\Gamma(\widetilde{B} \Gamma+\widetilde{\Gamma} B)^{-1}\right)-\left((\widetilde{B} \Gamma+\widetilde{\Gamma} B)^{-1} \widetilde{\Gamma}\right) B^{\prime}(\widetilde{B} \Gamma+\widetilde{\Gamma} B)^{-1}
$$

Lemma 2.40. If $\psi \in \Psi\left(S_{v}^{o}\right)$, then $D \ni z \mapsto \psi\left(T_{B_{z}}\right)$ is holomorphic.
Proof. Let $\gamma$ be the unbounded contour $\left\{ \pm r e^{ \pm i \theta} ; r>0\right\}, \omega_{D}<\theta<\nu$, parametrised counterclockwise around $S_{\omega_{D}}$. By inspection of the proof of Lemma 2.39 we have

$$
\sup _{\lambda \in \gamma}|\lambda|\left\|\frac{1}{w}\left(\left(\lambda-T_{B_{z+w}}\right)^{-1}-\left(\lambda-T_{B_{z}}\right)^{-1}\right)-\partial_{z}\left(\lambda-T_{B_{z}}\right)^{-1}\right\| \rightarrow 0,
$$

as $w \rightarrow 0$. Thus

$$
\frac{1}{w}\left(\psi\left(T_{B_{z+w}}\right)-\psi\left(T_{B_{z}}\right)\right) \rightarrow \frac{1}{2 \pi i} \int_{\gamma} \psi(\lambda) \partial_{z}\left(\lambda-T_{B_{z}}\right)^{-1} d \lambda
$$

since $\left.\int_{\gamma}|\psi(\lambda)| \frac{d \lambda}{\lambda} \right\rvert\,<\infty$.
Lemma 2.41. Assume that $T_{B_{z}}$ satisfy quadratic estimates $\left\|\left\|Q_{t}^{B_{z}} f\right\|\right\|\|f\|, f \in \mathcal{H}$, locally uniformly for $z \in D$. If $b \in H_{\infty}\left(S_{v}^{o}\right), \psi \in \Psi\left(S_{v}^{o}\right)$ and $0<\alpha<\beta<\infty$, then the following operators depend holomorphically on $z \in D$.
(i) $u(x) \mapsto\left(b\left(T_{B_{z}}\right) u\right)(x): \mathcal{H} \rightarrow \mathcal{H}$;
(ii) $u(x) \mapsto v(t, x)=\left(b\left(t T_{B_{z}}\right) u\right)(x): \mathcal{H} \rightarrow L_{2}\left(\mathbf{R}^{n} \times(\alpha, \beta)\right.$; $\left.\wedge\right)$;
(iii) $u(x) \mapsto v(t, x)=\left(\psi_{t}\left(T_{B_{z}}\right) u\right)(x): \mathcal{H} \rightarrow L_{2}\left(\mathbf{R}_{+}^{n+1}, \frac{d t d x}{t} ; \bigwedge\right)$.

Proof. (i) Take a uniformly bounded sequence $\psi_{k}(z) \in \Psi\left(S_{v}^{o}\right)$ which converges to $b(z) \in$ $H_{\infty}\left(S_{v}^{o}\right)$ uniformly on compact subsets of $S_{v}^{o}$. Lemma 2.38 then applies with $U_{z}^{k}=\psi_{k}\left(T_{B_{z}}\right)$ and $U_{z}=b\left(T_{B_{z}}\right)$, using Lemmas 2.40 and 2.28.
(ii) It suffices by Lemma 2.38 to show that $h(z)=\int_{\alpha}^{\beta} h_{t}(z) d t$ is holomorphic, where $h_{t}(z)=$ $\left(b\left(t T_{B_{z}}\right) f, G_{t}\right)$, for all $f(x) \in \mathcal{H}$ and $G(t, x) \in C_{0}^{\infty}\left(\mathbf{R}^{n} \times(\alpha, \beta) ; \wedge\right)$. That $h_{t}(z)$ is holomorphic for each $t$ is clear from (i), and for $h(z)$ this follows from an application of the Fubini theorem to the Cauchy integral formula for $h_{t}(z)$.
(iii) Consider the truncations $U_{z}^{k}: u(x) \mapsto v(t, x)=\chi_{k}(t)\left(\psi_{t}\left(T_{B_{z}}\right) u\right)(x)$, where $\chi_{k}$ denotes the characteristic function of the interval $(1 / k, k)$. It is clear from (ii) that $U_{z}^{k}$ is holomorphic, and letting $k \rightarrow \infty$ we deduce from Lemma 2.38 that $\psi_{t}\left(T_{B_{z}}\right)$ is holomorphic.

Proposition 2.42. Assume that $T_{B}$ satisfies quadratic estimates $\left\|Q_{t}^{B} f\right\| \approx\|f\|, f \in \mathcal{H}$, locally uniformly for all B such that $\left\|B-B_{0}\right\|_{\infty}<\varepsilon$, and let $\psi \in \Psi\left(S_{\nu}^{o}\right)$. Then we have the Lipschitz estimates

$$
\begin{gathered}
\left\|b\left(T_{B_{2}}\right)-b\left(T_{B_{1}}\right)\right\| \leqslant C\|b\|_{\infty}\left\|B_{2}-B_{1}\right\|_{\infty}, \quad b \in H_{\infty}\left(S_{v}^{o}\right), \\
\left\|\psi_{t}\left(T_{B_{2}}\right)-\psi_{t}\left(T_{B_{1}}\right)\right\|\|\leqslant\| B_{2}-B_{1} \|_{\infty},
\end{gathered}
$$

when $\left\|B_{i}-B_{0}\right\|<\varepsilon / 2, i=1,2$.
Proof. Let $B(z):=B_{1}+z\left(B_{2}-B_{1}\right) /\left\|B_{2}-B_{1}\right\|_{\infty}$, so that $B(z)$ is holomorphic in a neighbourhood of the interval $\left[0,\left\|B_{2}-B_{1}\right\|_{\infty}\right]$. In this neighbourhood, we have bounds $\left\|b\left(T_{B(z)}\right)\right\| \lesssim$ $\|b\|_{\infty}$ by Proposition 2.32 and holomorphic dependence on $z$ by Lemma 2.41(i). Schwarz' lemma now applies and proves that $\left\|\frac{d}{d z} b\left(T_{B(z)}\right)\right\| \lesssim\|b\|_{\infty}$ for all $z \in\left[0,\left\|B_{2}-B_{1}\right\|_{\infty}\right]$. This shows that

$$
\left\|b\left(T_{B_{2}}\right)-b\left(T_{B_{1}}\right)\right\| \leqslant \int_{0}^{\left\|B_{2}-B_{1}\right\|}\left\|\frac{d}{d t} b\left(T_{B(t)}\right)\right\| d t \leqslant C\|b\|_{\infty}\left\|B_{2}-B_{1}\right\|_{\infty} .
$$

The proof of the Lipschitz estimate for $\psi_{t}\left(T_{B}\right)$ similarly follows from Lemma 2.41(iii).

### 2.4. Decoupling of the Dirac equation

In Section 2.1 we introduced the Dirac type equation

$$
\begin{equation*}
\left(m d_{t, x}+B^{-1} m d_{t, x}^{*} B\right) F=0 \tag{2.9}
\end{equation*}
$$

satisfied by functions $F(t, x): \mathbf{R}_{ \pm}^{n+1} \rightarrow \bigwedge$. Of particular interest is when both terms vanish, i.e. when

$$
\left\{\begin{array} { l } 
{ d _ { t , x } F = 0 , }  \tag{2.10}\\
{ d _ { t , x } ^ { * } ( B F ) = 0 , }
\end{array} \quad \text { or equivalently } \quad \text { when } \left\{\begin{array}{l}
d F=-\mu \partial_{t} F, \\
d^{*}(B F)=\mu^{*} B \partial_{t} F .
\end{array}\right.\right.
$$

Consider a solution $F(t, x)=e^{\mp t\left|T_{B}\right|} f$ to (2.9) in $\mathbf{R}_{ \pm}^{n+1}$ as in Lemma 2.34, where $f, T_{B} f \in$ $E_{B}^{ \pm} \mathcal{V}_{B}$. Using $\partial_{t} F=-T_{B} F$ and Lemma 2.11, we get

$$
\begin{aligned}
\left.\left(m d_{t, x} F\right)\right|_{t=0} & =m\left(d-\mu T_{B}\right) f=\left(m d-N^{+} M_{B}^{-1}\left(m d+B^{-1} m d^{*} B\right)\right) f \\
& =M_{B}^{-1}\left(\left(M_{B}-B^{-1} N^{+} B\right) m d-B^{-1} N^{+} m d^{*} B\right) f \\
& =M_{B}^{-1}\left(-N^{-} m d-B^{-1} N^{+} m d^{*} B\right) f=M_{B}^{-1}\left(d \mu+B^{-1} d^{*} \mu^{*} B\right) f
\end{aligned}
$$

Thus we see that $d_{t, x} F=0=d_{t, x}^{*}(B F)$ at $t=0$ if and only if $d \mu f=0=d^{*} \mu^{*} B f$.
We may also rewrite Eq. (2.9) as

$$
\begin{equation*}
\left(d_{t, x} m+B^{-1} d_{t, x}^{*} m B\right) F(t, x)=0 \tag{2.11}
\end{equation*}
$$

by using $m d_{t, x}=\partial_{t}-d_{t, x} m$ and $m d_{t, x}^{*}=-\partial_{t}-d_{t, x}^{*} m$ from Lemma 2.7. Consider the case when both terms vanish, i.e. when

$$
\left\{\begin{array} { l } 
{ d _ { t , x } ( m F ) = 0 , }  \tag{2.12}\\
{ d _ { t , x } ^ { * } ( m B F ) = 0 , }
\end{array} \text { or equivalently } \quad \text { when } \left\{\begin{array}{l}
d F=\mu^{*} \partial_{t} F, \\
d^{*}(B F)=-\mu B \partial_{t} F .
\end{array}\right.\right.
$$

For $F(t, x)=e^{\mp t\left|T_{B}\right|} f$ solving (2.9) we also have

$$
\begin{aligned}
\left.d_{t, x}(m F)\right|_{t=0} & =\left.\partial_{t} F\right|_{t=0}-\left.m d_{t, x} F\right|_{t=0} \\
& =M_{B}^{-1}\left(d m+B^{-1} d^{*} m B\right) f-M_{B}^{-1}\left(d \mu+B^{-1} d^{*} \mu^{*} B\right) f \\
& =M_{B}^{-1}\left(d \mu^{*}+B^{-1} d^{*} \mu B\right) f .
\end{aligned}
$$

Thus we see that $d_{t, x}(m F)=0=d_{t, x}^{*}(m B F)$ at $t=0$ if and only if $d \mu^{*} f=0=d^{*} \mu B f$.
Definition 2.43. Introduce the closed, densely defined operators

$$
\hat{T}_{B}:=-M_{B}^{-1}\left(d \mu^{*}+B^{-1} d^{*} \mu B\right), \quad \check{T}_{B}:=-M_{B}^{-1}\left(d \mu+B^{-1} d^{*} \mu^{*} B\right)
$$

with domains

$$
\begin{aligned}
& \mathrm{D}\left(\hat{T}_{B}\right):=\left\{f \in \mathcal{H} ; \mu^{*} f \in \mathrm{D}(d), \mu B f \in \mathrm{D}\left(d^{*}\right)\right\}, \\
& \mathrm{D}\left(\check{T}_{B}\right):=\left\{f \in \mathcal{H} ; \mu f \in \mathrm{D}(d), \mu^{*} B f \in \mathrm{D}\left(d^{*}\right)\right\},
\end{aligned}
$$

respectively, and define the closed subspaces

$$
\begin{aligned}
& \hat{\mathcal{H}}_{B}:=\left\{f \in \mathcal{H} ; d(\mu f)=0=d^{*}\left(\mu^{*} B f\right)\right\}, \\
& \check{\mathcal{H}}_{B}:=\left\{f \in \mathcal{H} ; d\left(\mu^{*} f\right)=0=d^{*}(\mu B f)\right\} .
\end{aligned}
$$

Proposition 2.44. We have a topological splitting of $\mathcal{H}$ into closed subspaces

$$
\mathcal{H}=\hat{\mathcal{H}}_{B} \oplus \check{\mathcal{H}}_{B}
$$

Furthermore we have $T_{B}=\hat{T}_{B}+\check{T}_{B}$ with $\mathrm{D}\left(T_{B}\right)=\mathrm{D}\left(\hat{T}_{B}\right) \cap \mathrm{D}\left(\check{T}_{B}\right)$, and

$$
\overline{\mathrm{R}\left(\hat{T}_{B}\right)}=\hat{\mathcal{H}}_{B}=\mathrm{N}\left(\check{T}_{B}\right), \quad \overline{\mathrm{R}\left(\check{T}_{B}\right)}=\check{\mathcal{H}}_{B}=\mathrm{N}\left(\hat{T}_{B}\right)
$$

Thus, if we identify $\hat{T}_{B}$ with its restriction to $\hat{\mathcal{H}}_{B}$ and $\check{T}_{B}$ with its restriction to $\check{\mathcal{H}}_{B}$, then these are bisectorial operators with spectral and resolvent estimates as in Proposition 2.22.

Proof. Clearly $\mathrm{D}\left(\hat{T}_{B}\right) \cap \mathrm{D}\left(\check{T}_{B}\right)=\mathrm{D}(d) \cap B^{-1} \mathrm{D}\left(d^{*}\right)=\mathrm{D}\left(T_{B}\right)$ and $\hat{T}_{B}+\check{T}_{B}=T_{B}$. Furthermore Lemma 2.21(i) shows that $\mathrm{N}\left(\check{T}_{B}\right)=\hat{\mathcal{H}}_{B}$ and $\mathrm{N}\left(\hat{T}_{B}\right)=\check{\mathcal{H}}_{B}$.
(1) To show $\mathrm{R}\left(\hat{T}_{B}\right) \subset \hat{\mathcal{H}}_{B}$, let $f \in \mathrm{D}\left(\hat{T}_{B}\right)$ and use Lemma 2.11 to get

$$
\begin{aligned}
d \mu\left(\hat{T}_{B} f\right) & =-d \mu M_{B}^{-1}\left(d \mu^{*}+B^{-1} d^{*} \mu B\right) f \\
& =-d m N^{+} M_{B}^{-1}\left(d \mu^{*}+B^{-1} d^{*} \mu B\right) f \\
& =-d m M_{B}^{-1} B^{-1} N^{+} B\left(d \mu^{*}+B^{-1} d^{*} \mu B\right) f \\
& =-d m M_{B}^{-1}\left(B^{-1} N^{+} B-N^{-}\right) d \mu^{*} f \\
& =-d m d \mu^{*} f=m d^{2} \mu^{*} f=0,
\end{aligned}
$$

and similarly

$$
\begin{aligned}
d^{*} \mu^{*} B\left(\hat{T}_{B} f\right) & =-d^{*} m B\left(B^{-1} N^{-} B\right) M_{B}^{-1}\left(d \mu^{*}+B^{-1} d^{*} \mu B\right) f \\
& =-d^{*} m B M_{B}^{-1} N^{-}\left(d \mu^{*}+B^{-1} d^{*} \mu B\right) f \\
& =-d^{*} m B M_{B}^{-1}\left(N^{-}-B^{-1} N^{+} B\right) B^{-1} d^{*} \mu B f \\
& =d^{*} m d^{*} \mu B f=-m\left(d^{*}\right)^{2} \mu B f=0 .
\end{aligned}
$$

A similar calculation shows that $\mathrm{R}\left(\check{T}_{B}\right) \subset \check{\mathcal{H}}_{B}$.
(2) Next we show that for all $f \in \mathrm{D}\left(T_{B}\right)$, we have

$$
\left\|T_{B} f\right\| \approx\left\|\hat{T}_{B} f\right\|+\left\|\check{T}_{B} f\right\| .
$$

Using that $M_{B}$ and $B$ are isomorphisms, Lemma 2.21(i) with $\Gamma=d$ and $\widetilde{\Gamma}=d^{*}$ and orthogonality $\|\mu g\|^{2}+\left\|\mu^{*} g\right\|^{2}=\|g\|^{2}$, we obtain

$$
\begin{aligned}
\left\|\hat{T}_{B} f\right\|+\left\|\check{T}_{B} f\right\| & \approx\left(\left\|\mu^{*} d f\right\|+\left\|\mu d^{*} B f\right\|\right)+\left(\|\mu d f\|+\left\|\mu^{*} d^{*} B f\right\|\right) \\
& \approx\|d f\|+\left\|d^{*} B f\right\| \approx\left\|T_{B} f\right\| .
\end{aligned}
$$

(3) Clearly $\hat{\mathcal{H}}_{B} \cap \check{\mathcal{H}}_{B}=\{0\}$ and $\mathrm{R}\left(T_{B}\right) \subset \mathrm{R}\left(\hat{T}_{B}\right)+\mathrm{R}\left(\check{T}_{B}\right)$. Taking closures, using that $T_{B}$ has dense range in $\mathcal{H}$ and using (2) yields

$$
\mathcal{H}=\overline{\mathrm{R}\left(\hat{T}_{B}\right)} \oplus \overline{\mathrm{R}\left(\check{T}_{B}\right)}
$$

Thus from (1) it follows that $\overline{\mathrm{R}\left(\hat{T}_{B}\right)}=\hat{\mathcal{H}}_{B}$ and $\overline{\mathrm{R}\left(\check{T}_{B}\right)}=\check{\mathcal{H}}_{B}$ and that $\mathcal{H}=\hat{\mathcal{H}}_{B} \oplus \check{\mathcal{H}}_{B}$.
It follows that $T_{B}$ is diagonal in the splitting $\mathcal{H}=\hat{\mathcal{H}}_{B} \oplus \check{\mathcal{H}}_{B}$. We shall now further decompose the subspace $\hat{\mathcal{H}}_{B}$ into the subspaces of homogeneous $k$-vector field $\bigwedge^{k}$, which also are preserved by $T_{B}$. The same decomposition can be made for $\check{\mathcal{H}}_{B}$, but this is not useful since $T_{B}$ does not preserve these subspaces.

Definition 2.45. Let $\mathcal{H}^{k}:=L_{2}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\bigwedge^{k}\right)\right)$ and $\hat{\mathcal{H}}_{B}^{k}:=\mathcal{H}^{k} \cap \hat{\mathcal{H}}_{B}$.
Lemma 2.46. The operator $T_{B}$ preserve all subspaces $\hat{\mathcal{H}}_{B}^{0}, \hat{\mathcal{H}}_{B}^{1}, \hat{\mathcal{H}}_{B}^{2}, \ldots, \hat{\mathcal{H}}_{B}^{n+1}$ and $\check{\mathcal{H}}_{B}$ (but not the subspaces $\mathcal{H}^{k}$ ) and we have a splitting

$$
\hat{\mathcal{H}}_{B}=\hat{\mathcal{H}}_{B}^{0} \oplus \hat{\mathcal{H}}_{B}^{1} \oplus \hat{\mathcal{H}}_{B}^{2} \oplus \cdots \oplus \hat{\mathcal{H}}_{B}^{n+1} .
$$

Furthermore, for the operators $\hat{N}_{B}, \check{N}_{B}$ and $N$, we have mapping properties

$$
\hat{N}_{B}: \hat{\mathcal{H}}_{B}^{k} \rightarrow \hat{\mathcal{H}}_{B}^{k}, \quad \check{N}_{B}: \check{\mathcal{H}}_{B} \rightarrow \check{\mathcal{H}}_{B} \quad \text { and } \quad N: \hat{\mathcal{H}}^{1} \rightarrow \hat{\mathcal{H}}^{1} .
$$

Proof. As $d \mu^{*}, d^{*} \mu, B$ and $N^{ \pm}$all preserve $\mathcal{H}^{k}$ it follows that $T_{B}\left(\hat{\mathcal{H}}_{B}^{k}\right)=\hat{T}_{B}\left(\hat{\mathcal{H}}_{B}^{k}\right) \subset \hat{\mathcal{H}}_{B}^{k}$. To show the splitting of $\hat{\mathcal{H}}_{B}$ it suffices to prove

$$
\hat{\mathcal{H}}_{B} \subset \hat{\mathcal{H}}_{B}^{0} \oplus \hat{\mathcal{H}}_{B}^{1} \oplus \hat{\mathcal{H}}_{B}^{2} \oplus \cdots \oplus \hat{\mathcal{H}}_{B}^{n+1} .
$$

To this end, let $f \in \hat{\mathcal{H}}_{B}$ and write $f_{k}$ for the $\bigwedge^{k}$ part of $f$. Since $d \mu f=0=d^{*} \mu^{*} B f$ we get

$$
\begin{gathered}
0=(d \mu f)_{k+2}=d\left((\mu f)_{k+1}\right)=d \mu\left(f_{k}\right) \\
0=\left(d^{*} \mu^{*} B f\right)_{k-2}=d^{*}\left(\left(\mu^{*} B f\right)_{k-1}\right)=d^{*} \mu^{*}\left((B f)_{k}\right)=d^{*} \mu^{*} B\left(f_{k}\right)
\end{gathered}
$$

for all $k$. Thus $f_{k} \in \hat{\mathcal{H}}_{B}^{k}$ and $f \in \bigoplus_{k} \hat{\mathcal{H}}_{B}^{k}$.
To prove the mapping properties for $\hat{N}_{B}$, note that if $f=f_{1}+f_{2}$ in the splitting $\mathcal{H}=$ $B^{-1} N^{+} \mathcal{H} \oplus N^{-} \mathcal{H}$, then $f \in \hat{\mathcal{H}}_{B}$ if and only if $\mu f_{2} \in \mathrm{D}(d)$ and $\mu B f_{1} \in \mathrm{D}\left(d^{*}\right)$, according to Definition 2.43. Clearly $f_{1}-f_{2}=\hat{N}_{B}(f) \in \hat{\mathcal{H}}_{B}$ if $f \in \hat{\mathcal{H}}_{B}$. Since $\hat{N}_{B}$ preserves $\mathcal{H}^{k}$, the desired mapping property follows. The proofs for $\check{N}_{B}$ and $N$ are similar.

Lemma 2.47. With $\langle\cdot, \cdot\rangle_{B}$ denoting the duality from Definition 2.16, we have dual operators

$$
\left(\hat{T}_{B}\right)^{\prime}=-\hat{T}_{B^{*}}, \quad\left(\check{T}_{B}\right)^{\prime}=-\check{T}_{B^{*}}
$$

and restricted dualities $\left\langle\hat{\mathcal{H}}_{B}, \hat{\mathcal{H}}_{B^{*}}\right\rangle_{B},\left\langle\check{\mathcal{H}}_{B}, \check{\mathcal{H}}_{B^{*}}\right\rangle_{B}$ and $\left\langle\hat{\mathcal{H}}_{B}^{k}, \hat{\mathcal{H}}_{B^{*}}^{k}\right\rangle_{B}$ for all $k$. In the case $k=1$, we shall write the duality as $\left\langle\hat{\mathcal{H}}^{1}, \hat{\mathcal{H}}^{1}\right\rangle_{A}$.

Proof. The proofs of $\hat{T}_{B}^{\prime}=-\hat{T}_{B^{*}}$ and $\check{T}_{B}^{\prime}=-\check{T}_{B^{*}}$ are similar to that of $T_{B}^{\prime}=-T_{B^{*}}$ in Proposition 2.17. From this we get that the annihilator of $\hat{\mathcal{H}}_{B}=\overline{\mathrm{R}\left(\hat{T}_{B}\right)}$ is $\mathrm{N}\left(\hat{T}_{B^{*}}\right)=\check{\mathcal{H}}_{B^{*}}$ which is a complement of $\hat{\mathcal{H}}_{B^{*}}$. Thus we see from Remark 2.19 that $\left\langle\hat{\mathcal{H}}_{B}, \hat{\mathcal{H}}_{B^{*}}\right\rangle_{B}$ is a duality. The proof of the duality $\left\langle\check{\mathcal{H}}_{B}, \mathscr{\mathcal { H }}_{B^{*}}\right\rangle_{B}$ is similar. Also, since $B N^{+}-N^{-} B$ in the definition of $\langle\cdot, \cdot\rangle_{B}$ preserves $\mathcal{H}^{k}$ we also have dualities $\left\langle\hat{\mathcal{H}}_{B}^{k}, \hat{\mathcal{H}}_{B^{*}}^{k}\right\rangle_{B}$ for all $k$.

## Remark 2.48.

- The subspace of vector fields with a curl-free tangential part

$$
\hat{\mathcal{H}}_{B}^{1}=\left\{f \in L_{2}\left(\mathbf{R}^{n} ; \bigwedge^{1}\right) ; d\left(e_{0} \wedge f\right)=0\right\}
$$

is independent of $B$ and coincides with the space $\hat{\mathcal{H}}^{1}$ from Section 1.1. Furthermore, the operator $T_{A}$ there coincides with $\left.T_{B}\right|_{\hat{\mathcal{H}}^{1}}=\left.\hat{T}_{B}\right|_{\hat{\mathcal{H}}^{1}}$.

- The dense subspace $\mathcal{V}_{B} \subset \mathcal{H}$ splits $\mathcal{V}_{B}=\hat{\mathcal{V}}_{B} \oplus \breve{\mathcal{V}}_{B}$ with Proposition 2.44, where $\hat{\mathcal{V}}_{B}:=\mathcal{V}_{B} \cap$ $\hat{\mathcal{H}}_{B} \subset \hat{\mathcal{H}}_{B}$ and $\check{\mathcal{V}}_{B}:=\mathcal{V}_{B} \cap \check{\mathcal{H}}_{B} \subset \check{\mathcal{H}}_{B}$ are dense subspaces, and we can further decompose $\hat{\mathcal{V}}_{B}$ into homogeneous $k$-vector fields, $\hat{\mathcal{V}}_{B}=\hat{\mathcal{V}}_{B}^{0} \oplus \hat{\mathcal{V}}_{B}^{1} \oplus \cdots \oplus \hat{\mathcal{V}}_{B}^{n+1}$, where $\hat{\mathcal{V}}_{B}^{k}:=\mathcal{V}_{B} \cap \hat{\mathcal{H}}_{B}^{k}$.
- Furthermore, all these dense subspaces $\hat{\mathcal{V}}_{B}^{k}, \hat{\mathcal{V}}_{B}, \check{\mathcal{V}}_{B}$ of $\hat{\mathcal{H}}_{B}^{k}, \hat{\mathcal{H}}_{B}, \check{\mathcal{H}}_{B}$ splits algebraically into Hardy spaces similar to Lemma 2.34, e.g. $\hat{\mathcal{V}}_{B}^{k}=E_{B}^{+} \hat{\mathcal{V}}_{B}^{k}+E_{B}^{-} \hat{\mathcal{V}}_{B}^{k}$. Here $f \in E_{B}^{ \pm} \hat{\mathcal{V}}_{B}$ if and only if $F(t, x)$ satisfies (2.10), and $f \in E_{B}^{ \pm} \check{\mathcal{V}}_{B}$ if and only if $F(t, x)$ satisfies (2.12).


### 2.5. Operator equations and estimates for solutions

Our objective in this section is to set up our boundary operator method for solving the boundary value problems $(\mathrm{Neu}-A),(\operatorname{Reg}-A),\left(\mathrm{Neu}^{\perp}-A\right)$ and (Dir- $A$ ) as well as the transmission problem ( $\operatorname{Tr}-B^{k} \alpha^{ \pm}$). Under the assumption that $T_{B}$ has quadratic estimates, which is made throughout this section, we first show that solutions $F(t, x)$ are determined by their traces $f$ through the reproducing Cauchy type formula $F_{t}=e^{-t\left|T_{B}\right|} f$. Recall from Lemma 2.46 that both operators $E_{B}$ and $N_{B}$ preserve all subspaces $\hat{\mathcal{H}}_{B}^{k}$. We shall write $E_{B^{k}}=\left.E_{B}\right|_{\hat{\mathcal{H}}_{B}^{k}}=\operatorname{sgn}\left(\left.T_{B}\right|_{\hat{\mathcal{H}}_{B}^{k}}\right)$ and $N_{B^{k}}=\left.N_{B}\right|_{\hat{\mathcal{H}}_{B}^{k}}$ for the restrictions. In particular, we write $A=B^{1}$ when $k=1$.

Lemma 2.49. Assume that $T_{B}$ satisfies quadratic estimates, let $f \in \mathcal{H}^{k}$ and let $(0, \infty) \ni t \mapsto$ $F_{t}(x)=F(t, x) \in \mathcal{H}^{k}$ be a family of functions. Then the following are equivalent.
(i) $f \in E_{B}^{+} \hat{\mathcal{H}}_{B}^{k}$ and $F_{t}=e^{-t\left|T_{B}\right|} f$.
(ii) $F_{t} \in C^{1}\left(\mathbf{R}_{+} ; \mathcal{H}^{k}\right)$ and satisfies the equations

$$
\left\{\begin{array}{l}
d_{t, x}^{*}(B(x) F(t, x))=0 \\
d_{t, x} F(t, x)=0
\end{array}\right.
$$

and have $L_{2}$ limits $\lim _{t \rightarrow 0^{+}} F_{t}=f$ and $\lim _{t \rightarrow \infty} F_{t}=0$.
In fact, such $F_{t}$ belong to $C^{j}\left(\mathbf{R}_{+} ; \mathcal{H}^{k}\right)$ for all $j \geqslant 1$ and are in one-to-one correspondence with the trace $f \in E_{B}^{+} \hat{\mathcal{H}}_{B}^{k}$, and we have equivalences of norms

$$
\|f\| \approx \sup _{t>0}\left\|F_{t}\right\| \approx\left\|t \partial_{t} F_{t}\right\|
$$

The corresponding reproducing formula $F_{t}=e^{t\left|T_{B}\right|} f, f=\left.F\right|_{\mathbf{R}^{n}} \in E_{B}^{-} \hat{\mathcal{H}}_{B}^{k}$ is also valid for $F$ solving the equations in $\mathbf{R}_{-}^{n+1}$, and the corresponding estimates hold.

Proof. (i) implies (ii). As in the proof of Lemma 2.34, from Lemma 2.28 it follows that $\lim _{t \rightarrow 0} F_{t}=f$ and $\lim _{t \rightarrow \infty} F_{t}=0$, and also that $\partial_{t}^{j} F_{t}=\left(-\left|T_{A}\right|\right)^{j} e^{-t\left|T_{A}\right|} f$. Therefore $F_{t} \in$ $C^{j}\left(\mathbf{R}_{+} ; \mathcal{H}^{k}\right)$ for all $j$. For $j=1$, we get $0=\left(\partial_{t}+\left|T_{B}\right|\right) F_{t}=\left(\partial_{t}+T_{B}\right) F_{t}$, since $F_{t} \in E_{B}^{+} \hat{\mathcal{H}}_{B}^{k}$. As explained in Section 2.4 this is equivalent with the two equations $d_{t, x} F=0$ and $d_{t, x}^{*}(B F)=0$.
(ii) implies (i). The two equations can be written $\left(-\mu^{*} \partial_{t}+d^{*}\right) B F=0$ and $\left(\mu \partial_{t}+d\right) F=0$. Applying $\mu^{*}$ and $\mu$, respectively, to these equations and using nilpotence, we obtain $\mu^{*} d^{*} B F=0$ and $\mu d F=0$. Therefore $F_{t} \in \hat{\mathcal{H}}_{B}^{k}$, and since $\hat{\mathcal{H}}_{B}^{k}$ is closed we also have $f \in \hat{\mathcal{H}}_{B}^{k}$.

Next we write $F_{t}=F_{t}^{+}+F_{t}^{-}$, where $F_{t}^{ \pm}:=E_{B}^{ \pm} F_{t} \in E_{B}^{ \pm} \hat{\mathcal{H}}_{B}^{k}$ and similarly $f=f^{+}+f^{-}$. We rewrite the equations satisfied by $F_{t}$ as $\partial_{t} F_{t}+T_{B} F_{t}=0$. Applying the projections $E_{B}^{ \pm}$to this equation yields

$$
\partial_{t} F_{t}^{+}+\left|T_{B}\right| F_{t}^{+}=0, \quad \partial_{t} F_{t}^{-}-\left|T_{B}\right| F_{t}^{-}=0
$$

since $T_{B} F_{t}^{ \pm}= \pm\left|T_{B}\right| F_{t}^{ \pm}$. Fix $t>0$. Then it follows that $e^{(t-s)\left|T_{B}\right|} F_{s}^{-}$is constant for $s \in(t, \infty)$ and that $e^{(s-t)\left|T_{B}\right|} F_{s}^{+}$is constant for $s \in(0, t)$. Taking limits, this shows that

$$
\begin{gathered}
F_{t}^{-}=\lim _{s \rightarrow t^{+}} e^{(t-s)\left|T_{B}\right|} F_{s}^{-}=\lim _{s \rightarrow \infty} e^{(t-s)\left|T_{B}\right|} F_{s}^{-}=0 \quad \text { and } \\
F_{t}^{+}=\lim _{s \rightarrow t^{-}} e^{(s-t)\left|T_{B}\right|} F_{s}^{+}=\lim _{s \rightarrow 0} e^{(s-t)\left|T_{B}\right|} F_{s}^{+}=e^{-t\left|T_{B}\right|} f^{+} .
\end{gathered}
$$

Therefore $f=f^{+} \in E_{B}^{+} \hat{\mathcal{H}}_{B}^{k}$ and $F_{t}=e^{-t\left|T_{B}\right|} f$.
To prove the norm estimates we note that $\|f\|=\lim _{t \rightarrow 0}\left\|F_{t}\right\|$. By Proposition 2.32, $e^{-t\left|T_{B}\right|}$ are uniformly bounded and thus $\sup _{t>0}\left\|F_{t}\right\| \lesssim\|f\|$. Furthermore, using Proposition 2.31 with $\psi(z)=z e^{-|z|}$ shows that $\left\|t \partial_{t} F_{t}\right\|\|\|f\|$.

Remark 2.50. In proving that (ii) implies (i), it suffices to assume that $\left\|F_{t}\right\|$ grows at most polynomially when $t \rightarrow \infty$. Indeed, from the equation $F_{t}^{-}=e^{-(s-t)\left|T_{B}\right|} F_{s}^{-}$for $s>t$, it then follows that

$$
\left\|T_{B}^{k} F_{t}^{-}\right\|=(s-t)^{-k}\left\|\left((s-t) T_{B}\right)^{k} e^{-(s-t)\left|T_{B}\right|} F_{s}^{-}\right\| \leqslant C(s-t)^{-k}\left\|F_{s}\right\| \rightarrow 0
$$

when $s \rightarrow \infty$. Since $T_{B}$ is injective, this shows that $F_{t}^{-}=0$ and therefore that $F_{t}=e^{-t\left|T_{B}\right|} f \in$ $E_{B}^{+} \hat{\mathcal{H}}_{B}^{k}$ as before.

We now proceed by showing how, given data $g$ in $\left(\operatorname{Tr}-B^{k} \alpha^{ \pm}\right)$, we can solve for the trace $f=\left.F\right|_{\mathbf{R}^{n}}$ by using the boundary operators $E_{B}$ and $N_{B}$.

Lemma 2.51. Assume that $T_{B}$ satisfies quadratic estimates, so that the Hardy projections $E_{B}^{ \pm}$ are bounded by Proposition 2.32, let $\alpha^{ \pm} \in \mathbf{C}$ be given jump parameters and define the associated spectral point $\lambda:=\left(\alpha^{+}+\alpha^{-}\right) /\left(\alpha^{+}-\alpha^{-}\right)$. Then

$$
\lambda-E_{B^{k}} N_{B^{k}}: \hat{\mathcal{H}}_{B}^{k} \rightarrow \hat{\mathcal{H}}_{B}^{k}
$$

is an isomorphism if and only if the transmission problem $\left(\operatorname{Tr}-B^{k} \alpha^{ \pm}\right)$is well posed.
Proof. If we identify the $k$-vector fields $F^{ \pm}(t, x)$ in Transmission problem $\left(\operatorname{Tr}-B^{k} \alpha^{ \pm}\right)$with the boundary traces $f^{ \pm}$and write $f=f^{+}+f^{-}$using Lemma 2.49, then we see that the transmission problem is equivalent with the system of equations

$$
\left\{\begin{array}{l}
N_{B^{k}}^{+}\left(\alpha^{-} E_{B^{k}}^{+}-\alpha^{+} E_{B^{k}}^{-}\right) f=N_{B^{k}}^{+} g, \\
N_{B^{k}}^{-}\left(\alpha^{+} E_{B^{k}}^{+}-\alpha^{-} E_{B^{k}}^{-}\right) f=N_{B^{k}}^{-} g .
\end{array}\right.
$$

Using $E_{B^{k}}^{ \pm}=\frac{1}{2}\left(I \pm E_{B^{k}}\right)$ and adding up the equations, we see that the system is equivalent with the equation

$$
\left(\lambda-E_{B^{k}} N_{B^{k}}\right) f=\frac{2}{\alpha^{+}-\alpha^{-}} E_{B^{k}} g
$$

This proves the lemma.

Next we consider $k=1$ and the boundary value problems (Neu- $A$ ), (Reg- $A$ ) and ( $\mathrm{Neu}^{\perp}-A$ ). By Lemma 2.49, we have the following.
(1) (Neu-A) is well posed if and only if the restricted projection $N_{A}^{-}: E_{A}^{+} \hat{\mathcal{H}}^{1} \rightarrow N_{A}^{-} \hat{\mathcal{H}}^{1}$ is an isomorphism.
(2) (Reg-A) is well posed if and only if the restricted projection $N_{A}^{+}: E_{A}^{+} \hat{\mathcal{H}}^{1} \rightarrow N_{A}^{+} \hat{\mathcal{H}}^{1}$ is an isomorphism, or equivalently if and only if $N^{+}: E_{A}^{+} \hat{\mathcal{H}}^{1} \rightarrow N^{+} \hat{\mathcal{H}}^{1}$ is an isomorphism. Note that both $N_{A}^{+}$and $N^{+}$project along $N^{-} \mathcal{H}^{1}$.
(3) $\left(\mathrm{Neu}^{\perp}-A\right)$ is well posed if and only if the restricted projection $N^{-}: E_{A}^{+} \hat{\mathcal{H}}^{1} \rightarrow N^{-} \hat{\mathcal{H}}^{1}$ is an isomorphism.

Proposition 2.52. Assume that $T_{A}$ satisfies quadratic estimates. Then (Reg-A) is well posed if and only if $\left(\mathrm{Neu}^{\perp}-A^{*}\right)$ is well posed.

Proof. We need to show that if $N^{+}: E_{A}^{+} \hat{\mathcal{H}}^{1} \rightarrow N^{+} \hat{\mathcal{H}}^{1}$ is an isomorphism, then so is $N^{-}: E_{A^{*}}^{+} \hat{\mathcal{H}}^{1} \rightarrow N^{-} \hat{\mathcal{H}}^{1}$. The proof uses two facts. First that we have adjoint operators $\left(E_{A}\right)^{\prime}=$ $-E_{A^{*}}$ and $N^{\prime}=N$ according to Proposition 2.17 and Lemma 2.47. As in Remark 2.19, this shows that we have dual spaces $\left\langle E_{A}^{+} \hat{\mathcal{H}}^{1}, E_{A^{*}}^{-} \hat{\mathcal{H}}^{1}\right\rangle_{A}$ and $\left\langle N^{+} \hat{\mathcal{H}}^{1}, N^{+} \hat{\mathcal{H}}^{1}\right\rangle_{A}$, and we see that

$$
\left\langle N^{+} f, g\right\rangle_{A}=\langle f, g\rangle_{A}=\left\langle f, E_{A^{*}}^{-} g\right\rangle_{A},
$$

for all $f \in E_{A}^{+} \hat{\mathcal{H}}^{1}$ and $g \in N^{+} \hat{\mathcal{H}}^{1}$. Therefore, the restricted projections $N^{+}: E_{A}^{+} \hat{\mathcal{H}}^{1} \rightarrow N^{+} \hat{\mathcal{H}}^{1}$ and $E_{A^{*}}^{-}: N^{+} \hat{\mathcal{H}}^{1} \rightarrow E_{A^{*}}^{-} \hat{\mathcal{H}}^{1}$ are adjoint.

Secondly, if $R_{1}^{ \pm}$and $R_{2}^{ \pm}$are two pairs of complementary projections in a Hilbert space $\mathcal{H}$, as in Definition 2.13, then $R_{1}^{-}: R_{2}^{+} \mathcal{H} \rightarrow R_{1}^{-} \mathcal{H}$ has a priori estimates, as in Remark 2.18, if and only if $R_{2}^{-}: R_{1}^{+} \mathcal{H} \rightarrow R_{2}^{-} \mathcal{H}$ has a priori estimates. Indeed, both statements are seen to be equivalent with that the subspaces $R_{1}^{+} \mathcal{H}$ and $R_{2}^{+} \mathcal{H}$ are transversal, i.e. that the estimate $\left\|f_{1}+f_{2}\right\| \approx\left\|f_{1}\right\|+\left\|f_{2}\right\|$ holds for all $f_{1} \in R_{1}^{+} \mathcal{H}$ and $f_{2} \in R_{2}^{+} \mathcal{H}$. To see this, assume that $\left\|R_{1}^{-} f_{2}\right\| \gtrsim\left\|f_{2}\right\|$ holds for all $f_{2} \in R_{2}^{+} \mathcal{H}$. Then $\left\|f_{2}\right\| \lesssim\left\|R_{1}^{-}\left(f_{1}+f_{2}\right)\right\| \lesssim\left\|f_{1}+f_{2}\right\|$ for all $f_{i} \in R_{i}^{+} \mathcal{H}$, which proves transversality. Conversely, assume that $R_{1}^{+} \mathcal{H}$ and $R_{2}^{+} \mathcal{H}$ are transversal. Then $f_{2}-R_{1}^{-} f_{2}=R_{1}^{+} f_{2}=$ : $f_{1} \in R_{1}^{+} \mathcal{H}$ for all $f_{2} \in R_{2}^{+} \mathcal{H}$. Therefore $\left\|R_{1}^{-} f_{2}\right\|=\left\|f_{2}-f_{1}\right\| \approx\left\|f_{2}\right\|+\left\|f_{1}\right\| \geqslant\left\|f_{2}\right\|$. The same argument can be used to show that transversality also holds if and only if $\left\|R_{2}^{-} f_{1}\right\| \gtrsim\left\|f_{1}\right\|$ holds for all $f_{1} \in R_{1}^{+} \mathcal{H}$.

To prove the proposition, assume that $N^{+}: E_{A}^{+} \hat{\mathcal{H}}^{1} \rightarrow N^{+} \hat{\mathcal{H}}^{1}$ is an isomorphism. It follows that the adjoint operator $E_{A^{*}}^{-}: N^{+} \hat{\mathcal{H}}^{1} \rightarrow E_{A^{*}}^{-} \hat{\mathcal{H}}^{1}$ also is an isomorphism. Using the second fact above twice, shows that

$$
E_{A}^{-}: N^{-} \hat{\mathcal{H}}^{1} \rightarrow E_{A}^{-} \hat{\mathcal{H}}^{1} \quad \text { and } \quad N^{-}: E_{A^{*}}^{+} \hat{\mathcal{H}}^{1} \rightarrow N^{-} \hat{\mathcal{H}}^{1}
$$

have a priori estimates. As these are adjoint operators as well, both must in fact be isomorphisms. In particular we have shown that $\left(\mathrm{Neu}^{\perp}-A^{*}\right)$ is well posed. The proof of the converse implication is similar.

When perturbing $A$, it is preferable not to have operators defined on spaces like $E_{A}^{+} \hat{\mathcal{H}}$, which varies with $A$. We have the following.

Lemma 2.53. Assume that $T_{A}$ satisfies quadratic estimates.
(1) If $I-E_{A} N_{A}: \hat{\mathcal{H}}^{1} \rightarrow \hat{\mathcal{H}}^{1}$ is an isomorphism, then Neumann problem (Neu-A) is well posed.
(2) If $I+E_{A} N_{A}: \hat{\mathcal{H}}^{1} \rightarrow \hat{\mathcal{H}}^{1}$ is an isomorphism, or if $I+E_{A} N: \hat{\mathcal{H}}^{1} \rightarrow \hat{\mathcal{H}}^{1}$ is an isomorphism, then Regularity problem $(\operatorname{Reg}-A)$ is well posed.
(3) If $I-E_{A} N: \hat{\mathcal{H}}^{1} \rightarrow \hat{\mathcal{H}}^{1}$ is an isomorphism, then Neumann problem $\left(\mathrm{Neu}^{\perp}-A\right)$ is well posed.

Proof. Assume for example that $I+E_{A} N_{A}$ is an isomorphism. We need to prove that $N_{A}^{+}: E_{A}^{+} \hat{\mathcal{H}}^{1} \rightarrow N_{A}^{+} \hat{\mathcal{H}}^{1}$ is an isomorphism. Note that if $N_{A}^{+} f=g$ where $f \in E_{A}^{+} \hat{\mathcal{H}}^{1}$, then

$$
g=N_{A}^{+} f=\frac{1}{2}\left(I+N_{A}\right) f=\frac{1}{2}\left(E_{A}+N_{A}\right) f=\frac{1}{2} E_{A}\left(I+E_{A} N_{A}\right) f .
$$

If $g \in N_{A}^{+} \hat{\mathcal{H}}^{1}$, let $f:=2\left(E_{A}+N_{A}\right)^{-1} g$. Then it follows that $0=N_{A}^{-} g=\frac{1}{2}\left(E_{A}+N_{A}\right)\left(E_{A}^{-} f\right)$, since $N_{A}^{-}\left(E_{A}+N_{A}\right)=\frac{1}{2}\left(E_{A}+N_{A}-I-N_{A} E_{A}\right)=\left(E_{A}+N_{A}\right) E_{A}^{-}$, so $E_{A}^{-} f=0$ and therefore $f \in E_{A}^{+} \hat{\mathcal{H}}^{1}$.

A similar calculation proves well-posedness of the other boundary value problems.
Remark 2.54. More generally, letting $k=1$ and $\left(\alpha^{+}, \alpha^{-}\right)=(1,0)$, i.e. $\lambda=1$, in Lemma 2.51, we see that $I-E_{A} N_{A}$ is an isomorphism if and only if the restricted projections

$$
N_{A}^{+}: E_{A}^{-} \hat{\mathcal{H}}^{1} \rightarrow N_{A}^{+} \hat{\mathcal{H}}^{1} \quad \text { and } \quad N_{A}^{-}: E_{A}^{+} \hat{\mathcal{H}}^{1} \rightarrow N_{A}^{-} \hat{\mathcal{H}}^{1}
$$

are isomorphisms. Similarly, if $\left(\alpha^{+}, \alpha^{-}\right)=(0,1)$, i.e. $\lambda=-1$ in Lemma 2.51, we see that $I+$ $E_{A} N_{A}$ is an isomorphism if and only if the restricted projections $N_{A}^{+}: E_{A}^{+} \hat{\mathcal{H}}^{1} \rightarrow N_{A}^{+} \hat{\mathcal{H}}^{1}$ and $N_{A}^{-}: E_{A}^{-} \hat{\mathcal{H}}^{1} \rightarrow N_{A}^{-} \hat{\mathcal{H}}^{1}$ are isomorphisms.

We next turn to the Dirichlet problem (Dir-A), where we aim to prove an analogue of Lemma 2.49 which characterises the solution $U_{t}$ as a Poisson integral of the boundary trace $u$. As discussed in the introduction, we shall use $\left(\mathrm{Neu}^{\perp}-A\right)$ to construct the solution $U_{t}$. Given Dirichlet data $u \in L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}\right)$, we form $u e_{0} \in N^{-} \hat{\mathcal{H}}^{1}$. It then follows from Lemma 2.12 that the vector field $F_{t}$ solving $\left(\mathrm{Neu}^{\perp}-A\right)$ with data $u e_{0}$, has a normal component $U:=F_{0}$ which satisfies the second order equation (1.1). We now define the Poisson integral of $u$ to be

$$
\mathcal{P}_{t}(u):=\left(F_{t}, e_{0}\right), \quad \text { when } F_{t}=e^{-t\left|T_{A}\right|} f \text { and }\left(f, e_{0}\right)=u
$$

Lemma 2.55. Assume that $T_{A}$ satisfies quadratic estimates and that Neumann problem ( $\mathrm{Neu}^{\perp}$ A) is well posed. Let $u \in L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}\right)$ and let $(0, \infty) \ni t \mapsto U_{t}(x)=U(t, x) \in L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}\right)$ be a family of functions. Then the following are equivalent.
(i) $U_{t}=\mathcal{P}_{t} u$ for all $t>0$.
(ii) $U_{t} \in C^{2}\left(\mathbf{R}_{+} ; L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}\right)\right), \nabla_{t, x} U_{t} \in C^{1}\left(\mathbf{R}_{+} ; L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}^{n+1}\right)\right)$ and $U$ satisfies the equation

$$
\operatorname{div}_{t, x} A(x) \nabla_{t, x} U(t, x)=0
$$

and we have $L_{2}$ limits $\lim _{t \rightarrow 0^{+}} U_{t}=u, \lim _{t \rightarrow \infty} U_{t}=0$ and $\lim _{t \rightarrow \infty} \nabla_{t, x} U_{t}=0$.
If this holds, then $U_{t}, \nabla_{t, x} U_{t} \in C^{j}\left(\mathbf{R}_{+} ; L_{2}\left(\mathbf{R}^{n}\right)\right)$ for all $j \geqslant 1$. Furthermore, $U_{t}$ is in one-to-one correspondence with the trace $u$, and we have equivalences of norms

$$
\|u\| \approx \sup _{t>0}\left\|U_{t}\right\| \approx\left\|t \partial_{t} U_{t}\right\| .
$$

Proof. (i) implies (ii). Assume that $U_{t}=\left(F_{t}, e_{0}\right)$, where $F_{t}=e^{-t\left|T_{A}\right|} f$ and $\left(f, e_{0}\right)=u$. As in the proof of Lemma 2.34, from Lemma 2.28 it follows that $F_{t} \rightarrow f$, and therefore $U_{t} \rightarrow u$, and also that $\partial_{t}^{j} F_{t}=\left(-\left|T_{A}\right|\right)^{j} e^{-t\left|T_{A}\right|} f$. Therefore $F_{t} \in C^{j}\left(\mathbf{R}_{+} ; L_{2}\left(\mathbf{R}^{n}\right)\right)$ for all $j$, and so does $U_{t}$. For $j=1$, we see that $F$ satisfies the Dirac type equation since $\left|T_{A}\right| F_{t}=T_{A} F_{t}$, and Lemma 2.12 thus shows that $U_{t}$ satisfies the second order equation. Furthermore, we note from the expression (1.5) for $T_{A}$, that $\nabla_{x} U_{t}=\left(\partial_{t} F_{t}\right)_{\|}$. Thus $\nabla_{x} U_{t} \in C^{j}\left(\mathbf{R}_{+} ; L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}\right)\right)$ for all $j$, and yet another application of Lemma 2.28 shows that $U_{t}=o(1)$ and $\nabla_{t, x} U_{t}=o(1 / t)$ when $t \rightarrow \infty$.
(ii) implies (i). Assume $U_{t}$ has the stated properties and boundary trace $u$. Consider the family of vector fields $G_{t}:=\nabla_{t, x} U_{t}$. Since these satisfy Lemma 2.49(ii) for $t \geqslant s>0$ with boundary trace $G_{s}$, we obtain that $G_{s+t}=e^{-t\left|T_{A}\right|} G_{s}$ for all $s, t>0$. For the normal components, this means that

$$
\partial_{0} U_{s+t}=\mathcal{P}_{t}\left(\partial_{0} U_{s}\right),
$$

or equivalently that $\partial_{s}\left(U_{s+t}-\mathcal{P}_{t}\left(U_{s}\right)\right)=0$. Since $\lim _{s \rightarrow \infty} U_{s}=0$, we must have $U_{s+t}=\mathcal{P}_{t}\left(U_{s}\right)$ for all $s, t>0$. Letting $s \rightarrow 0$, we conclude that $U_{t}=\mathcal{P}_{t}(u)$.

The equivalence of norms $\|u\| \approx \sup _{t>0}\left\|U_{t}\right\|$ follows from the uniform boundedness of the operators $\mathcal{P}_{t}$. For the equivalence $\|u\| \approx\left\|t \partial_{t} U_{t}\right\|$ we use that $\left(\mathrm{Neu}^{\perp}-A\right)$ is well posed and the corresponding square function estimate for $F_{t}$ from Lemma 2.49 to get $\|u\| \approx\|f\| \approx\left\|t \partial_{t} F_{t}\right\| \approx$ $\left\|t \partial_{t} U_{t}\right\| \|$, since for all $t>0$ we have $\left\|\partial_{t} U_{t}\right\|=\left\|N^{-}\left(\partial_{t} F_{t}\right)\right\| \approx\left\|\partial_{t} F_{t}\right\|$.

We end this section with the proof of the non-tangential estimate $\left\|\widetilde{N}_{*}(F)\right\| \approx\|f\|$ in Theorem 1.1.

Proposition 2.56. Assume that $T_{A}$ satisfies quadratic estimates. Let $F_{t}=e^{-t\left|T_{A}\right|} f$, where $f \in$ $E_{A}^{+} \hat{\mathcal{H}}^{1}$. Then $\|f\| \approx\left\|\widetilde{N}_{*}(F)\right\|$, where the non-tangential maximal function is

$$
\tilde{N}_{*}(F)(x):=\sup _{t>0}\left({\left.\underset{D(t, x)}{ } f|F(s, y)|^{2} d s d y\right)^{1 / 2},, ~ \text {, }}\right.
$$

and $D(t, x):=\left\{(s, y) \in \mathbf{R}_{+}^{n+1} ;|s-t|<c_{0} t,|y-x|<c_{1} t\right\}$, for given constants $c_{0} \in(0,1)$ and $c_{1}>0$.

The proof uses the following lemma.
Lemma 2.57. Let $f \in \hat{\mathcal{H}}^{1}$ and define $H_{t}=\left(1+i t T_{B}\right)^{-1} f \in \hat{\mathcal{H}}^{1}$. Write $H_{t}=H_{t}^{1,0} e_{0}+H_{t}^{1, \|}$ and $f=f^{1,0} e_{0}+f^{1, \|}$.
(i) The normal component $H_{t}^{1,0}$ satisfies the second order divergence form equation

$$
\left[\begin{array}{ll}
1 & i t \operatorname{div}]
\end{array}\left[\begin{array}{cc}
a_{\perp \perp} & a_{\perp \|} \\
a_{\| \perp} & a_{\| \|}
\end{array}\right]\left[\begin{array}{c}
1 \\
i t \nabla
\end{array}\right] H_{t}^{1,0}=\left[\begin{array}{cc}
1 & \text { it div }]
\end{array}\left[\begin{array}{c}
a_{\perp \perp} f^{1,0} \\
-a_{\| \|} f^{1, \|}
\end{array}\right]\right. \text {, }\right.
$$

where we identify normal vectors $u e_{0}$ with scalars $u$, and the tangential component $H_{t}{ }^{1, \|}$ satisfies

$$
H_{t}^{1, \|}=f^{1, \|}+i t \nabla H_{t}^{1,0}
$$

(ii) There exist $p<2$ and $q>2$ such that for any fixed $r_{0}<\infty$ we have

$$
\left(f_{B\left(x, r_{0} t\right)}\left|H_{t}^{1,0}\right|^{q}\right)^{1 / q}+\left(\underset{B\left(x, r_{0} t\right)}{ }\left|H_{t}^{1, \|}\right|^{p}\right)^{1 / p} \lesssim M\left(|f|^{p}\right)^{1 / p}(x),
$$

for all $x \in \mathbf{R}^{n}$ and $t>0$. Here $M(f)(x):=\sup _{r>0} f_{B(x, r)}|f(y)|$ dy denotes the HardyLittlewood maximal function.

Proof. (i) Multiplying the equation $\left(1+i t T_{B}\right) H_{t}=f$ with $M_{B}$ we get

$$
\left(m(\mu+i t d)+B^{-1}\left(m\left(-\mu^{*}+i t d^{*}\right)\right) B\right) H_{t}=M_{B} f .
$$

Similar to the proof of Lemma 2.12 we can now use the anticommutation relations from Lemma 2.7 to rewrite this equation as

$$
\left(-(\mu+i t d) m+B^{-1}\left(-\left(-\mu^{*}+i t d^{*}\right) m\right) B\right) H_{t}=M_{B} f
$$

since $\{m, \mu+i t d\}=I,\left\{m,-\mu^{*}+i t d^{*}\right\}=-I$ and $I-B^{-1} I B=0$. Then apply $\left(-\mu^{*}+i t d^{*}\right) B$ to obtain

$$
\left(\mu^{*}-i t d^{*}\right) B(\mu+i t d)\left(m H_{t}\right)=\left(-\mu^{*}+i t d^{*}\right)\left(B N^{+}-N^{-} B\right) f,
$$

since $-\mu^{*}+i t d^{*}$ is nilpotent. Evaluating the scalar part of this equation, we get the desired identity.

To find the identity for $H_{t}^{1, \|}$, we use the expression for $T_{A}=\left.\hat{T}_{B}\right|_{\hat{\mathcal{H}}^{1}}$ from Definition 2.43. Multiplying the equation $\left(I+i t \hat{T}_{A}\right) H_{t}=f$ by $A N^{+}-N^{-} A$ yields

$$
a_{\| \|} H_{t}^{1, \|}-a_{00} H_{t}^{1,0} e_{0}-i t\left(A \nabla H_{t}^{1,0}+d^{*} \mu A H_{t}\right)=a_{\| \| \|} f^{1, \|}-a_{00} f_{t}^{1,0} e_{0} .
$$

Evaluating the tangential part of this equation gives the desired identity.
(ii) By rescaling, we see from (i) that it suffices to show that

$$
\begin{equation*}
\left(\int_{B\left(x, r_{0}\right)}|u(y)|^{q} d y\right)^{1 / q}+\left(\int_{B\left(x, r_{0}\right)}|\nabla u(y)|^{p} d y\right)^{1 / p} \lesssim M\left(|g|^{p}\right)^{1 / p}(x), \tag{2.13}
\end{equation*}
$$

for all $x \in \mathbf{R}^{n}$ and all $u$ and $g$ satisfying an equation

$$
\left[\begin{array}{ll}
1 & i \operatorname{div}]
\end{array}\left[\begin{array}{cc}
a_{\perp \perp}^{\prime} & a_{\perp \|}^{\prime} \\
a_{\| \perp}^{\prime} & a_{\| \|}^{\prime}
\end{array}\right]\left[\begin{array}{c}
1 \\
i \nabla
\end{array}\right] u=\left[\begin{array}{ll}
1 & i \operatorname{div}]
\end{array}\right]\left[\begin{array}{c}
a_{\perp \perp}^{\prime} g^{1,0} \\
-a_{\| \|}^{\prime} g^{1, \|}
\end{array}\right]\right.
$$

where $A^{\prime}$ is a matrix with same norm and accretivity constant as for $A$. Indeed, by rescaling we see that $u(x)=H_{t}^{1,0}(t x), g(x)=f(t x)$ and $A^{\prime}(x)=A(t x)$ satisfies this hypothesis. To prove (2.13), we use that the maps $g \mapsto u$ and $g \mapsto \nabla u$ have $L_{p}\left(\mathbf{R}^{n}\right) \rightarrow L_{q}\left(\mathbf{R}^{n}\right)$ and $L_{p}\left(\mathbf{R}^{n}\right) \rightarrow$ $L_{p}\left(\mathbf{R}^{n}\right)$ off-diagonal bounds, respectively, with exponent $M$ for any $M>0$, i.e. there exists $C_{M}<\infty$ such that

$$
\begin{equation*}
\|u\|_{L_{q}(E)}+\|\nabla u\|_{L_{p}(E)} \leqslant C_{M}\langle\operatorname{dist}(E, F)\rangle^{-M}\|g\|_{p} \tag{2.14}
\end{equation*}
$$

whenever $E, F \subset \mathbf{R}^{n}$ and $\operatorname{supp} g \subset F$. To see this, let $L:=\left[\begin{array}{ll}1 i \operatorname{div}] A^{\prime}[1 i \nabla\end{array}\right]^{t}$. Note that $L: W_{2}^{1}\left(\mathbf{R}^{n}\right) \rightarrow W_{2}^{-1}\left(\mathbf{R}^{n}\right)$ is an isomorphism and that $L: W_{p}^{1}\left(\mathbf{R}^{n}\right) \rightarrow W_{p}^{-1}\left(\mathbf{R}^{n}\right)$ is bounded. Then by the stability result of Šneĭberg [24], it follows that there exists $\varepsilon>0$ such that $L: W_{p}^{1}\left(\mathbf{R}^{n}\right) \rightarrow W_{p}^{-1}\left(\mathbf{R}^{n}\right)$ is an isomorphism when $|p-2|<\varepsilon$. We then fix $p_{0} \in(2-\varepsilon, 2)$ and use Sobolev's embedding theorem to see that

$$
\|u\|_{q_{0}}+\|\nabla u\|_{p_{0}} \lesssim\|u\|_{W_{p_{0}}^{1}} \lesssim\|L u\|_{W_{p_{0}}^{-1}} \lesssim\|g\|_{p_{0}} .
$$

By choosing $p_{0}$ close to 2 , we may assume that $q_{0}>2$. Thus we have bounded maps $g \mapsto$ $u: L_{p_{0}}(F) \rightarrow L_{q_{0}}(E)$ and $g \mapsto \nabla u: L_{p_{0}}(F) \rightarrow L_{p_{0}}(E)$, with norms $\leqslant C$. Also, by Proposition 2.25, the norms of $g \mapsto u: L_{2}(F) \rightarrow L_{2}(E)$ and $g \mapsto \nabla u: L_{2}(F) \rightarrow L_{2}(E)$ are bounded by $\langle\operatorname{dist}(E, F)\rangle^{-M_{0}}$. Interpolation now proves (2.14) for some $p_{0}<p<2,2<q<q_{0}$ and $M=M_{0}\left(1-p_{0} / p\right) /\left(1-p_{0} / 2\right)$.

Finally we show how (2.14) implies (2.13). Let $E=F_{0}:=B\left(x, r_{0}\right)$ and for $k \geqslant 1$ let $F_{k}:=$ $B\left(x, 2^{k} r_{0}\right) \backslash B\left(x, 2^{k-1} r_{0}\right)$. This gives

$$
\begin{aligned}
\left(\int_{B\left(x, r_{0}\right)}|u(y)|^{q} d y\right)^{1 / q}+\left(\int_{B\left(x, r_{0}\right)}|\nabla u(y)|^{p} d y\right)^{1 / p} & \lesssim \sum_{k=0}^{\infty} 2^{-M k}\left(\int_{F_{k}}|g(y)|^{p} d y\right)^{1 / p} \\
& \lesssim \sum_{k=0}^{\infty} 2^{(n / p-M) k}\left(\int_{B\left(x, 2^{k} r_{0}\right)}|g(y)|^{p} d y\right)^{1 / p} \\
& \lesssim\left(M\left(|g|^{p}\right)(x)\right)^{1 / p}
\end{aligned}
$$

provided we chose $M>n / p$.
Proof of Proposition 2.56. To prove that $\left\|\widetilde{N}_{*}(F)\right\| \gtrsim\|f\|$, we calculate

$$
\begin{aligned}
\left\|\tilde{N}_{*}(F)\right\|^{2} & \gtrsim \sup _{t>0} \int_{\mathbf{R}^{n}} f_{|y-x|<c_{1} t} f_{|s-t|<c_{0} t}|F(s, y)|^{2} d s d y d x \\
& =\sup _{t>0} f_{|s-t|<c_{0} t}\left\|F_{s}\right\|^{2} d s \gtrsim \sup _{t>0}\left\|F_{\left(1+c_{0}\right) t}\right\|^{2} \approx\|f\|^{2} .
\end{aligned}
$$

We have here used that $F_{\left(1+c_{0}\right) t}=e^{-\left(\left(1+c_{0}\right) t-s\right)\left|T_{A}\right|} F_{s}$, which by Proposition 2.32 shows that $\left\|F_{\left(1+c_{0}\right) t}\right\| \lesssim\left\|{\underset{\sim}{\sim}}_{s}\right\|$ when $|s-t|<c_{0} t$.

To prove $\left\|\widetilde{N}_{*}(F)\right\| \lesssim\|f\|$, we note that since $\operatorname{curl}_{t, x} F=0$, we can write $F=\nabla_{t, x} U$ for some scalar potential $U$, and we see that $U$ solves the second order equation (1.1). With notation $\widetilde{D}(t, x)=\left\{(s, y) ;|s-t|<\tilde{c}_{0} t,|y-x|<\tilde{c}_{1} t\right\}$, for constants $c_{0}<\tilde{c}_{0}<1$ and $c_{1}<\tilde{c}_{1}<\infty$, and $\bar{U}:=f f_{\widetilde{D}(t, x)} U(s, y) d s d y$, we have

$$
\begin{aligned}
\left(f_{D(t, x)}|t \nabla U(s, y)|^{2} d s d y\right)^{1 / 2} & \lesssim\left(\int_{\widetilde{D}(t, x)} f|U(s, y)-\bar{U}|^{2} d s d y\right)^{1 / 2} \\
& \lesssim\left(\underset{\widetilde{D}(t, x)}{ } f|t \nabla U(s, y)|^{p} d s d y\right)^{1 / p}
\end{aligned}
$$

with $2(n+1) /(n+3)<p<2$. The first estimate uses Caccioppoli's inequality and the second estimate uses Poincare's inequality. Thus it suffices to bound the $L_{2}$ norm of

$$
\sup _{t>0}\left({\underset{\widetilde{D}}{(t, x)}}_{f f}^{f}|F(s, y)|^{p} d s d y\right)^{1 / p} .
$$

To this end, we write $F_{s}=H_{s}+\psi_{s}\left(T_{A}\right) f$, where $H_{s}:=\left(I+i s T_{A}\right)^{-1} f$ and $\psi(z):=e^{-|z|}-$ $(1+i z)^{-1}$. Using the quadratic estimates, the second term has estimates

$$
\begin{aligned}
& \int_{\mathbf{R}^{n}} \sup _{t>0}\left(\int_{\tilde{D}(t, x)}^{f f}\left|\psi_{s}\left(T_{A}\right) f(y)\right|^{p} d s d y\right)^{2 / p} d x \\
& \quad \lesssim \int_{\mathbf{R}^{n}} \sup _{t>0}\left(\int_{|s-t|<\tilde{c}_{0} t} \int_{|y-x|<\tilde{c}_{1} t}\left|\psi_{s}\left(T_{A}\right) f(y)\right|^{2} \frac{d y d s}{t^{n+1}}\right) d x \\
& \quad \lesssim \int_{0}^{\infty} \int_{\mathbf{R}^{n}} \int_{|y-x|<\tilde{c}_{1} s /\left(1-\tilde{c}_{0}\right)}\left|\psi_{s}\left(T_{A}\right) f(y)\right|^{2} s^{-(n+1)} d y d x d s \\
& \quad \lesssim \int_{0}^{\infty}\left\|\psi_{s}\left(T_{A}\right) f\right\|^{2} \frac{d s}{s} \lesssim\|f\|^{2} .
\end{aligned}
$$

For the first term, we use Lemma 2.57(ii) and obtain

$$
\begin{aligned}
& \int_{\mathbf{R}^{n}} \sup _{t>0}({\left.\underset{\widetilde{D}(t, x)}{ } f f\left|H_{s}(y)\right|^{p} d s d y\right)^{2 / p} d x}^{\quad \lesssim \int_{\mathbf{R}^{n}} \sup _{t>0} \sup _{|s-t|<\tilde{c}_{0} t}(\underbrace{f}_{B\left(x, \tilde{c}_{1} s /\left(1-\tilde{c}_{0}\right)\right)}\left|H_{s}(y)\right|^{p} d y)^{2 / p} d x}
\end{aligned}
$$

$$
\lesssim\left\|M\left(|f|^{p}\right)\right\|_{2 / p}^{2 / p} \lesssim\left\||f|^{p}\right\|_{2 / p}^{2 / p}=\|f\|_{2}^{2}
$$

using the boundedness of the Hardy-Littlewood maximal function on $L_{2 / p}\left(\mathbf{R}^{n}\right)$. This completes the proof.

## 3. Invertibility of unperturbed operators

In this section we prove Theorems 1.4 and 1.1 for the unperturbed problem, i.e. for $B_{0}^{k}$ and $A_{0}$, respectively. We do this by verifying the hypothesis in Lemmas 2.51 and 2.53, i.e. we prove that $T_{B_{0}^{k}}$ satisfies quadratic estimates and $\lambda-E_{B_{0}^{k}} N_{B_{0}^{k}}$ is an isomorphism, and that $T_{A_{0}}$ satisfies quadratic estimates and $I \pm E_{A_{0}} N_{A_{0}}$ are isomorphisms, respectively.

### 3.1. Block coefficients

In this section we assume that $B=B_{0} \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}(\bigwedge)\right)$ have properties as in Definition 2.8 with the extra property that it is a block matrix, i.e.

$$
B=\left[\begin{array}{cc}
B_{\perp \perp} & 0 \\
0 & B_{\| \| \|}
\end{array}\right]
$$

in the splitting $\mathcal{H}=N^{-} \mathcal{H} \oplus N^{+} \mathcal{H}$. Note that $B$ being of this form is equivalent with the commutation relations $N^{ \pm} B=B N^{ \pm}$.

Lemma 3.1. Let B be a block matrix as above. Then

$$
T_{B}=\Gamma+B^{-1} \Gamma^{*} B
$$

where $\Gamma=N m d=-i N \underline{d}$ is a nilpotent first order, homogeneous partial differential operator with constant coefficients.

Proof. Since $N^{ \pm} B=B N^{ \pm}$, it follows that $M_{B}=N^{+}-B^{-1} N^{-} B=N^{+}-N^{-}=N$ and

$$
T_{B}=N\left(m d+B^{-1} m d^{*} B\right)=\Gamma+B^{-1} \Gamma^{*} B
$$

since $N^{2}=I$ and $N B^{-1}=B^{-1} N$. The operator $\Gamma$ is nilpotent since $\Gamma^{2}=N m d N m d=$ $N m N d m d=-N m N m d^{2}=0$.

Remark 3.2. Note that if $\Pi_{B}=\Gamma+\Gamma_{B}^{*}$, where $\Gamma_{B}^{*}=B^{-1} \Gamma^{*} B$ and $\Gamma$ is nilpotent, is an operator of the form considered in [10], then $\Pi_{B}$ intertwines $\Gamma$ and $\Gamma_{B}^{*}$ in the sense that $\Pi_{B} \Gamma u=\Gamma_{B}^{*} \Pi_{B} u$ for all $u \in \mathrm{D}\left(\Gamma_{B}^{*} \Pi_{B}\right)$ and $\Pi_{B} \Gamma_{B}^{*} u=\Gamma \Pi_{B} u$ for all $u \in \mathrm{D}\left(\Gamma \Pi_{B}\right)$. Thus $\Pi_{B}^{2}$ commutes with both $\Gamma$ and $\Gamma_{B}^{*}$ on appropriate domains. In particular, if $P_{t}^{B}=\left(1+t^{2} \Pi_{B}^{2}\right)^{-1}$ and $Q_{t}^{B}=t \Pi_{B}(1+$ $\left.t^{2} \Pi_{B}^{2}\right)^{-1}$, then we find that $\Gamma P_{t}^{B} u=P_{t}^{B} \Gamma u, \Gamma_{B}^{*} Q_{t}^{B} u=Q_{t}^{B} \Gamma u$ for all $u \in \mathrm{D}(\Gamma)$ and $\Gamma_{B}^{*} P_{t}^{B} u=$ $P_{t}^{B} \Gamma_{B}^{*} u, \Gamma Q_{t}^{B} u=Q_{t}^{B} \Gamma_{B}^{*} u$ for all $u \in \mathrm{D}\left(\Gamma_{B}^{*}\right)$.

Theorem 3.3. Let B be a block matrix as above. Then
(i) $T_{B}$ satisfies quadratic estimates, and
(ii) $E_{B} N_{B}+N_{B} E_{B}=0$ and $N_{B}=N$. In particular $\lambda-E_{B} N_{B}$ is an isomorphism whenever $\lambda \notin\{i,-i\}$ with

$$
\left(\lambda-E_{B} N_{B}\right)^{-1}=\frac{1}{\lambda^{2}+1}\left(\lambda-N_{B} E_{B}\right) .
$$

Proof. For operators of the form $\Gamma+B^{-1} \Gamma^{*} B$, quadratic estimates were proved in [10], with essentially the same methods as we use here in Section 4.1.

To prove (ii), note that since $B$ is a block matrix, it follows that $B^{-1} N^{ \pm} \mathcal{H}=N^{ \pm} \mathcal{H}$. Thus the projections $N_{B}^{ \pm}$associated with the splitting $\mathcal{H}=N^{-} \mathcal{H} \oplus B^{-1} N^{+} \mathcal{H}$ are $N^{ \pm}$and the associated reflection operator is $N_{B}=N$.

To prove invertibility of $\lambda-E_{B} N_{B}$, we note that

$$
T_{B} N_{B}=\left(\Gamma+B^{-1} \Gamma^{*} B\right) N=-N\left(\Gamma+B^{-1} \Gamma^{*} B\right)=-N_{B} T_{B}
$$

since $N$ commutes with $d, d^{*}$ and $B$, and anticommutes with $m$. Thus

$$
E_{B} N_{B}=\operatorname{sgn}\left(T_{B}\right) N_{B}=N_{B} \operatorname{sgn}\left(N_{B} T_{B} N_{B}\right)=N_{B} \operatorname{sgn}\left(-T_{B}\right)=-N_{B} E_{B},
$$

since $\operatorname{sgn}(z)$ is odd. Using this anticommutation formula we obtain

$$
\left(\lambda-N_{B} E_{B}\right)\left(\lambda-E_{B} N_{B}\right)=\lambda^{2}+1-\lambda\left(E_{B} N_{B}+N_{B} E_{B}\right)=\lambda^{2}+1,
$$

and similarly $\left(\lambda-E_{B} N_{B}\right)\left(\lambda-N_{B} E_{B}\right)=\lambda^{2}+1$, from which the stated formula for the inverse follows.

### 3.2. Constant coefficients

We here collect results in the case when $B(x)=B \in \mathcal{L}(\bigwedge)$ is a constant accretive matrix. In this case we make use of the Fourier transform

$$
\mathcal{F} u(\xi)=\hat{u}(\xi):=\frac{1}{2 \pi i} \int_{\mathbf{R}^{n}} u(x) e^{-i(x, \xi)} d x
$$

acting componentwise. If we let

$$
\left.\mu_{\xi} f(\xi):=\xi \wedge f(\xi), \quad \mu_{\xi}^{*} f(\xi):=\xi\right\lrcorner f(\xi)
$$

then $T_{B}$, conjugated with $\mathcal{F}$, is the multiplication operator

$$
M_{\xi} f(\xi):=M_{B}^{-1}\left(i m \mu_{\xi}-i B^{-1} m \mu_{\xi}^{*} B\right) f(\xi), \quad \xi \in \mathbf{R}^{n}
$$

Lemma 3.4. For all $t \in \mathbf{R}$ and $\xi \in \mathbf{R}^{n}$ we have

$$
\left|\left(i t+M_{\xi}\right)^{-1}\right| \approx\left(t^{2}+|\xi|^{2}\right)^{-1 / 2}
$$

Proof. Let $u=\left(i t+M_{\xi}\right) f$. It suffices to prove that $\|u\|^{2} \gtrsim\left(t^{2}+|\xi|^{2}\right)\|f\|^{2}$. With $\widetilde{B}:=m B m$, it follows from the definition of $M_{\xi}$ that

$$
m M_{B} u=\left(\Gamma+\widetilde{B}^{-1} \Gamma^{*} B\right) f
$$

where $\Gamma=i\left(t \mu+\mu_{\xi}\right)$. Using Lemma 2.7 we get

$$
\left\|\left(\Gamma+\Gamma^{*}\right) g\right\|^{2}=\|\Gamma g\|^{2}+\left\|\Gamma^{*} g\right\|^{2}=\left(\left(\Gamma^{*} \Gamma+\Gamma^{*} \Gamma\right) g, g\right)=\left(t^{2}+|\xi|^{2}\right)\|g\|^{2} .
$$

Therefore our estimate follows from Lemma 2.21(ii).
Proposition 3.5. If $B(x)=B \in \mathcal{L}(\bigwedge)$ is a constant, accretive matrix, then $T_{B}$ satisfies quadratic estimates.

Proof. Using the lemma, we obtain the estimate

$$
\left|\frac{t M_{\xi}}{1+t^{2} M_{\xi}^{2}}\right| \leqslant t\left|M_{\xi}\right|\left|\left(i-M_{t \xi}\right)^{-1}\right|\left|\left(i+M_{t \xi}\right)^{-1}\right| \lesssim \frac{t|\xi|}{1+t^{2}|\xi|^{2}}
$$

Thus using Plancherel's formula we obtain

$$
\begin{aligned}
\int_{0}^{\infty}\left\|\frac{t T_{B}}{1+t^{2} T_{B}^{2}} u\right\|^{2} \frac{d t}{t} & \approx \int_{0}^{\infty}\left\|\frac{t M_{\xi}}{1+t^{2} M_{\xi}^{2}} \hat{u}\right\|^{2} \frac{d t}{t} \\
& \lesssim \int_{\mathbf{R}^{n}}\left(\int_{0}^{\infty}\left(\frac{t|\xi|}{1+t^{2}|\xi|^{2}}\right)^{2} \frac{d t}{t}\right)|\hat{u}(\xi)|^{2} d \xi \approx\|u\|^{2}
\end{aligned}
$$

where the last step follows from a change of variables $s=t|\xi|$.
Next we prove that Neumann and Regularity problems are well posed in the case of a complex, constant, accretive matrix

$$
A=\left[\begin{array}{ll}
a_{00} & a_{0 \|} \\
a_{\| 0} & a_{\| \|}
\end{array}\right]
$$

For this it suffices to consider $B=I \oplus A \oplus I \oplus I \oplus \cdots \oplus I$ and the action of $T_{A}=\left.T_{B}\right|_{\hat{\mathcal{H}}^{1}}$ on the invariant subspace $\hat{\mathcal{H}}^{1}$. Recall that $f \in \hat{\mathcal{H}}^{1}$ means that $f$ is a vector field $f: \mathbf{R}^{n} \rightarrow \bigwedge^{1}=\mathbf{C}^{n+1}$ such that $d f_{\|}=0$. On the Fourier transform side $f \in \hat{\mathcal{H}}^{1}$ is seen to correspond to a vector field $\hat{f}: \mathbf{R}^{n} \rightarrow \bigwedge^{1}$ such that $\xi \wedge \hat{f}_{\|}=0$, i.e. $\hat{f}$ is such that its tangential part $f_{\|}$is a radial vector field. Thus the space

$$
\mathcal{F}\left(\hat{\mathcal{H}}^{1}\right)=\left\{\hat{f} \in L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}^{n+1}\right) ; \xi \wedge \hat{f}_{\|}=0\right\}
$$

can be identified with $L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}^{2}\right)$ if we use $\left\{e_{0}, \xi /|\xi|\right\}$ as basis for the two-dimensional space $\hat{\mathcal{H}}_{\xi}^{1}$ to which $\hat{f}(\xi)$ belongs. Furthermore, the operator $T_{A}$ is conjugated to the multiplication operator

$$
M_{\xi}=M_{A}^{-1}\left(i \mu^{*} \mu_{\xi}-i A^{-1} \mu \mu_{\xi}^{*} A\right): \mathcal{F}\left(\hat{\mathcal{H}}^{1}\right) \rightarrow \mathcal{F}\left(\hat{\mathcal{H}}^{1}\right)
$$

under the Fourier transform, where $M_{A}:=N^{+}-A^{-1} N^{-} A$. We see from the expression (1.5) for $T_{A}$ that, at a fixed point $\xi \in S^{n-1}$, the matrix for $M_{\xi}$ in the basis $\left\{e_{0}, \xi\right\}$ is

$$
\left[\begin{array}{cc}
\frac{i}{a_{00}}\left(a_{0 \|}+a_{\| 0}, \xi\right) & \frac{i}{a_{00}}\left(a_{\| \|} \xi, \xi\right) \\
-i & 0
\end{array}\right]
$$

Theorem 3.6. Let $A(x)=A \in \mathcal{L}(\bigwedge)$ be a constant, complex, accretive matrix. Then

$$
I \pm E_{A} N_{A}: \hat{\mathcal{H}}^{1} \rightarrow \hat{\mathcal{H}}^{1}
$$

are isomorphisms. In particular, Neumann and Regularity problems (Neu- $A$ ) and (Reg-A) are well posed.

Proof. From Remark 2.54 we see that it suffices to prove that all four restricted projections

$$
N_{A}^{ \pm}: E_{A}^{ \pm} \hat{\mathcal{H}}^{1} \rightarrow N_{A}^{ \pm} \hat{\mathcal{H}}^{1}
$$

are isomorphisms, or equivalently that the constant multiplication operators $N_{A}^{ \pm}$are isomorphisms on $\mathcal{F}\left(\hat{\mathcal{H}}^{1}\right)$. For the projections $\chi_{ \pm}\left(M_{\xi}\right)$ conjugated to the Hardy projection operators $E_{A}^{ \pm}$, we observe that $M_{t \xi}=t M_{\xi}$ and hence $\chi_{ \pm}\left(t M_{\xi}\right)=\chi_{ \pm}\left(M_{\xi}\right)$ for all $t>0$. Thus it suffices to verify that

$$
N_{A}^{ \pm}: \chi_{ \pm}\left(M_{\xi}\right) \hat{\mathcal{H}}_{\xi}^{1} \rightarrow N_{A}^{ \pm} \hat{\mathcal{H}}_{\xi}^{1}
$$

are isomorphisms for each $\xi \in S^{n-1}$. For such fixed $\xi$, using the basis $\left\{e_{0}, \xi\right\}$ from above, $N_{A}^{-} \hat{\mathcal{H}}_{\xi}^{1}=\left\{\hat{f} \in \hat{\mathcal{H}}_{\xi}^{1} ; e_{0} \wedge \hat{f}=0\right\}$ is spanned by $[10]^{t}$ and $N_{A}^{+} \hat{\mathcal{H}}_{\xi}^{1}=\left\{\hat{f} \in \hat{\mathcal{H}}_{\xi}^{1} ;\left(A \hat{f}, e_{0}\right)=0\right\}$ is spanned by $\left[\left(a_{0 \|}, \xi\right)-a_{00}\right]^{t}$. Indeed $\left(A\left(z e_{0}+w \xi\right), e_{0}\right)=z a_{00}+w\left(a_{0 \|}, \xi\right)$.

If we call $e_{\xi}^{+}$and $e_{\xi}^{-}$the two eigenvectors of $M_{\xi}$, it follows that $\chi_{ \pm}\left(M_{\xi}\right) \hat{\mathcal{H}}_{\xi}^{1}$ are spanned by these, as these subspaces are one-dimensional. It suffices to show that $\left[\begin{array}{ll}1 & 0\end{array}\right]^{t}$ and $\left[\left(a_{0 \|}, \xi\right)-a_{00}\right]^{t}$ are not eigenvectors of $M_{\xi}$. We have

$$
\begin{aligned}
M_{\xi}\left[\begin{array}{l}
1 \\
0
\end{array}\right] & =i\left[\begin{array}{c}
\frac{1}{a_{00}}\left(a_{0 \|}+a_{\| 0}, \xi\right) \\
-1
\end{array}\right] \\
M_{\xi}\left[\begin{array}{c}
\left(a_{0 \|}, \xi\right) \\
-a_{00}
\end{array}\right] & =i\left[\begin{array}{c}
\frac{1}{a_{00}}\left(a_{0 \|}+a_{\| 0}, \xi\right)\left(a_{0 \|}, \xi\right)-\left(a_{\| \|} \xi, \xi\right) \\
-\left(a_{0 \|}, \xi\right)
\end{array}\right] .
\end{aligned}
$$

Clearly the normal vector is not an eigenvector. To prove that the second is not an eigenvector, note that the cross product is

$$
-\left(a_{0 \|}, \xi\right)^{2}+\left(a_{0 \|}+a_{\| 0}, \xi\right)\left(a_{0 \|}, \xi\right)-a_{00}\left(a_{\| \| \|} \xi, \xi\right)=\left(a_{\| 0}, \xi\right)\left(a_{0 \|}, \xi\right)-a_{00}\left(a_{\| \|} \xi, \xi\right)
$$

The right-hand side is non-zero since

$$
\left(\left[\begin{array}{ll}
a_{00} & a_{0 \|} \\
a_{\| 0} & a_{\| \|}
\end{array}\right]\left[\begin{array}{c}
z \\
w \xi
\end{array}\right],\left[\begin{array}{c}
z \\
w \xi
\end{array}\right]\right)=\left(\left[\begin{array}{cc}
a_{00} & \left(a_{0 \|}, \xi\right) \\
\left(a_{\| 0}, \xi\right) & \left(a_{\| \|} \xi, \xi\right)
\end{array}\right]\left[\begin{array}{c}
z \\
w
\end{array}\right],\left[\begin{array}{c}
z \\
w
\end{array}\right]\right)
$$

is a non-degenerate quadratic form as $A$ is accretive.
Remark 3.7. We note that the method above also can be used to show that $I \pm E_{A} N: \hat{\mathcal{H}}^{1} \rightarrow \hat{\mathcal{H}}^{1}$, with the unperturbed operator $N$, are isomorphisms when $A$ is constant. Here we also need to observe that the tangential vector [ $\left.\begin{array}{ll}0 & 1\end{array}\right]^{\mathrm{t}}$ is not an eigenvector to $M_{\xi}$.

### 3.3. Real symmetric coefficients

In this section, we assume that $B^{*}=B$. We first prove a Rellich type estimate.
Proposition 3.8. Assume that $B^{*}=B$ and that $f \in E_{B}^{+} \mathcal{V}_{B}$ or $f \in E_{B}^{-} \mathcal{V}_{B}$. Then

$$
\left.\left.(B f, f)=2 \operatorname{Re}\left(e_{0}\right\lrcorner(B f), e_{0}\right\lrcorner f\right)=2 \operatorname{Re}\left(e_{0} \wedge(B f), e_{0} \wedge f\right)
$$

In particular $\|f\| \approx\left\|\hat{N}_{B}^{-} f\right\| \approx\left\|\check{N}_{B}^{-} f\right\| \approx\left\|\check{N}_{B}^{+} f\right\| \approx\left\|\hat{N}_{B}^{+} f\right\|$.
Proof. It suffices to consider $f \in E_{B}^{+} \mathcal{V}_{B}$ as the case $f \in E_{B}^{-} \mathcal{V}_{B}$ is treated similarly. We use Lemma 2.34 and write $F_{t}:=e^{-t\left|T_{B}\right|} f$. Hence $(0, \infty) \ni t \mapsto F_{t} \in \mathcal{H}$ is differentiable, $\lim _{t \rightarrow 0} F_{t}=f$ and $\lim _{t \rightarrow \infty} F_{t}=0$. Furthermore $F_{t} \in \mathrm{D}\left(d_{x}\right)$ and $B F_{t} \in \mathrm{D}\left(d_{x}^{*}\right)$ for all $t \in(0, \infty)$. We note the formulae

$$
\begin{aligned}
d_{t, x}^{*} m G_{t}+m d_{t, x}^{*} G_{t} & =-\partial_{t} G_{t}, \\
B m d_{t, x} F_{t} & =-m d_{t, x}^{*} B F_{t} .
\end{aligned}
$$

The first identity, which we apply with $G_{t}=B F_{t}$, follows from Lemma 2.7, whereas the second is equivalent to $\partial_{t} F_{t}+T_{B} F_{t}=0$, and follows from Lemma 2.34. We get

$$
\begin{aligned}
(B f, f) & =-\int_{0}^{\infty}\left(\partial_{t} B F_{t}, F_{t}\right)+\left(B F_{t}, \partial_{t} F_{t}\right)=-2 \operatorname{Re} \int_{0}^{\infty}\left(\partial_{t} B F_{t}, F_{t}\right) \\
& =2 \operatorname{Re} \int_{0}^{\infty}\left(d_{t, x}^{*} m B F_{t}+m d_{t, x}^{*} B F_{t}, F_{t}\right)=2 \operatorname{Re} \int_{0}^{\infty}\left(d_{t, x}^{*} m B F_{t}-B m d_{t, x} F_{t}, F_{t}\right) \\
& =2 \operatorname{Re} \int_{0}^{\infty}\left(d_{t, x}^{*} m B F_{t}, F_{t}\right)-\left(m B F_{t}, d_{t, x} F_{t}\right) \\
& =2 \operatorname{Re} \int_{0}^{\infty}\left(\left(d^{*}-\mu^{*} \partial_{t}\right) m B F_{t}, F_{t}\right)-\left(m B F_{t},\left(d+\mu \partial_{t}\right) F_{t}\right)
\end{aligned}
$$

$$
\begin{aligned}
& =-2 \operatorname{Re} \int_{0}^{\infty}\left(\partial_{t} m B F_{t}, \mu F_{t}\right)+\left(m B F_{t}, \partial_{t} \mu F_{t}\right) \\
& =2 \operatorname{Re}(m B f, \mu f)=2 \operatorname{Re}\left(e_{0} \wedge(B f), e_{0} \wedge f\right)
\end{aligned}
$$

Note that all integrals are convergent since $\left\|F_{t}\right\| \lesssim \min \left(1, t^{-s}\right)$ and since $\partial_{t} F_{t}=-T_{B} F_{t}$ with $\left\|T_{B} F_{t}\right\| \lesssim \min \left(t^{s-1}, t^{-1}\right)$ if $f \in \mathrm{D}\left(\left|T_{B}\right|^{s}\right) \cap \mathrm{R}\left(\left|T_{B}\right|^{-s}\right)$. This follows as in the proof of Lemma 2.34. Furthermore, we note that

$$
\begin{aligned}
\left(e_{0} \wedge B f, e_{0} \wedge f\right) & \left.=\left(B f, e_{0}\right\lrcorner\left(e_{0} \wedge f\right)\right) \\
& \left.\left.\left.=\left(B f, f-e_{0} \wedge\left(e_{0}\right\lrcorner f\right)\right)=(B f, f)-\left(e_{0}\right\lrcorner(B f), e_{0}\right\lrcorner f\right)
\end{aligned}
$$

Together with the calculation above this proves that $\left.\left.(B f, f)=2 \operatorname{Re}\left(e_{0}\right\lrcorner(B f), e_{0}\right\lrcorner f\right)$.
To prove that $\|f\| \approx\left\|\hat{N}_{B}^{+} f\right\| \approx\left\|\hat{N}_{B}^{-} f\right\|$, it suffices to show that $\|f\| \lesssim\left\|\hat{N}_{B}^{ \pm} f\right\|$ since $\hat{N}_{B}^{ \pm}$ are bounded. From the Rellich type identities above we have $\left.\|f\|^{2} \lesssim \| e_{0}\right\lrcorner(B f)\left\|\| e_{0}\right\lrcorner f \| \lesssim$ $\left\|\hat{N}_{B}^{-} f\right\|\|f\|$ which proves $\|f\| \lesssim\left\|\hat{N}_{B}^{-} f\right\|$, and using the other identity we obtain $\|f\|^{2} \lesssim \| e_{0} \wedge$ $(B f)\left\|\left\|e_{0} \wedge f\right\| \lesssim\right\| f\left\|\left\|\hat{N}_{B}^{+} f\right\|\right.$. The proof of $\| f\|\approx\| \check{N}_{B}^{+} f\|\approx\| \check{N}_{B}^{-} f \|$ is similar.

We note that Proposition 3.8 also proves that the norms of the two components of $f \in E_{B}^{ \pm} \mathcal{V}_{B}$ in the unperturbed splitting $\mathcal{H}=N^{-} \mathcal{H} \oplus N^{+} \mathcal{H}$ are comparable.

Corollary 3.9. Assume that $B^{*}=B$ and that $f \in E_{B}^{+} \mathcal{V}_{B}$ or $f \in E_{B}^{-} \mathcal{V}_{B}$, and decompose $B$ in the splitting $\mathcal{H}=N^{-} \mathcal{H} \oplus N^{+} \mathcal{H}$ as

$$
B=\left[\begin{array}{cc}
B_{\perp \perp} & B_{\perp_{\|}} \\
B_{\| \perp} & B_{\| \|}
\end{array}\right]
$$

If $f_{\|}=N^{+} f$ and $f_{\perp}=N^{-} f$, then

$$
\left(B_{\perp \perp} f_{\perp}, f_{\perp}\right)=\left(B_{\| \|} f_{\|}, f_{\|}\right)
$$

In particular $\|f\| \approx\left\|N^{+} f\right\| \approx\left\|N^{-} f\right\|$.
Proof. We obtain from Proposition 3.8 that $\operatorname{Re}\left(B_{\perp \perp} f_{\perp}+B_{\perp \|} f_{\|}, f_{\perp}\right)=\operatorname{Re}\left(B_{\| \perp} f_{\perp}+B_{\| \|} f_{\|}, f_{\|}\right)$. The corollary now follows since $B_{\perp \perp}, B_{\| \|} \geqslant \kappa_{B}$ and $B_{\perp \|}^{*}=B_{\| \perp}$.

Next we turn to quadratic estimates for the operator $T_{B}$. As before we assume that $B$ is as in Definition 2.8 and that $B=B^{*}$. For the rest of this section we shall also assume that

$$
\begin{equation*}
B_{\| \perp}^{j}=B_{\perp \|}^{j}=0, \quad \text { for } j \geqslant 2 \tag{3.1}
\end{equation*}
$$

Note that in particular this is true if $B=I \oplus A \oplus I \oplus \cdots \oplus I$, where $A \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\bigwedge^{1}\right)\right)$.
To prove quadratic estimates, we shall use Proposition 2.36 where we verify the hypothesis separately on the subspaces $\hat{\mathcal{H}}_{B}^{k}$ and $\check{\mathcal{H}}_{B}$, which is possible by Lemma 2.46.

Lemma 3.10. Assume that $B^{*}=B$ and (3.1) holds. Then

$$
\|f\| \lesssim\left\|t \partial_{t} F_{t}\right\|_{ \pm}, \quad \text { for all } f \in E_{B}^{ \pm} \check{\mathcal{V}}_{B}
$$

where $F_{t}=e^{\mp t\left|T_{B}\right|} f$.
Proof. As the two estimates are similar, we only show the estimate for $f \in E_{B}^{+} \breve{\mathcal{V}}_{B}$. Note that $F(t, x)$ satisfies (2.12), i.e.

$$
\left\{\begin{array}{l}
\left.d F=e_{0}\right\lrcorner \partial_{t} F, \\
d^{*}(B F)=-e_{0} \wedge\left(B \partial_{t} F\right) .
\end{array}\right.
$$

Splitting both sides of both equations into normal and tangential parts, with notation as in Corollary 3.9, we obtain

$$
\begin{align*}
d F_{\|} & =m \partial_{t} F_{\perp},  \tag{3.2}\\
d F_{\perp} & =0,  \tag{3.3}\\
d^{*}\left(B_{\perp \perp} F_{\perp}+B_{\perp \|} F_{\|}\right) & =-m\left(B_{\| \perp} \partial_{t} F_{\perp}+B_{\| \|} \partial_{t} F_{\|}\right),  \tag{3.4}\\
d^{*}\left(B_{\| \perp} F_{\perp}+B_{\| \|} F_{\|}\right) & =0 . \tag{3.5}
\end{align*}
$$

We have here used that $\left.e_{0}\right\lrcorner f=m f$ if $f$ is normal, and $e_{0} \wedge f=m f$ if $f$ is tangential. A key observation is that the first term on the right-hand side in (3.4) vanishes since $B_{\| \perp} F_{\perp}=0$. To see this, note that

$$
B_{\| \perp} F_{\perp}=B_{\| \perp}^{1} F_{\perp}^{1}+B_{\| \perp}^{2} F_{\perp}^{2}+\cdots+B_{\| \perp}^{n+1} F_{\perp}^{n+1} .
$$

By hypothesis (3.1), $B_{\| \perp}^{2}=\cdots=B_{\| \perp}^{n+1}=0$. Furthermore, writing $F_{\perp}^{1}=F_{0} e_{0}$ with the function $F_{0}$ being scalar, we get from (3.3) that $0=d\left(F_{0} e_{0}\right)=\left(\nabla F_{0}\right) \wedge e_{0}$, so $F_{0}$ is constant and therefore vanishes, and thus so does $F_{\perp}^{1}$. Eq. (3.4) reduces to

$$
\begin{equation*}
B_{\| \|} \partial_{t} F_{\|}=-m d^{*}\left(B_{\perp \perp} F_{\perp}+B_{\perp \|} F_{\|}\right) . \tag{3.6}
\end{equation*}
$$

We calculate

$$
\begin{aligned}
& \|f\|^{2} \lesssim\left(B_{\| \|} f_{\|}, f_{\|}\right)=-2 \operatorname{Re} \int_{0}^{\infty}\left(B_{\| \|} \partial_{t} F_{\|}, F_{\|}\right) d t=2 \operatorname{Re} \int_{0}^{\infty}\left(d^{*}\left(B_{\perp \perp} F_{\perp}+B_{\perp \|} F_{\|}\right), m F_{\|}\right) d t \\
& \quad=-2 \operatorname{Re} \int_{0}^{\infty}\left(\left(d^{*}\left(B_{\perp \perp} \partial_{t} F_{\perp}+B_{\perp \|} \partial_{t} F_{\|}\right), m F_{\|}\right)+\left(d^{*}\left(B_{\perp \perp} F_{\perp}+B_{\perp \|} F_{\| \|}\right), m \partial_{t} F_{\|}\right)\right) t d t \\
& \quad=2 \operatorname{Re} \int_{0}^{\infty}\left(\left(B_{\perp \perp} \partial_{t} F_{\perp}+B_{\perp \|} \partial_{t} F_{\|}, \partial_{t} F_{\perp}\right)+\left(m B_{\| \|} \partial_{t} F_{\|}, m \partial_{t} F_{\| \|}\right)\right) t d t
\end{aligned}
$$

$$
=2 \operatorname{Re} \int_{0}^{\infty}\left(B \partial_{t} F, \partial_{t} F\right) t d t \approx\left\|t \partial_{t} F\right\|^{2}
$$

Here in the third step we use (3.6), in the fourth step we integrate by parts, in the fifth step we use duality and that $d m F_{\|}=-m d F_{\|}=-\partial_{t} F_{\perp}$ by (3.2) for the first term and again (3.6) for the second term. Finally in step 6 we use that $m$ is an isometry and that $B_{\| \perp} F_{\perp}=0$.

Next we turn to the subspace $\hat{\mathcal{H}}^{1}$.
Lemma 3.11. Let $A \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\bigwedge^{1}\right)\right)$ be real symmetric and $A \geqslant \kappa>0$, and let $B=I \oplus A \oplus$ $I \oplus \cdots \oplus I$. Then

$$
\|f\| \lesssim\left\|t \partial_{t} F\right\|_{ \pm}, \quad \text { for all } f \in E_{B}^{ \pm} \hat{\mathcal{V}}_{B}^{1}
$$

where $F=e^{\mp t\left|T_{B}\right|} f$.
Proof. Recall that if $f \in E_{B}^{ \pm} \hat{\mathcal{V}}_{B}^{1}$ and $F=F_{0} e_{0}+F_{\|}=e^{\mp t\left|T_{B}\right|} f$, then Lemma 2.12 shows that $F_{0}$ satisfies the equation

$$
\operatorname{div}_{t, x} A(x) \nabla_{t, x} F_{0}(t, x)=0
$$

Therefore, by the square function estimate of Dahlberg, Jerison and Kenig [13] and the estimates of harmonic measure of Jerison and Kenig [17], we have estimates $\left\|f_{0}\right\| \lesssim\left\|t \nabla_{t, x} F_{0}\right\|_{ \pm}$. Hence applying the Rellich estimates in Proposition 3.8, we obtain

$$
\|f\| \approx\left\|f_{0}\right\| \lesssim\left\|t \nabla_{t, x} F_{0}\right\|_{ \pm} \lesssim\left\|t t \partial_{t} F\right\|_{ \pm}
$$

since $\partial_{i} F_{0}=\partial_{0} F_{i}$, as $d_{t, x} F=0$.
We are now in position to prove the main result of this section.
Theorem 3.12. Let $A \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\bigwedge^{1}\right)\right)$ be real symmetric and $A \geqslant \kappa>0$, and let $B=I \oplus A \oplus$ $I \oplus \cdots \oplus I$. Then $T_{B}$ has quadratic estimates, so that in particular $E_{B}=\operatorname{sgn}\left(T_{B}\right): \mathcal{H} \rightarrow \mathcal{H}$ is bounded. Furthermore we have isomorphisms

$$
I \pm E_{B} N_{B}: \mathcal{H} \rightarrow \mathcal{H}
$$

Proof. To prove that $T_{B}$ has quadratic estimates it suffices by Proposition 2.36 to prove

$$
\|f\| \lesssim\left\|t \partial_{t} F\right\|_{ \pm}, \quad \text { for all } f \in E_{B}^{ \pm} \mathcal{V}_{B}
$$

To this end, we split

$$
f=f_{0}+f_{1}+\cdots+f_{n+1}+\check{f}
$$

where $f_{j} \in \hat{\mathcal{V}}_{B}^{j}$ and $\check{f} \in \check{\mathcal{V}}_{B}$. Lemma 2.46 shows that $e^{\mp t\left|T_{B}\right|}$ preserves these subspaces so that similarly

$$
F=F_{0}+F_{1}+\cdots+F_{n+1}+\check{F}
$$

where $\left(F_{j}\right)_{t} \in \hat{\mathcal{V}}_{B}^{j}$ and $\check{F}_{t} \in \check{\mathcal{V}}_{B}$. It thus suffices to prove that $\left\|f_{j}\right\| \lesssim\left\|t \partial_{t} F_{j}\right\|_{ \pm}, j=0,1, \ldots$, $n+1$, and $\|\check{f}\| \lesssim\left\|t \partial_{t} \check{F}\right\|_{ \pm}$. Lemma 3.10 shows that $\|\check{f}\| \lesssim\left\|t \partial_{t} \check{F}\right\|_{ \pm}$and Lemma 3.11 shows that $\left\|f_{1}\right\| \lesssim\left\|t \partial_{t} F_{1}\right\|_{ \pm}$. Furthermore, for $j \neq 1$ we observe that $T_{B}=\hat{T}_{B}=T_{I}$ on the subspace $\hat{\mathcal{V}}_{B}^{j}$, where $I$ denotes the identity matrix. Thus $F_{j}=e^{\mp t\left|T_{I}\right|} f_{j}$ and it follows from Proposition 3.5 that $\left\|f_{j}\right\| \lesssim\left\|t \partial_{t} F_{j}\right\|_{ \pm}$.

Applying Propositions 2.36 and 2.32 now shows that $E_{B}=\operatorname{sgn}\left(T_{B}\right): \mathcal{H} \rightarrow \mathcal{H}$ is a bounded operator. To show that $I \pm E_{B} \hat{N}_{B}: \mathcal{H} \rightarrow \mathcal{H}$ is invertible, note that the adjoint with respect to the duality $\langle\cdot, \cdot\rangle_{B}$ from Definition 2.16 is

$$
I \mp \check{N}_{B} E_{B}=\check{N}_{B}\left(I \mp E_{B} \check{N}_{B}\right) \check{N}_{B}
$$

according to Proposition 2.17. Having established the boundedness of $E_{B}$, the Rellich estimates clearly extends to a priori estimates for all eight restricted projections

$$
\hat{N}_{B}^{ \pm}: E_{B}^{ \pm} \mathcal{H} \rightarrow \hat{N}_{B}^{ \pm} \mathcal{H}, \quad \check{N}_{B}^{ \pm}: E_{B}^{ \pm} \mathcal{H} \rightarrow \check{N}_{B}^{ \pm} \mathcal{H}
$$

As in Remark 2.54 this translates to a priori estimates for $I \pm E_{B} \hat{N}_{B}$ and their adjoints, which proves that they are isomorphisms as in Remark 2.18.

Remark 3.13. Using instead Corollary 3.9, we also prove as in Theorem 3.12 that $I \pm E_{B} N$ : $\mathcal{H} \rightarrow \mathcal{H}$, where $N$ is the unperturbed operator, are isomorphisms.

## 4. Quadratic estimates for perturbed operators

In Section 3 we proved that $T_{B_{0}}$ satisfies quadratic estimates

$$
\begin{equation*}
\left\|Q_{t}^{B_{0}} f\right\| \approx\|f\|, \quad f \in \mathcal{H} \tag{4.1}
\end{equation*}
$$

for certain unperturbed coefficients $B_{0}$ :
(b) Block coefficients

$$
B_{0}=\left[\begin{array}{cc}
\left(B_{0}\right)_{\perp \perp} & 0 \\
0 & \left(B_{0}\right)_{\| \|}
\end{array}\right]
$$

(c) Constant coefficients $B_{0}(x)=B_{0}, x \in \mathbf{R}^{n}$, of the form $B_{0}=I \oplus A_{0} \oplus I \oplus \cdots \oplus I$, i.e. $B_{0}$ only acts non-trivially on the vector part.
(s) Real symmetric coefficients of the form $B_{0}=I \oplus A_{0} \oplus I \oplus \cdots \oplus I$.

Note that for (c), we did prove quadratic estimates for general constant coefficients $B_{0}$ in Proposition 3.5. However, since we only prove invertibility $I \pm E_{B_{0}} N_{B_{0}}$ on the subspace $\hat{\mathcal{H}}^{1}$
in Theorem 3.6, we shall only prove perturbation results for constant coefficients of the form $B_{0}=I \oplus A_{0} \oplus I \oplus \cdots \oplus I$.

In this section we let $B_{0} \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}(\bigwedge)\right)$ be a fixed accretive coefficient matrix with properties (b), (c) or (s). Constants $C$ in estimates or implicit in notation $\lesssim$ and $\approx$ in this section will be allowed to depend only on $\left\|B_{0}\right\|_{\infty}, \kappa_{B_{0}}$ and dimension $n$. Note that this is indeed the case for the constants implicit in (4.1).

We now consider a small perturbation $B \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}(\bigwedge)\right)$ of $B_{0}$ :

$$
\left\|B-B_{0}\right\|_{L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}(\wedge)\right)} \leqslant \varepsilon_{0}
$$

Throughout this section, we assume in particular that $\varepsilon_{0}$ is chosen small enough so that $B$ has properties as in Definition 2.8 with $\|B\|_{\infty} \leqslant 2\left\|B_{0}\right\|_{\infty}$ and $\kappa_{B} \geqslant \frac{1}{2} \kappa_{B_{0}}$. The goal is to show that

$$
\begin{equation*}
\left\|Q_{t}^{B} f\right\|\|\leqslant C\| f \|, \quad \text { whenever }\left\|B-B_{0}\right\|_{\infty} \leqslant \varepsilon, f \in \mathcal{H} \tag{4.2}
\end{equation*}
$$

where $C=C\left(\left\|B_{0}\right\|_{\infty}, \kappa_{B_{0}}, n\right)$ and $\varepsilon=\varepsilon\left(\left\|B_{0}\right\|_{\infty}, \kappa_{B_{0}}, n\right) \leqslant \varepsilon_{0}$. Since the properties of $B_{0}$ are stable under taking adjoints, we see that the quadratic estimate (2.7) follows from (4.2) and Lemma 2.35. In order to prove (4.2) we make use of the following identity, where we note that each of the six terms consists of three factors, the first being one of the operators from Corollary 2.23 , the second being a multiplication operator $\mathcal{E}$ with norm $\|\mathcal{E}\|_{\infty} \lesssim \varepsilon_{0}$ and the last factor being one of the operators from Corollary 2.23 but with $B$ replaced by $B_{0}$ :

$$
\begin{align*}
Q_{t}^{B}-Q_{t}^{B_{0}}= & P_{t}^{B}\left(\left(M_{B}^{-1}-M_{B_{0}}^{-1}\right) M_{B_{0}}\right) Q_{t}^{B_{0}}  \tag{4.3}\\
& -i P_{t}^{B}\left(M_{B}^{-1}\left(B^{-1}-B_{0}^{-1}\right) B_{0}\right)\left(t \underline{d}_{B_{0}}^{*} P_{t}^{B_{0}}\right)  \tag{4.4}\\
& -i\left(P_{t}^{B} M_{B}^{-1} t \underline{d}_{B}^{*}\right)\left(B^{-1}\left(B-B_{0}\right)\right) P_{t}^{B_{0}}  \tag{4.5}\\
& -Q_{t}^{B}\left(\left(M_{B}^{-1}-M_{B_{0}}^{-1}\right) M_{B_{0}}\right)\left(t T_{B_{0}} Q_{t}^{B_{0}}\right)  \tag{4.6}\\
& +i Q_{t}^{B}\left(M_{B}^{-1}\left(B^{-1}-B_{0}^{-1}\right) B_{0}\right)\left(t \underline{d}_{B_{0}}^{*} Q_{t}^{B_{0}}\right)  \tag{4.7}\\
& +i\left(Q_{t}^{B} M_{B}^{-1} t \underline{d}_{B}^{*}\right)\left(B^{-1}\left(B-B_{0}\right)\right) Q_{t}^{B_{0}} \tag{4.8}
\end{align*}
$$

Recall that

$$
T_{B}=-i M_{B}^{-1}\left(\underline{d}+B^{-1} \underline{d}^{*} B\right), \quad \text { where } \underline{d}=i m d, M_{B}=N^{+}-B^{-1} N^{-} B
$$

$P_{t}^{B}=\left(1+t^{2} T_{B}^{2}\right)^{-1}$ and $Q_{t}^{B}=t T_{B}\left(1+t^{2} T_{B}^{2}\right)^{-1}$ and similarly for $B_{0}$. The identity (4.3)-(4.8) is established by using

$$
Y\left(1+Y^{2}\right)^{-1}-X\left(1+X^{2}\right)^{-1}=\left(1+Y^{2}\right)^{-1}((Y-X)-Y(Y-X) X)\left(1+X^{2}\right)^{-1}
$$

with $Y=t T_{B}$ and $X=t T_{B_{0}}$, and then inserting

$$
\begin{aligned}
Y-X= & T_{B}-T_{B_{0}} \\
= & \left(M_{B}^{-1}-M_{B_{0}}^{-1}\right) M_{B_{0}} t T_{B_{0}}-i\left(M_{B}^{-1}\left(B^{-1}-B_{0}^{-1}\right) B_{0}\right)\left(t \underline{d}_{B_{0}}^{*}\right) \\
& -i\left(M_{B}^{-1} t \underline{d}_{B}^{*}\right)\left(B^{-1}\left(B-B_{0}\right)\right) .
\end{aligned}
$$

Lemma 4.1. Assume that for all $\mathcal{E} \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}(\bigwedge)\right)$ and all operator families $\left(\widetilde{Q}_{t}\right)_{t>0}$ with $L_{2}$ off-diagonal bounds $\left\|\widetilde{Q}_{t}\right\|_{\text {off }} \leqslant C$ as in Definition 2.24 , we have the two estimates

$$
\begin{gather*}
\left\|\widetilde{Q}_{t} \mathcal{E}\left(t \underline{d} Q_{t}^{B_{0}}\right)\right\|\left\|_{\mathrm{op}} \leqslant C\left(\left\|\widetilde{Q}_{t}\right\|_{\mathrm{op}}+1\right)\right\| \mathcal{E} \|_{\infty},  \tag{4.9}\\
\left\|\widetilde{Q}_{t} \mathcal{E}\left(t \underline{d}_{B_{0}}^{*} Q_{t}^{B_{0}}\right)\right\|\left\|_{\mathrm{op}} \leqslant C\left(\left\|\widetilde{Q}_{t}\right\|_{\mathrm{op}}+1\right)\right\| \mathcal{E} \|_{\infty} . \tag{4.10}
\end{gather*}
$$

Then there exists $\varepsilon>0$ such that (4.2) holds.
Proof. First consider the terms (4.3), (4.4) and (4.8). Using the uniform boundedness of $P_{t}^{B}$ and $Q_{t}^{B} M_{B}^{-1} t \underline{d}_{B}^{*}$ from Corollary 2.23 , we deduce that

$$
\begin{aligned}
\|(4.3)\|_{\mathrm{op}} & \lesssim\left\|B-B_{0}\right\|_{\infty}\left\|Q_{t}^{B_{0}}\right\|_{\mathrm{op}}, \\
\|(4.4)\|_{\mathrm{op}} & \lesssim\left\|B-B_{0}\right\|_{\infty}\left\|t \underline{d}_{B_{0}}^{*} P_{t}^{B_{0}}\right\|_{\mathrm{op}}, \\
\|(4.8)\|_{\mathrm{op}} & \lesssim B-B_{0}\left\|_{\infty}\right\| Q_{t}^{B_{0}} \|_{\mathrm{op}} .
\end{aligned}
$$

Observe that $t \underline{d}_{B_{0}}^{*} P_{t}^{B_{0}}=i \mathbf{P}_{B_{0}}^{2} M_{B_{0}} Q_{t}^{B_{0}}$ as in Corollary 2.23, where the Hodge projection $\mathbf{P}_{B_{0}}^{2}$ and $M_{B_{0}}$ are bounded. Thus we see from (4.1) that

$$
\|(4.3)\|_{\mathrm{op}}+\|(4.4)\|_{\mathrm{op}}+\|(4.8)\|_{\mathrm{op}} \lesssim\left\|B-B_{0}\right\|_{\infty}
$$

To handle the terms (4.5)-(4.7) we introduce the truncated operator families $\widetilde{Q}_{t}^{1}:=\chi(t) Q_{t}^{B}$ and $\widetilde{Q}_{t}^{2}:=\chi(t) P_{t}^{B} M_{B}^{-1} t \underline{d}_{B}^{*}$, where $\chi(t)$ denotes the characteristic function of the interval $\left[\tau^{-1}, \tau\right]$ for some large $\tau$. Note that $P_{t}^{B} M_{B}^{-1} t \underline{d}_{B}^{*}=i Q_{t}^{B} \mathbf{P}_{B}^{1}$ as in Corollary 2.23, and therefore $\left\|\left\|\widetilde{Q}_{t}^{2}\right\|_{\mathrm{op}}=\right.$ $\left\|\left\|\widetilde{Q}_{t}^{1} \mathbf{P}_{B}^{1}\right\|_{\mathrm{op}} \lesssim\right\| \widetilde{Q}_{t}^{1} \|_{\mathrm{op}}$. To use the hypothesis on the terms (4.5)-(4.7) we note that the last factors are

$$
\begin{aligned}
P_{t}^{B_{0}} & =I+i M_{B_{0}}^{-1}\left(t \underline{d} Q_{t}^{B_{0}}\right)+i M_{B_{0}}^{-1}\left(t \underline{d}_{B_{0}}^{*} Q_{t}^{B_{0}}\right), \\
t T_{B_{0}} Q_{t}^{B_{0}} & =-i M_{B_{0}}^{-1}\left(t \underline{d} Q_{t}^{B_{0}}\right)-i M_{B_{0}}^{-1}\left(t \underline{d}_{B_{0}}^{*} Q_{t}^{B_{0}}\right), \\
-i t \underline{d}_{B_{0}}^{*} Q_{t}^{B_{0}} & =-i\left(t \underline{d}_{B_{0}}^{*} Q_{t}^{B_{0}}\right),
\end{aligned}
$$

respectively. Thus we get from (4.3)-(4.8), after multiplication with $\chi(t)$, that

$$
\left\|\chi(t)\left(Q_{t}^{B}-Q_{t}^{B_{0}}\right)\right\|_{\mathrm{op}} \leqslant C\left(\| \| \chi(t) Q_{t}^{B} \|_{\mathrm{op}}+1\right)\left\|B-B_{0}\right\|_{\infty},
$$

and thus if $\left\|B-B_{0}\right\|_{\infty} \leqslant \varepsilon:=1 /(2 C)$ that

$$
\left\|\chi Q_{t}^{B}\right\|_{\mathrm{op}} \leqslant \frac{C\left\|B-B_{0}\right\|_{\infty}+\left\|Q_{t}^{B_{0}}\right\|_{\mathrm{op}}}{1-C\left\|B-B_{0}\right\|_{\infty}} \leqslant C^{\prime},
$$

since $\left\|\chi \chi(t) Q_{t}^{B}\right\|_{\mathrm{op}}<\infty$. Since this estimate is independent of $\tau$, it follows that $\left\|Q_{t}^{B} f\right\| \lesssim$ $\|f\|$.

Before turning to the proofs of (4.9) and (4.10), we summarise fundamental techniques from harmonic analysis that we shall need. We use the following dyadic decomposition of $\mathbf{R}^{n}$. Let $\Delta=$ $\bigcup_{j=-\infty}^{\infty} \Delta_{2^{j}}$ where $\Delta_{t}:=\left\{2^{j}\left(k+(0,1]^{n}\right): k \in \mathbf{Z}^{n}\right\}$ if $2^{j-1}<t \leqslant 2^{j}$. For a dyadic cube $Q \in \Delta_{2^{j}}$, denote by $l(Q)=2^{j}$ its sidelength, by $|Q|=2^{n j}$ its volume and by $R_{Q}:=Q \times\left(0,2^{j}\right] \subset \mathbf{R}_{+}^{n+1}$ the associated Carleson box. Let the dyadic averaging operator $A_{t}: \mathcal{H} \rightarrow \mathcal{H}$ be given by

$$
A_{t} u(x):=u_{Q}:=\int_{Q} u(y) d y=\frac{1}{|Q|} \int_{Q} u(y) d y
$$

for every $x \in \mathbf{R}^{n}$ and $t>0$, where $Q \in \Delta_{t}$ is the unique dyadic cube that contains $x$.
We now survey known results for a family of operators $\Theta_{t}: \mathcal{H} \rightarrow \mathcal{H}, t>0$. For the proofs we refer to [4] and [10].

Definition 4.2. By the principal part of $\left(\Theta_{t}\right)_{t>0}$ we mean the multiplication operators $\gamma_{t}$ defined by

$$
\gamma_{t}(x) w:=\left(\Theta_{t} w\right)(x)
$$

for every $w \in \bigwedge$. We view $w$ on the right-hand side of the above equation as the constant function defined on $\mathbf{R}^{n}$ by $w(x):=w$. We identify $\gamma_{t}(x)$ with the (possibly unbounded) multiplication operator $\gamma_{t}: f(x) \mapsto \gamma_{t}(x) f(x)$.

Lemma 4.3. Assume that $\Theta_{t}$ has $L_{2}$ off-diagonal bounds with exponent $M>n$. Then $\Theta_{t}$ extends to a bounded operator $L_{\infty} \rightarrow L_{2}^{\text {loc }}$. In particular we have well-defined functions $\gamma_{t} \in$ $L_{2}^{\text {loc }}\left(\mathbf{R}^{n} ; \mathcal{L}(\bigwedge)\right)$ with bounds

$$
\int_{Q}\left|\gamma_{t}(y)\right|^{2} d y \lesssim\left\|\Theta_{t}\right\|_{\mathrm{off}}^{2}
$$

for all $Q \in \Delta_{t}$. Moreover $\left\|\gamma_{t} A_{t}\right\| \lesssim\left\|\Theta_{t}\right\|_{\text {off }}$ uniformly for all $t>0$.
We have the following principal part approximation $\Theta_{t} \approx \gamma_{t}$.
Lemma 4.4. Assume that $\Theta_{t}$ has $L_{2}$ off-diagonal bounds with exponent $M>3 n$ and let $F_{t}$ : $\mathbf{R}^{n} \rightarrow \bigwedge$ be a family of functions. Then

$$
\left\|\left(\Theta_{t}-\gamma_{t} A_{t}\right) F_{t}\right\|\|\lesssim\| \Theta_{t}\left\|_{\text {off }}\right\| t \nabla F_{t}\| \|
$$

where $\nabla f=\nabla \otimes f=\sum_{j=1}^{n} e_{i} \otimes\left(\partial_{j} f\right)$ denotes the full differential of $f$. Moreover, if $P_{t}$ is a standard Fourier mollifier (we shall use $P_{t}=\left(1+t^{2} \Pi^{2}\right)^{-1}$ where $\Pi=\Gamma+\Gamma^{*}$ and $\Gamma$ is an exact nilpotent, homogeneous first order partial differential operator with constant coefficients as in [10]) and $F_{t}=P_{t} f$ for some $f \in \mathcal{H}$, then $\left\|\gamma_{t} A_{t}\left(P_{t}-I\right) f\right\| \lesssim\left\|\Theta_{t}\right\|_{\text {off }}\|f\|$ and $\left\|t \nabla P_{t}\right\|_{\mathrm{op}} \leqslant C$. Thus

$$
\left\|\mid\left(\Theta_{t} P_{t}-\gamma_{t} A_{t}\right) f\right\|\|\lesssim\| \Theta_{t}\left\|_{\text {off }}\right\| f \| .
$$

Definition 4.5. A function $\gamma(t, x): \mathbf{R}_{+}^{n+1} \rightarrow \bigwedge$ is called a Carleson function if there exists $C<$ $\infty$ such that

$$
\iint_{R_{Q}}|\gamma(t, x)|^{2} \frac{d x d t}{t} \leqslant C^{2}|Q|
$$

for all cubes $Q \subset \mathbf{R}^{n}$. Here $R_{Q}:=Q \times(0, l(Q)]$ is the Carleson box over $Q$. We define the Carleson norm $\left\|\gamma_{t}\right\|_{C}$ to be the smallest constant $C$.

We use Carleson's lemma in the following form.
Lemma 4.6. Let $\gamma_{t}(x)=\gamma(t, x): \mathbf{R}_{+}^{n+1} \rightarrow \bigwedge$ be a Carleson function and let $F_{t}(x)=F(t, x)$ : $\mathbf{R}_{+}^{n+1} \rightarrow \bigwedge$ be a family of functions. Then

$$
\left\|\gamma_{t} F_{t}\right\| \lesssim \lesssim \gamma_{t}\left\|_{C}\right\| N_{*}\left(F_{t}\right) \|,
$$

where $N_{*}\left(F_{t}\right)(x):=\sup _{|y-x|<t}\left|F_{t}(y)\right|$ denotes the non-tangential maximal function of $F_{t}$. In particular, if $F_{t}=A_{t} f$ for some $f \in \mathcal{H}$, then $\left\|N_{*}\left(A_{t} f\right)\right\| \lesssim\|M(f)\| \lesssim\|f\|$, where $M(f)(x):=$ $\sup _{r>0} f_{B(x, r)}|f(y)|$ dy denotes the Hardy-Littlewood maximal function, and thus

$$
\left\|\gamma_{t} A_{t} f\right\| \lesssim\left\|\gamma_{t}\right\|_{C}\|f\| .
$$

Lemma 4.7. Assume that $\Theta_{t}$ has $L_{2}$ off-diagonal bounds with exponent $M>n$. Then

$$
\left\|\Theta_{t} f\right\|_{C} \lesssim\left(\left\|\Theta_{t}\right\|_{\mathrm{op}}+\left\|\Theta_{t}\right\|_{\text {off }}\right)\|f\|_{\infty}
$$

for every $f \in L_{\infty}\left(\mathbf{R}^{n} ; \bigwedge\right)$. In particular, choosing $f=w=$ constant we obtain

$$
\left\|\gamma_{t}\right\|_{C} \lesssim\left\|\Theta_{t}\right\|_{\mathrm{op}}+\left\|\Theta_{t}\right\|_{\mathrm{off}} .
$$

### 4.1. Perturbation of block coefficients

In this section we assume that $B_{0}$ is a block matrix, i.e. we assume that

$$
B_{0}=\left[\begin{array}{cc}
\left(B_{0}\right)_{\perp \perp} & 0 \\
0 & \left(B_{0}\right)_{\| \|}
\end{array}\right]
$$

in the splitting $\mathcal{H}=N^{-} \mathcal{H} \oplus N^{+} \mathcal{H}$. Our goal is to prove Theorem 1.4 by verifying the hypothesis of Lemma 4.1. We recall from Lemma 3.1 that

$$
T_{B_{0}}=\Pi_{B_{0}}=\Gamma+B_{0}^{-1} \Gamma^{*} B_{0}, \quad \text { where } \Gamma:=-i N \underline{d},
$$

is an operator of the form treated in [10]. Thus with a slight change of notation for $\mathcal{E}$, we need to prove the following.

Theorem 4.8. If $B_{0}$ is a block matrix and if $\left(\widetilde{Q}_{t}\right)_{t>0}$ is an operator family with $L_{2}$ off-diagonal bounds $\left\|\widetilde{Q}_{t}\right\|_{\text {off }} \leqslant C$ as in Definition 2.24 , then with notation as above

$$
\begin{align*}
& \left\|\widetilde{Q}_{t} \mathcal{E}\left(t^{2} \Gamma \Gamma_{B_{0}}^{*} P_{t}^{B_{0}}\right)\right\|_{\mathrm{op}} \leqslant C\left(\left\|\widetilde{Q}_{t}\right\|_{\mathrm{op}}+1\right)\|\mathcal{E}\|_{\infty}  \tag{4.11}\\
& \left\|\widetilde{Q}_{t} \mathcal{E}\left(t^{2} \Gamma_{B_{0}}^{*} \Gamma P_{t}^{B_{0}}\right)\right\|_{\mathrm{op}} \leqslant C\left(\left\|\widetilde{Q}_{t}\right\|_{\mathrm{op}}+1\right)\|\mathcal{E}\|_{\infty} \tag{4.12}
\end{align*}
$$

where $P_{t}^{B_{0}}=\left(1+t^{2} \Pi_{B_{0}}^{2}\right)^{-1}$.
We note that if (4.11) holds with $\Gamma$ replaced by $\Gamma^{*}$ and with $B_{0}$ replaced by $B_{0}^{-1}$, then also (4.12) holds. This follows from the conjugation formula

$$
\begin{aligned}
& \widetilde{Q}_{t} \mathcal{E}\left(t^{2} \Gamma_{B_{0}}^{*} \Gamma P_{t}^{B_{0}}\right) \\
& \quad=B_{0}^{-1}\left(B_{0} \widetilde{Q}_{t} B_{0}^{-1}\right)\left(B_{0} \mathcal{E} B_{0}^{-1}\right)\left(t^{2} \Gamma^{*} \Gamma_{B_{0}^{-1}}\left(1+t^{2}\left(\Gamma^{*}+B_{0} \Gamma B_{0}^{-1}\right)^{2}\right)^{-1}\right) B_{0}
\end{aligned}
$$

Thus it suffices to prove (4.11), as long as we only use properties of ( $\Gamma, B_{0}$ ) shared with $\left(\Gamma^{*}, B_{0}^{-1}\right)$. To this end, we let $\Theta_{t}$ be the operator

$$
\Theta_{t}:=\widetilde{Q}_{t} \mathcal{E}\left(t^{2} \Gamma \Gamma_{B_{0}}^{*} P_{t}^{B_{0}}\right)
$$

and denote by $\gamma_{t}(x)$ its principal part as in Definition 4.2. We note that we have a Hodge type splitting $\mathcal{H}=\mathrm{N}(\Gamma) \oplus \mathrm{N}\left(\Gamma_{B_{0}}^{*}\right)$ by Lemma 2.21(i), and since $\left.\Theta_{t}\right|_{\mathrm{N}\left(\Gamma_{B_{0}}^{*}\right)}=0$ it suffices to bound $\left\|\left\|\Theta_{t} f\right\|\right.$ for $f \in \mathrm{~N}(\Gamma)$. We do this by writing

$$
\begin{equation*}
\Theta_{t} f=\Theta_{t}\left(I-P_{t}\right) f+\left(\Theta_{t} P_{t}-\gamma_{t} A_{t}\right) f+\gamma_{t} A_{t} f \tag{4.13}
\end{equation*}
$$

where $\Pi:=\Gamma+\Gamma^{*}$ is the corresponding unperturbed operator and $P_{t}:=\left(1+t^{2} \Pi^{2}\right)^{-1}$ and $Q_{t}:=t \Pi\left(1+t^{2} \Pi^{2}\right)^{-1}$.

Lemma 4.9. We have, for all $f \in \mathrm{~N}(\Gamma)$, the estimate

$$
\left\|\Theta_{t}\left(I-P_{t}\right) f\right\| \leqslant C\left(\left\|\widetilde{Q}_{t}\right\|_{\mathrm{op}}+1\right)\|\mathcal{E}\|_{\infty}\|f\| .
$$

Proof. If $f \in \mathrm{~N}(\Gamma)$, then $\left(I-P_{t}\right) f=t \Gamma Q_{t} f \in \mathrm{~N}(\Gamma)$, which shows that $t^{2} \Gamma \Gamma_{B_{0}}^{*} P_{t}^{B_{0}}(I-$ $\left.P_{t}\right) f=\left(I-P_{t}^{B_{0}}\right)\left(I-P_{t}\right) f$. To prove the estimate, we write

$$
\Theta_{t}\left(I-P_{t}\right) f=\widetilde{Q}_{t} \mathcal{E}\left(I-P_{t}^{B_{0}}\right) t \Gamma Q_{t} f=\widetilde{Q}_{t} \mathcal{E} f-\widetilde{Q}_{t} \mathcal{E} P_{t} f-\widetilde{Q}_{t} \mathcal{E} Q_{t}^{B_{0}}\left(\Pi_{B_{0}}^{-1} \Gamma\right) Q_{t} f
$$

where we recall that $Q_{t}^{B_{0}}=t \Pi_{B_{0}} P_{t}^{B_{0}}$. Clearly $\left\|\widetilde{Q}_{t} \mathcal{E} f\right\| \lesssim\left\|\widetilde{Q}_{t}\right\|\left\|\|\mathcal{E}\|_{\infty}\right\| f \|$. For the second term, we write $\widetilde{\Theta}_{t}:=\widetilde{Q}_{t} \mathcal{E}$ with principal part $\widetilde{\gamma}_{t}$ as in Definition 4.2 and estimate

$$
\begin{aligned}
\left\|\widetilde{\Theta}_{t} P_{t} f\right\| \| & \leqslant\left\|\left(\widetilde{\Theta}_{t} P_{t}-\widetilde{\gamma}_{t} A_{t}\right) f\right\|\|+\| \widetilde{\gamma}_{t} A_{t} f \| \\
& \lesssim\|\mathcal{E}\|_{\infty}\|f\|+\left\|\widetilde{\gamma}_{t}\right\|_{C}\|f\| \lesssim\|\mathcal{E}\|_{\infty}\|f\|+\left(\left\|\widetilde{Q}_{t}\right\|_{\mathrm{op}}+1\right)\|\mathcal{E}\|_{\infty}\|f\| .
\end{aligned}
$$

In the second step, we used for the first term that $\left\|\widetilde{\Theta}_{t}\right\|_{\text {off }} \lesssim\|\mathcal{E}\|_{\infty}$ in Lemma 4.4 and for the second term we used Lemma 4.6. In the last step, we used Lemma 4.7 on the last term and that $\left\|\widetilde{\Theta}_{t}\right\|_{\text {op }} \lesssim\| \| \widetilde{Q}_{t}\left\|_{\text {op }}\right\| \mathcal{E} \|_{\infty}$. Finally we note that

$$
\left\|\widetilde{Q}_{t} \mathcal{E} Q_{t}^{B_{0}}\left(\Pi_{B_{0}}^{-1} \Gamma\right) Q_{t} f\right\| \lesssim\|\mathcal{E}\|_{\infty}\left\|Q_{t} f\right\| \lesssim\|\mathcal{E}\|_{\infty}\|f\|
$$

since the Hodge projection $\Pi_{B_{0}}^{-1} \Gamma$ is bounded by Lemma 2.21(i).
To estimate the last term in (4.13) we shall apply a local $T(b)$ theorem as in $[4,10]$. We here give an alternative construction of test functions to those used in [10], more in the spirit of the original proof of the Kato square root problem [4].

Lemma 4.10. Let $\Gamma$ be a nilpotent operator in $\mathcal{H}$, which is a homogeneous, first order partial differential operator with constant coefficients. Denote by $\bigwedge_{\Gamma} \subset \bigwedge$ the image of the linear functions $u: \mathbf{R}^{n} \rightarrow \bigwedge$ under $\Gamma$, where we identify $\bigwedge$ with the constant functions $\mathbf{R}^{n} \rightarrow \bigwedge$.

Then for each $w \in \bigwedge_{\Gamma}$ with $|w|=1$, each cube $Q \subset \mathbf{R}^{n}$ and each $\varepsilon>0$, there exists a test function $f_{Q, \varepsilon}^{w} \in \mathcal{H}$ such that $f_{Q, \varepsilon}^{w} \in \mathrm{R}(\Gamma),\left\|f_{Q, \varepsilon}^{w}\right\| \lesssim|Q|^{1 / 2}$,

$$
\left\|\Gamma_{B_{0}}^{*} f_{Q, \varepsilon}^{w}\right\| \lesssim \frac{1}{\varepsilon l(Q)}|Q|^{1 / 2} \quad \text { and } \quad\left|f_{Q} f_{Q, \varepsilon}^{w}-w\right| \lesssim \varepsilon^{1 / 2}
$$

Proof. Let $u(x)$ be a linear function such that $w=\Gamma u$ and $\sup _{3 Q}|u(x)| \lesssim l(Q)$, and define $w_{Q}:=\Gamma\left(\eta_{Q} u\right)$, where $\eta_{Q}$ is a smooth cutoff such that $\left.\eta_{Q}\right|_{2 Q}=1, \operatorname{supp}\left(\eta_{Q}\right) \subset 3 Q$ and $\left\|\nabla \eta_{Q}\right\|_{\infty} \lesssim 1 / l(Q)$. It follows that

$$
w_{Q} \in \mathrm{R}(\Gamma),\left.\quad w_{Q}\right|_{2 Q}=w, \quad \operatorname{supp} w_{Q} \subset 3 Q \quad \text { and } \quad\left\|w_{Q}\right\|_{\infty} \leqslant C .
$$

Next we define the test function $f_{Q, \varepsilon}^{w}:=P_{\varepsilon l}^{B_{0}} w_{Q}$, where we write $l=l(Q)$. Using Corollary 2.23, it follows that $\left\|f_{Q, \varepsilon}^{w}\right\| \lesssim|Q|^{1 / 2}$ and $\left\|\Gamma_{B_{0}}^{*} f_{Q, \varepsilon}^{w}\right\| \lesssim \frac{1}{\varepsilon l(Q)}|Q|^{1 / 2}$ and since $\Gamma$ commutes with $P_{\varepsilon l}^{B_{0}}$, it follows that $f_{Q, \varepsilon}^{w} \in \mathrm{R}(\Gamma)$. To verify the accretivity property, we make use of [10, Lemma 5.6] which shows that

$$
\begin{aligned}
\left|f_{Q} f_{Q, \varepsilon}^{w}-w\right| & =\left|f_{Q}\left(I-P_{\varepsilon l}^{B_{0}}\right) w_{Q}\right|=\left|f_{Q} \varepsilon l \Gamma\left(Q_{\varepsilon l}^{B_{0}} w_{Q}\right)\right| \\
& \lesssim \varepsilon^{1 / 2}\left(f_{Q}\left|Q_{\varepsilon l}^{B_{0}} w_{Q}\right|^{2}\right)^{1 / 4}\left(f_{Q}\left|\varepsilon l \Gamma Q_{\varepsilon l}^{B_{0}} w_{Q}\right|^{2}\right)^{1 / 4} \lesssim \varepsilon^{1 / 2}
\end{aligned}
$$

where we used that $w_{Q} \in \mathrm{R}(\Gamma)$ in the second step.
Proof of Theorem 4.8. We have seen that it suffices to prove (4.11), and to bound each term in (4.13) for $f \in \mathrm{~N}(\Gamma)$. The first term is estimated by Lemma 4.9 and the second by Lemma 4.4.

To prove that the last term has estimate $\left\|\gamma_{t} A_{t} f\right\| \lesssim\left(\left\|\widetilde{Q}_{t}\right\|_{\text {op }}+1\right)\|\mathcal{E}\|_{\infty}\|f\|$, it suffices by Lemma 4.6 to prove that

$$
\left\|\gamma_{t}\right\|_{C} \lesssim\left(\left\|\widetilde{Q}_{t}\right\|_{\text {op }}+1\right)\|\mathcal{E}\|_{\infty} .
$$

To this end, we apply the local $T(b)$ argument and stopping time argument in [10, Section 5.3]. Note that $\mathrm{R}(\Gamma) \subset L_{2}\left(\mathbf{R}^{n} ; \bigwedge_{\Gamma}\right)$ and thus, since $\Gamma$ is an exact nilpotent operator, it follows that $\mathrm{N}(\Gamma) \subset L_{2}\left(\mathbf{R}^{n} ; \bigwedge_{\Gamma}\right)$. Furthermore $A_{t} f \in L_{2}\left(\mathbf{R}^{n} ; \bigwedge_{\Gamma}\right)$ if $f \in L_{2}\left(\mathbf{R}^{n} ; \bigwedge_{\Gamma}\right)$. Thus it suffices to bound the Carleson norm of $\gamma_{t}(x)$ seen as a linear operator $\gamma_{t}(x): \bigwedge_{\Gamma} \rightarrow \bigwedge$. The conical decomposition $\bigcup_{\nu \in \mathcal{V}} K_{\nu}$ of the space of matrices performed in [10, Section 5.3], here decomposes the space $\mathcal{L}\left(\bigwedge_{\Gamma} ; \bigwedge\right)$ and for the fixed unit matrix $\nu \in \mathcal{L}\left(\bigwedge_{\Gamma} ; \bigwedge\right)$ we choose $w \in \bigwedge_{\Gamma}$ and $\hat{w} \in \Lambda$ such that $|\hat{w}|=|w|=1$ and $v^{*}(\hat{w})=w$. With the stopping time argument in [10, Section 5.3], using the new test functions $f_{Q}^{w}$ from Lemma 4.10, we obtain

$$
\left\|\gamma_{t}\right\|_{C}^{2} \lesssim \sup _{Q, w \in \bigwedge_{\Gamma},|w|=1} \frac{1}{|Q|} \iint_{R_{Q}}\left|\gamma_{t}(x) A_{t} f_{Q}^{w}(x)\right|^{2} \frac{d x d t}{t}
$$

where $f_{Q}^{w}=f_{Q, \varepsilon}^{w}$ for a small enough but fixed $\varepsilon$. To estimate the right-hand side we use (4.13) with $f$ replaced with the test function $f_{Q}^{w}$, estimate the first two terms with Lemma 4.9 (which works since $\left.f_{Q}^{w} \in \mathrm{R}(\Gamma)\right)$ and Lemma 4.4, and obtain

$$
\begin{aligned}
\iint_{R_{Q}}\left|\gamma_{t} A_{t} f_{Q}^{w}\right|^{2} \frac{d x d t}{t} & \lesssim \iint_{R_{Q}}\left|\gamma_{t} A_{t} f_{Q}^{w}-\Theta_{t} f_{Q}^{w}\right|^{2} \frac{d x d t}{t}+\iint_{R_{Q}}\left|\Theta_{t} f_{Q}^{w}\right|^{2} \frac{d x d t}{t} \\
& \lesssim\left(\left\|\widetilde{Q}_{t}\right\|_{\mathrm{op}}+1\right)^{2}\|\mathcal{E}\|_{\infty}^{2}|Q|+\iint_{R_{Q}}\left|\Theta_{t} f_{Q}^{w}\right|^{2} \frac{d x d t}{t}
\end{aligned}
$$

Using that $\left\|\Gamma_{B_{0}}^{*} f_{Q}^{w}\right\| \lesssim \frac{1}{l(Q)}|Q|^{1 / 2}$ we then get

$$
\left\|\Theta_{t} f_{Q}^{w}\right\|=\left\|\widetilde{Q}_{t} \mathcal{E}\left(t \Gamma P_{t}^{B_{0}}\right) t\left(\Gamma_{B_{0}}^{*} f_{Q}^{w}\right)\right\| \lesssim\|\mathcal{E}\|_{\infty} \frac{t}{l(Q)}|Q|^{1 / 2}
$$

This yields

$$
\iint_{R_{Q}}\left|\Theta_{t} f_{Q}^{w}\right|^{2} \frac{d x d t}{t} \lesssim|Q|\|\mathcal{E}\|_{\infty}^{2} \int_{0}^{l(Q)}\left(\frac{t}{l(Q)}\right)^{2} \frac{d t}{t} \lesssim|Q|\|\mathcal{E}\|_{\infty}^{2}
$$

which proves that

$$
\iint_{R_{Q}}\left|\gamma_{t} A_{t} f_{Q}^{w}\right|^{2} \frac{d x d t}{t} \lesssim\left(\| \| \widetilde{Q}_{t} \|_{\mathrm{op}}+1\right)^{2}\|\mathcal{E}\|_{\infty}^{2}|Q|
$$

Remark 4.11. (i) Note that using the new test function from Lemma 4.10 simplifies the estimate of the term $\left(\gamma_{t} A_{t}-\Theta_{t}\right) f_{Q}^{w}$ in the above proof as compared with the proof of [10, Proposition 5.9]. The useful new property of the test functions from Lemma 4.10 is that they belong to $\mathrm{R}(\Gamma)$.
(ii) In the proof of Theorem 4.8, we only estimate the Carleson norm of the restriction of the matrix $\gamma_{t}$ to the subspace $\bigwedge_{\Gamma}$ as this suffices since we want to bound the quadratic norm of $\gamma_{t} A_{t} f$, and $A_{t} f$ is always $\bigwedge_{\Gamma}$-valued. However, to prove (4.11) and (4.12), we use $\Gamma$ being either $N m d$ or $N m d^{*}$. In these two cases, the space $\bigwedge_{\Gamma}$ is either the orthogonal complement of $\operatorname{span}\left\{1, e_{0}\right\}$ or $\operatorname{span}\left\{e_{0,1, \ldots, n}, e_{1, \ldots, n}\right\}$, respectively. Note also that block matrices preserve these spaces $\bigwedge_{\Gamma}^{\perp}$. It is seen that in these two cases $\gamma_{t}=0$ on $\bigwedge_{\Gamma}^{\perp}$, so actually we do get an estimate of the Carleson norm of the whole matrix $\gamma_{t}$.

### 4.2. Perturbation of vector coefficients

In this section we assume that the unperturbed coefficients $B_{0}$ are of the form

$$
B_{0}=I \oplus A_{0} \oplus I \oplus \cdots \oplus I
$$

i.e. $B_{0}$ only acts non-trivially on the vector part, and that $A_{0}$ is a matrix such that $T_{B_{0}}$ has quadratic estimates. Note that this hypothesis is true if $A_{0}$ is either real symmetric, constant or of block form, by Theorem 3.12, Proposition 3.5 and Theorem 3.3, respectively. Our goal is to prove Theorem 1.1, by verifying the hypothesis of Lemma 4.1, as well as proving Theorem 1.3. We start by reformulating Lemma 4.1 in terms of $e^{-t\left|T_{B_{0}}\right|}$, acting only on functions $f$ in one of the two Hardy spaces $E_{B_{0}}^{ \pm} \mathcal{H}$, instead of $P_{t}^{B_{0}}$.

Theorem 4.12. If $B_{0}$ is as above and if $\left\|\widetilde{Q}_{t}\right\|_{\text {off }} \leqslant C$, then for all $f \in E_{B_{0}}^{+} \mathcal{H}$ we have estimates

$$
\begin{align*}
\left\|\widetilde{Q}_{t} \mathcal{E}\left(F_{t}-f\right)\right\| & \leqslant C\left(\left\|\widetilde{Q}_{t}\right\|_{\mathrm{op}}+1\right)\|\mathcal{E}\|_{\infty}\|f\|  \tag{4.14}\\
\left\|\widetilde{Q}_{t} \mathcal{E} \underline{d} \int_{0}^{t} F_{s} d s\right\| & \leqslant C\left(\left\|\widetilde{Q}_{t}\right\|_{\mathrm{op}}+1\right)\|\mathcal{E}\|_{\infty}\|f\|, \tag{4.15}
\end{align*}
$$

where $F_{t}:=e^{-t\left|T_{B_{0}}\right|} f$ is the extension of $f$ as in Lemma 2.49. The corresponding estimates for $f \in E_{B_{0}}^{-} \mathcal{H}$ also hold.

Proof that Theorem 4.12 implies (4.9) and (4.10). We first note that it suffices to prove (4.9) and (4.10) for all $f \in E_{B_{0}}^{+} \mathcal{H}$ and all $f \in E_{B_{0}}^{-} \mathcal{H}$ since we have a Hardy space splitting $\mathcal{H}=$ $E_{B_{0}}^{+} \mathcal{H} \oplus E_{B_{0}}^{-} \mathcal{H}$. We only consider $f \in E_{B_{0}}^{+} \mathcal{H}$ since the proof for $f \in E_{B_{0}}^{-} \mathcal{H}$ is similar.

Now let $f \in E_{B_{0}}^{+} \mathcal{H}$ and use Proposition 2.31, which shows that if $\psi \in \Psi\left(S_{v}^{o}\right)$, then $\left\|\left\|\psi_{t}\left(T_{B_{0}}\right)\right\|_{\mathrm{op}} \lesssim\right\| Q_{t}^{B_{0}} \|_{\mathrm{op}} \leqslant C$. For the estimate (4.9), we write

$$
\widetilde{Q}_{t} \mathcal{E} t \underline{d} Q_{t}^{B_{0}}=\widetilde{Q}_{t} \mathcal{E}\left(\underline{d}_{B_{0}}^{-1}\right)\left(I-e^{-t\left|T_{B_{0}}\right|}\right)+\widetilde{Q}_{t} \mathcal{E}\left(\underline{d} T_{B_{0}}^{-1}\right) \psi_{t}\left(T_{B_{0}}\right),
$$

where $\psi(z)=e^{-|z|}-\left(1+z^{2}\right)^{-1}$. Note for the first term that $T_{B_{0}}^{-1}\left(I-e^{-t\left|T_{B_{0}}\right|}\right) f=T_{B_{0}}^{-1}(f-$ $\left.F_{t}\right)=-T_{B_{0}}^{-1} \int_{0}^{t} \partial_{s} F_{s} d s=\int_{0}^{t} F_{s} d s$. Therefore Theorem 4.12, the boundedness of $\underline{d} T_{B_{0}}^{-1}$ and Proposition 2.31 give the estimate (4.9).

For the estimate (4.10), note that it suffices to estimate

$$
\widetilde{Q}_{t} \mathcal{E} t T_{B_{0}} Q_{t}^{B_{0}}=\widetilde{Q}_{t} \mathcal{E}\left(I-P_{t}^{B_{0}}\right)
$$

since $i M_{B_{0}} T_{B_{0}}=\underline{d}+\underline{d}_{B_{0}}^{*}$. But this follows immediately from (4.14) since

$$
\widetilde{Q}_{t} \mathcal{E}\left(I-P_{t}^{B_{0}}\right) f=\widetilde{Q}_{t} \mathcal{E}\left(f-F_{t}\right)+\widetilde{Q}_{t} \mathcal{E} \psi_{t}\left(T_{B_{0}}\right) f
$$

with the same $\psi$ as above.
Remark 4.13. In the case when $B_{0}$ is a constant matrix, we can estimate $\widetilde{Q}_{t} \mathcal{E}\left(I-P_{t}^{B_{0}}\right) f$ directly. We prove that

$$
\left\|\widetilde{Q}_{t} \mathcal{E}\left(I-P_{t}^{B_{0}}\right) f\right\| \lesssim\left(\left\|\widetilde{Q}_{t}\right\|_{\mathrm{op}}+1\right)\|\mathcal{E}\|_{\infty}\|f\|
$$

as follows. Clearly $\left\|\widetilde{Q}_{t} \mathcal{E}\right\|_{\text {op }} \lesssim\| \| \widetilde{Q}_{t}\left\|_{\text {op }}\right\| \mathcal{E} \|_{\infty}$. For the second term we write $\Theta_{t}:=\widetilde{Q}_{t} \mathcal{E} P_{t}^{B_{0}}$. Inserting a standard Fourier mollifier $P_{t}$, we write $\Theta_{t}=\Theta_{t}\left(I-P_{t}\right)+\Theta_{t} P_{t}$. Here $\| \Theta_{t}(I-$ $\left.P_{t}\right) f\|\lesssim\| P_{t}^{B_{0}}\left(I-P_{t}\right) f\|\lesssim\| f \|$ is easily verified using the Fourier transform. On the other hand,

$$
\left\|\Theta_{t} P_{t} f\right\| \leqslant\left\|\left(\Theta_{t} P_{t}-\gamma_{t} A_{t}\right) f\right\|\|+\| \gamma_{t} A_{t} f\|\lesssim\| f\|+\| \gamma_{t}\left\|_{C}\right\| f \|,
$$

using Lemmas 4.4 and 4.6. However, since $B_{0}$ is constant, we have that $T_{B_{0}} w=0$ if $w$ is a constant function, and therefore

$$
\gamma_{t} w=\widetilde{Q}_{t} \mathcal{E}\left(P_{t}^{B_{0}} w\right)=\widetilde{Q}_{t} \mathcal{E} w
$$

Therefore Lemma 4.7 shows that $\left\|\gamma_{t}\right\|_{C} \lesssim\left(\| \| \widetilde{Q}_{t} \|_{\text {op }}+1\right)\|\mathcal{E}\|_{\infty}$.
We now set up some notation for the proof of Theorem 4.12. We decompose the function $F_{t}:=e^{-t\left|T_{B_{0}}\right|} f$, where $f \in E_{B_{0}}^{+} \mathcal{H}$, as

$$
\begin{equation*}
F_{t}=F_{t}^{0}+\left(F_{t}^{1, \perp}+F_{t}^{1, \|}\right)+F_{t}^{2}+\cdots+F_{t}^{n+1} \tag{4.16}
\end{equation*}
$$

and similarly for $f=\lim _{t \rightarrow 0^{+}} F_{t}$. It is important to note the special property that the normal component of the vector part $F_{t}^{1, \perp}=F_{t}^{1,0} e_{0}$ has by Lemma 2.12: it satisfies the divergence form second order equation $\operatorname{div}_{t, x} A_{0}(x) \nabla_{t, x} F^{1,0}=0$. Furthermore, we decompose the matrix $A_{0}$ as

$$
A_{0}=\left[\begin{array}{cc}
a_{\perp \perp} & a_{\perp \|} \\
a_{\| \perp} & a_{\| \|}
\end{array}\right]
$$

in the splitting $\mathcal{H}=N^{-} \mathcal{H} \oplus N^{+} \mathcal{H}$. We view the components $a_{\perp \perp}, a_{\perp \|}, a_{\| \perp}$ and $a_{\| \| \|}$as operators, and write $a_{\perp \perp}\left(f^{1,0} e_{0}\right)=\left(a_{00} f^{1,0}\right) e_{0}, a_{\perp \|} f^{1, \|}=\left(a_{0 \|} \cdot f^{1, \|}\right) e_{0}$ and $a_{\| \perp}\left(f^{1,0} e_{0}\right)=f^{1,0} a_{\| 0}$, where $a_{00}$ is a scalar and $a_{0 \|}$ and $a_{\| 0}$ are vectors.

We introduce an auxiliary block matrix

$$
\hat{B}_{0}=I \oplus \hat{A}_{0} \oplus I \oplus \cdots \oplus I, \quad \hat{A}_{0}=\left[\begin{array}{cc}
a_{\perp \perp} & 0 \\
0 & a_{\| \|}
\end{array}\right]
$$

Lemma 4.14. Let $F_{t}:=e^{-t\left|T_{B_{0}}\right|} f$, where $f \in E_{B_{0}}^{+} \mathcal{H}$, so that $\left(\partial_{t}+T_{B_{0}}\right) F_{t}=0$. Then

$$
\left(\partial_{t}+T_{\hat{B}_{0}}\right) F_{t}=-i\left(a_{00}^{-1} a_{\perp \|}\right) \underline{d} F_{t}^{1, \perp}-i \underline{d}_{\hat{B}_{0}}^{*}\left(a_{\| \|}^{-1} a_{\| \perp}\right) F_{t}^{1, \perp}=\partial_{t} F_{t}^{1, \perp}+i{\underline{d_{\hat{B}_{0}}}}_{*} F_{t}^{1, \|} .
$$

Proof. To prove the first identity, note that

$$
\begin{aligned}
\left(\partial_{t}+T_{\hat{B}_{0}}\right) F_{t}= & \left(T_{\hat{B}_{0}}-T_{B_{0}}\right) F_{t} \\
= & -i\left(\hat{B}_{0} N^{+}-N^{-} \hat{B}_{0}\right)^{-1}\left(\hat{B}_{0} \underline{d}+\underline{d}^{*} \hat{B}_{0}\right) F_{t} \\
& +i\left(B_{0} N^{+}-N^{-} B_{0}\right)^{-1}\left(B_{0} \underline{d}+\underline{d}^{*} B_{0}\right) F_{t} \\
= & i\left(\hat{B}_{0} N^{+}-N^{-} \hat{B}_{0}\right)^{-1}\left(\left(B_{0}-\hat{B}_{0}\right) \underline{d}+\underline{d}^{*}\left(B_{0}-\hat{B}_{0}\right)\right) F_{t}
\end{aligned}
$$

since $\hat{B}_{0} N^{+}-N^{-} \hat{B}_{0}=B_{0} N^{+}-N^{-} B_{0}$. The vector part of this matrix is $\left[\begin{array}{cc}-a_{\perp \perp} & 0 \\ 0 & a_{\| \|}\end{array}\right]$. Furthermore

$$
B_{0}-\hat{B}_{0}=0 \oplus\left[\begin{array}{cc}
0 & a_{\perp \|} \\
a_{\| \perp} & 0
\end{array}\right] \oplus 0 \oplus \cdots \oplus 0
$$

which shows that

$$
\left(B_{0}-\hat{B}_{0}\right) \underline{d} F_{t}=a_{\perp \|} \underline{d} F_{t}^{1, \perp} \quad \text { and } \quad \underline{d}^{*}\left(B_{0}-\hat{B}_{0}\right) F_{t}=\underline{d}^{*} a_{\| \perp} F_{t}^{1, \perp}
$$

using the mapping properties of $\underline{d}$ and $\underline{d}^{*}$ from Remark 2.6. Thus

$$
\begin{aligned}
\left(\partial_{t}+T_{\hat{B}_{0}}\right) F_{t} & =-i a_{\perp \perp}^{-1}\left(a_{\perp \|} \underline{d} F_{t}^{1, \perp}+\underline{d}^{*} a_{\| \perp} F_{t}^{1, \perp}\right) \\
& =-i\left(a_{00}^{-1} a_{\perp \|}\right) \underline{d} F_{t}^{1, \perp}-i \underline{d}_{\hat{B}_{0}}^{*}\left(a_{\| \|}^{-1} a_{\| \perp}\right) F_{t}^{1, \perp}
\end{aligned}
$$

since $\hat{B}_{0} N^{+}-N^{-} \hat{B}_{0}=-a_{\perp \perp}$ on normal vector fields.
To prove the second identity, we multiply the Dirac equation $\left(\partial_{t}+T_{B_{0}}\right) F_{t}=0$ by $B_{0} N^{+}-$ $N^{-} B_{0}$ to obtain

$$
\left(B_{0} N^{+}-N^{-} B_{0}\right) \partial_{t} F_{t}-i\left(B_{0} \underline{d}+\underline{d}^{*} B_{0}\right) F_{t}=0
$$

The normal component of the vector part on the left-hand side is

$$
-a_{\perp \perp} \partial_{t} F_{t}^{1, \perp}-i\left(a_{\perp \|} \underline{d} F_{t}^{1, \perp}+\underline{d}^{*}\left(a_{\| \perp} F_{t}^{1, \perp}+a_{\| \|} F_{t}^{1, \|}\right)\right)=0 .
$$

Here we have used the expression for the vector part of $B_{0} N^{+}-N^{-} B_{0}$ above for the first term. For the other terms we recall from Remark 2.6 that the vector part of $\underline{d} F_{t}$ is $\underline{d} F_{t}^{1, \perp}$ which is
tangential, and that the normal vector part of $\underline{d}^{*} B_{0} F_{t}$ is $\underline{d}^{*}\left(B_{0} F_{t}\right)^{1, \|}$. Therefore, after multiplying the equation with $a_{00}^{-1}$, we obtain

$$
-i\left(a_{00}^{-1} a_{\perp \|}\right) \underline{d} F_{t}^{1, \perp}-i \underline{d}_{\hat{B}_{0}}^{*}\left(a_{\| \|}^{-1} a_{\| \perp}\right) F_{t}^{1, \perp}=\partial_{t} F_{t}^{1, \perp}+i \underline{d}_{\hat{B}_{0}}^{*} F_{t}^{1, \|}
$$

which proves the lemma.
As in Lemma 3.1 we note that since $\hat{B}_{0}$ is a block matrix

$$
T_{\hat{B}_{0}}=\Pi_{\hat{B}_{0}}=\Gamma+\hat{B}_{0}^{-1} \Gamma^{*} \hat{B}_{0}, \quad \text { where } \Gamma:=-i N \underline{d} .
$$

To prove Theorem 4.12 we shall need the following corollary to Theorem 4.8.
Corollary 4.15. Let $\hat{B}_{0}$ be the block matrix defined above, assume that $\left\|\widetilde{Q}_{t}\right\|_{\text {off }} \leqslant C$ and let $v \in$ $L_{\infty}\left(\mathbf{R}^{n} ; \bigwedge\right)$ be a function with norm $\|v\|_{\infty} \leqslant C$. Then for all $f \in E_{B_{0}}^{+} \mathcal{H}$ we have the estimates

$$
\begin{array}{r}
\left\|\widetilde{Q}_{t} \mathcal{E} t^{2} \Gamma \Gamma_{\hat{B}_{0}}^{*} P_{t}^{\hat{B}_{0}}\left(F_{t}^{1,0} v\right)\right\|\left\|\left(\left\|\widetilde{Q}_{t}\right\|_{\mathrm{op}}+1\right)\right\| \mathcal{E}\left\|_{\infty}\right\| f \|, \\
\left\|\widetilde{Q}_{t} \mathcal{E} P_{t}^{\hat{B}_{0}}\left(F_{t}^{1,0} v\right)\right\|  \tag{4.18}\\
\lesssim\left(\left\|\widetilde{Q}_{t}\right\|_{\mathrm{op}}+1\right)\|\mathcal{E}\|_{\infty}\|f\|
\end{array}
$$

where $F_{t}:=e^{-t\left|T_{B_{0}}\right|} f$ and $F_{t}^{1,0}=\left(F_{t}, e_{0}\right)$. The corresponding estimates for $f \in E_{B_{0}}^{-} \mathcal{H}$ also hold.

We defer the proof until the end of this section, and turn to a lemma in preparation for the proof of Theorem 4.12.

Lemma 4.16. If $f \in E_{B_{0}}^{+} \mathcal{H}$ and $F_{t}:=e^{-t\left|T_{B_{0}}\right|}$ f, then

$$
\begin{align*}
&\left\|\int_{0}^{t} \frac{s}{t}\left(s \partial_{s} F_{s}\right) \frac{d s}{s}\right\| \| \leqslant\left\|t \partial_{t} F_{t}\right\| \|  \tag{4.19}\\
& \lesssim\|f\|,  \tag{4.20}\\
&\left\|t t \underline{d} F_{t}\right\| \lesssim\left\|t \partial_{t} F_{t}\right\| \| \\
& \lesssim f \| .
\end{align*}
$$

The corresponding estimates for $f \in E_{B_{0}}^{-} \mathcal{H}$ also hold.
Proof. The proof of (4.19) uses Schur estimates. Applying Cauchy-Schwarz inequality, we estimate the square of the left-hand side by

$$
\begin{aligned}
\int_{0}^{\infty}\left\|\int_{0}^{t} \frac{s}{t}\left(s \partial_{s} F_{s}\right) \frac{d s}{s}\right\|^{2} \frac{d t}{t} & \leqslant \int_{0}^{\infty}\left(\int_{0}^{t} \frac{s}{t} \frac{d s}{s}\right)\left(\int_{0}^{t} \frac{s}{t}\left\|s \partial_{s} F_{s}\right\|^{2} \frac{d s}{s}\right)^{2} \frac{d t}{t} \\
& =\int_{0}^{\infty}\left(\int_{s}^{\infty} \frac{s}{t} \frac{d t}{t}\right)\left\|s \partial_{s} F_{s}\right\|^{2} \frac{d s}{s}=\| \| \psi_{t}\left(T_{B_{0}}\right) f\left\|^{2} \lesssim\right\| f \|^{2}
\end{aligned}
$$

where $\psi(z):=z e^{-|z|}$.

To prove (4.20), we note that $i M_{B_{0}} T_{B_{0}}=\underline{d}+\underline{d}_{B_{0}}^{*}$ and write $\mathbf{P}_{B_{0}}^{1}, \mathbf{P}_{B_{0}}^{2}$ for the two Hodge projections corresponding to the splitting $\mathcal{H}=\mathrm{N}(\underline{d}) \oplus \mathrm{N}\left(\underline{d}_{B_{0}}^{*}\right)$. Note that these projections are bounded by Lemma 2.21(i). We get

$$
t \underline{d} F_{t}=i \mathbf{P}_{B_{0}}^{1} M_{B_{0}}\left(t T_{B_{0}} F_{t}\right)=i \mathbf{P}_{B_{0}}^{1} M_{B_{0}} \psi_{t}\left(T_{B_{0}}\right) f
$$

from which (4.20) follows.
Proof that Corollary 4.15 implies the estimate (4.14). By inserting $I=P_{t}^{\hat{B}_{0}}+\left(I-P_{t}^{\hat{B}_{0}}\right)$ we write the left-hand side in (4.14) as

$$
\widetilde{Q}_{t} \mathcal{E} P_{t}^{\hat{B}_{0}}\left(F_{t}-f\right)+\widetilde{Q}_{t} \mathcal{E}\left(I-P_{t}^{\hat{B}_{0}}\right) F_{t}-\widetilde{Q}_{t} \mathcal{E}\left(I-P_{t}^{\hat{B}_{0}}\right) f=: X_{1}+X_{2}-X_{3}
$$

For $X_{3}$ we get from Theorem 4.8 the estimate $\left\|\left\|X_{3}\right\|\right\|\left(\left\|\widetilde{Q}_{t}\right\|_{\text {op }}+1\right)\|\mathcal{E}\|_{\infty}\|f\|$ since $I-P_{t}^{\hat{B}_{0}}=$ $t^{2} \Gamma \Gamma_{\hat{B}_{0}}^{*} P_{t}^{\hat{B}_{0}}+t^{2} \Gamma_{\hat{B}_{0}}^{*} \Gamma P_{t}^{\hat{B}_{0}}$. For the term $X_{2}$ we use the first identity in Lemma 4.14 to obtain

$$
\begin{aligned}
X_{2}= & \widetilde{Q}_{t} \mathcal{E} Q_{t}^{\hat{B}_{0}} t T_{\hat{B}_{0}} F_{t} \\
= & -\widetilde{Q}_{t} \mathcal{E} Q_{t}^{\hat{B}_{0}}\left(t \partial_{t} F_{t}\right)-i \widetilde{Q}_{t} \mathcal{E} Q_{t}^{\hat{B}_{0}} a_{00}^{-1} a_{\perp \|}\left(t \underline{d} F_{t}^{1, \perp}\right) \\
& -i \widetilde{Q}_{t} \mathcal{E}\left(Q_{t}^{\hat{B}_{0}} t \underline{d}_{\hat{B}_{0}}^{*}\right)\left(F_{t}^{1,0}\left(a_{\| \|}^{-1} a_{\| \perp} e_{0}\right)\right)=:-X_{4}-i X_{5}-i X_{6} .
\end{aligned}
$$

We have the estimate $\left\|X_{4}\right\|\|\lesssim\| \mathcal{E}\left\|_{\infty}\right\| t \partial_{t} F_{t}\|\lesssim \lesssim \mathcal{E}\|_{\infty}\|f\|$. For $X_{5}$, we see from Remark 2.6 that $\underline{d} F_{t}^{1, \perp}$ is the vector part of $\underline{d} F_{t}$. Thus $\left\|\left\|X_{5}\right\| \lesssim\right\| \mathcal{E}\left\|_{\infty}\right\| t \underline{d} F_{t}\|\lesssim \lesssim\| \mathcal{E}\left\|_{\infty}\right\| f \|$ by (4.20). To handle the term $X_{6}$ we note that $Q_{t}^{\hat{B}_{0}} t \underline{B}_{\hat{B}_{0}}^{*}=-Q_{t}^{\hat{B}_{0}} t \Gamma_{\hat{B}_{0}}^{*} N=-t^{2} \Gamma \Gamma_{\hat{B}_{0}}^{*} P_{t}^{\hat{B}_{0}} N$ using Remark 3.2. Thus we obtain from Corollary 4.15, with $v=a_{\| \|}^{-1} a_{\| \perp} e_{0}$, the estimate $\left\|\left|X_{6} \|\right| \lesssim\right.$ $\left(\left\|\left\|\widetilde{Q}_{t}\right\|_{\text {op }}+1\right)\|\mathcal{E}\|_{\infty}\|f\|\right.$.

It remains to estimate the term $X_{1}$. To handle this, we separate the normal vector part as

$$
X_{1}=\widetilde{Q}_{t} \mathcal{E} P_{t}^{\hat{B}_{0}}\left(F_{t}^{1,0} e_{0}\right)-\widetilde{Q}_{t} \mathcal{E} P_{t}^{\hat{B}_{0}}\left(f^{1,0} e_{0}\right)+\widetilde{Q}_{t} \mathcal{E} P_{t}^{\hat{B}_{0}}\left(G_{t}-g\right)=: X_{7}-X_{8}+X_{9}
$$

where $G_{t}:=F_{t}-F_{t}^{1, \perp}$ and $g=f-f^{1, \perp}$. From Corollary 4.15, with $v=e_{0}$, we get the estimate $\left\|X_{7}\right\| \lesssim\left(\left\|\left\|\widetilde{Q}_{t}\right\|_{\text {op }}+1\right)\|\mathcal{E}\|_{\infty}\|f\|\right.$. For the term $X_{8}$, we write $P_{t}^{\hat{B}_{0}}=I-t^{2} \Gamma \Gamma_{\hat{B}_{0}}^{*} P_{t}^{\hat{B}_{0}}-$ $t^{2} \Gamma_{\hat{B}_{0}}^{*} \Gamma P_{t}^{\hat{B}_{0}}$. From Theorem 4.8 we obtain the estimate $\left\|X_{8}\right\| \lesssim\left(\left\|\widetilde{Q}_{t}\right\|_{\mathrm{op}}+1\right)\|\mathcal{E}\|_{\infty}\|f\|$. For the term $X_{9}$, we integrate by parts to obtain

$$
X_{9}=\widetilde{Q}_{t} \mathcal{E} P_{t}^{\hat{B}_{0}}\left(t \partial_{t} G_{t}\right)-\widetilde{Q}_{t} \mathcal{E} P_{t}^{\hat{B}_{0}}\left(\int_{0}^{t} s \partial_{s}^{2} G_{s} d s\right)=: X_{10}-X_{11}
$$

We have the estimate $\left\|X_{10}\right\|\|\lesssim \mathcal{E}\|_{\infty}\| \| t \partial_{t} F_{t}\|\lesssim\| \mathcal{E}\left\|_{\infty}\right\| f \|$. For the term $X_{11}$ we apply $\partial_{t}$ to the last expression for $\left(\partial_{t}+T_{\hat{B}_{0}}\right) F_{t}$ in Lemma 4.14 and get

$$
\begin{equation*}
\partial_{s}^{2} G_{s}+\partial_{s} T_{\hat{B}_{0}} F_{s}=i{\underline{d_{\hat{B}_{0}}}}_{*}^{\partial_{s}} F_{s}^{1, \|} \tag{4.21}
\end{equation*}
$$

Thus

$$
\begin{aligned}
X_{11} & =-\widetilde{Q}_{t} \mathcal{E}\left(P_{t}^{\hat{B}_{0}} t T_{\hat{B}_{0}}\right)\left(\int_{0}^{t} \frac{s}{t}\left(s \partial_{s} F_{s}\right) \frac{d s}{s}\right)+i \widetilde{Q}_{t} \mathcal{E}\left(P_{t}^{\hat{B}_{0}} t \underline{d}_{\hat{B}_{0}}^{*}\right)\left(\int_{0}^{t} \frac{s}{t}\left(s \partial_{s} F_{s}^{1, \|}\right) \frac{d s}{s}\right) \\
& =:-X_{12}+i X_{13} .
\end{aligned}
$$

Both $\left\|\left|X_{12} \|\right|\right.$ and $\| \mid X_{13}\| \|$ can now be estimated with $\|\mathcal{E}\|_{\infty}\|f\|$ by (4.19) since $P_{t}^{\hat{B}_{0}} t T_{\hat{B}_{0}}=Q_{t}^{\hat{B}_{0}}$ and $P_{t}^{\hat{B}_{0}} t \underline{d}_{\hat{B}_{0}}^{*}$ are uniformly bounded by Corollary 2.23 . This proves the estimate (4.14).

Proof that Corollary 4.15 implies the estimate (4.15). We write the left-hand side in (4.15), using integration by parts, as

$$
\widetilde{Q}_{t} \mathcal{E} \underline{d} \int_{0}^{t} F_{s} d s=\widetilde{Q}_{t} \mathcal{E} t \underline{d} F_{t}-\widetilde{Q}_{t} \mathcal{E} \underline{d} \int_{0}^{t} s \partial_{s} F_{s} d s=: X_{1}-X_{2}
$$

For $X_{1}$ we have $\left\|\widetilde{Q}_{t} \mathcal{E}\left(t \underline{d} F_{t}\right)\right\| \lesssim\|\mathcal{E}\|_{\infty}\|f\|$ by (4.20). For $X_{2}$, we write $I=P_{t}^{\hat{B}_{0}}+\left(I-P_{t}^{\hat{B}_{0}}\right)$ and get

$$
X_{2}=\widetilde{Q}_{t} \mathcal{E}\left(t \underline{d} P_{t}^{\hat{B}_{0}}\right) \int_{0}^{t} \frac{s}{t}\left(s \partial_{s} F_{s}\right) \frac{d s}{s}+\widetilde{Q}_{t} \mathcal{E} t \underline{d} Q_{t}^{\hat{B}_{0}} \int_{0}^{t} s \partial_{s} T_{\hat{B}_{0}} F_{s} d s=: X_{3}+X_{4}
$$

Using that $\left\|t \underline{d} P_{t}^{\hat{B}_{0}}\right\| \leqslant C$ by Corollary 2.23, and (4.19) shows that $\left\|\left\|X_{3}\right\|\right\| \lesssim \mathcal{E}\left\|_{\infty}\right\| f \|$. To handle $X_{4}$, we use the identity (4.21), which gives

$$
X_{4}=\widetilde{Q}_{t} \mathcal{E} t \underline{d} Q_{t}^{\hat{B}_{0}}\left(-\int_{0}^{t} s \partial_{s}^{2} G_{s} d s+i \int_{0}^{t} s \underline{d}_{\hat{B}_{0}}^{*} \partial_{s} F_{s}^{1, \|} d s\right)=:-X_{5}+i X_{6}
$$

For $X_{6}$ we note that $\underline{d} Q_{t}^{\hat{B}_{0}} \underline{d}_{\hat{B}_{0}}^{*}=N \underline{d} Q_{t}^{\hat{B}_{0}} N \underline{d}_{\hat{B}_{0}}^{*}=-\Gamma Q_{t}^{\hat{B}_{0}} \Gamma_{\hat{B}_{0}}^{*}=-\Gamma^{2} Q_{t}^{\hat{B}_{0}}=0$ using Remark 3.2, and thus $X_{6}=0$. To handle $X_{5}$, we rewrite this with an integration by parts as

$$
X_{5}=\widetilde{Q}_{t} \mathcal{E}\left(t \underline{d} Q_{t}^{\hat{B}_{0}}\right)\left(t \partial_{t} G_{t}\right)-\widetilde{Q}_{t} \mathcal{E} t \underline{d} Q_{t}^{\hat{B}_{0}} G_{t}+\widetilde{Q}_{t} \mathcal{E}\left(t \underline{d} Q_{t}^{\hat{B}_{0}} g\right)=: X_{6}-X_{7}+X_{8}
$$

where $g=f-f^{1, \perp}$. Using Lemma 4.16 and that $\left\|t \underline{d} Q_{t}^{\hat{B}_{0}}\right\| \leqslant C$ by Corollary 2.23 , we get $\left\|X_{6}\right\| \lesssim\|\mathcal{E}\|_{\infty}\|f\|$. For $X_{8}$, we note that $t \underline{d} Q_{t}^{\hat{B}_{0}}=i N t^{2} \Gamma \Gamma_{\hat{B}_{0}}^{*} P_{t}^{\hat{B}_{0}}$. Thus we can apply Theorem 4.8 to obtain $\left\|X_{8}\right\| \lesssim\left(\left\|\left\|\widetilde{Q}_{t}\right\|_{\text {op }}+1\right)\|\mathcal{E}\|_{\infty}\|f\|\right.$.

We now write $G_{t}=F_{t}-F_{t}^{1,0} e_{0}$ and get

$$
X_{7}=\widetilde{Q}_{t} \mathcal{E} t \underline{d} P_{t}^{\hat{B}_{0}}\left(t T_{\hat{B}_{0}} F_{t}\right)-\widetilde{Q}_{t} \mathcal{E} t \underline{d} Q_{t}^{\hat{B}_{0}}\left(F_{t}^{1,0} e_{0}\right)=: X_{9}-X_{10} .
$$

Again noting that $t \underline{d} Q_{t}^{\hat{B}_{0}}=i N t^{2} \Gamma \Gamma_{\hat{B}_{0}}^{*} P_{t}^{\hat{B}_{0}}$ we obtain from Corollary 4.15, with $v=e_{0}$, the estimate $\left\|\left\|X_{10}\right\|\right\|\left(\left\|\left\|\widetilde{Q}_{t}\right\|_{\text {op }}+1\right)\|\mathcal{E}\|_{\infty}\|f\|\right.$. The term $X_{9}$ remains, on which we use the first identity in Lemma 4.14 to obtain

$$
\begin{aligned}
X_{9}= & -\widetilde{Q}_{t} \mathcal{E} t \underline{d} P_{t}^{\hat{B}_{0}}\left(t \partial_{t} F_{t}\right)-i \widetilde{\widetilde{Q}}_{t} \mathcal{E}\left(t \underline{d} P_{t}^{\hat{B}_{0}}\right) a_{00}^{-1} a_{\perp \|}\left(t \underline{d} F_{t}^{1, \perp}\right) \\
& -i \widetilde{Q}_{t} \mathcal{E}\left(t \underline{t} P_{t}^{\hat{B}_{0}} t \underline{d}_{\hat{B}_{0}}^{*}\right)\left(F_{t}^{1,0}\left(a_{\| \|}^{-1} a_{\| \perp} e_{0}\right)\right)=:-X_{11}-i X_{12}-i X_{13} .
\end{aligned}
$$

Using that $\left\|t \underline{d} P_{t}^{\hat{B}_{0}}\right\| \leqslant C$ by Corollary 2.23 , shows that $\left\|X_{11}\right\| \lesssim\|\mathcal{E}\|_{\infty}\|f\|$. For $X_{12}$ we see from Remark 2.6 that $\underline{d} F_{t}^{1, \perp}$ is the vector part of $\underline{d} F_{t}$. Thus $\left\|X_{12}\right\| \lesssim\|\mathcal{E}\|_{\infty}\| \| t{ }_{t} F_{t}\|\lesssim\| \mathcal{E}\left\|_{\infty}\right\| f \|$ by (4.20). To handle the final term $X_{13}$ we note that $\underline{d} P_{t}^{\hat{B}_{0}} \underline{d}_{\hat{B}_{0}}^{*}=\underline{d} N P_{t}^{\hat{B}_{0}} N \underline{d}_{\hat{B}_{0}}^{*}=-\Gamma P_{t}^{\hat{B}_{0}} \Gamma_{\hat{B}_{0}}^{*}=$ $-\Gamma \Gamma_{\hat{B}_{0}}^{*} P_{t}^{\hat{B}_{0}}$ using Remark 3.2. Thus we obtain from Corollary 4.15, with $v=a_{\| \|}^{-1} a_{\| \perp} e_{0}$, the estimate $\left\|X_{13}\right\| \lesssim\left(\left\|\left\|\widetilde{Q}_{t}\right\|_{\text {op }}+1\right)\|\mathcal{E}\|_{\infty}\|f\|\right.$. This proves the estimate (4.15).

Proof of Corollary 4.15. Let $f \in E_{B_{0}}^{+} \mathcal{H}$ and consider the functions $F_{t}=e^{-t\left|T_{B_{0}}\right|} f$ and $G_{t}:=$ $P_{t}^{B_{0}} f$. We note that $F_{t}-G_{t}=\psi_{t}\left(T_{B_{0}}\right) f$, where $\psi(z)=e^{-|z|}-\left(1+z^{2}\right)^{-1} \in \Psi\left(S_{\nu}^{o}\right)$. Thus it suffices to prove the estimate $\left\|\Theta_{t}^{i}\left(G_{t}^{1,0}\right)\right\| \lesssim\left(\left\|\left\|\widetilde{Q}_{t}\right\|_{\mathrm{op}}+1\right)\|\mathcal{E}\|_{\infty}\|f\|, i=1,2\right.$, for the two families of operators

$$
\begin{align*}
& \Theta_{t}^{1}:=\widetilde{Q}_{t} \mathcal{E} t^{2} \Gamma \Gamma_{\hat{B}_{0}}^{*} P_{t}^{\hat{B}_{0}} M_{v}  \tag{4.22}\\
& \Theta_{t}^{2}:=\widetilde{Q}_{t} \mathcal{E} P_{t}^{\hat{B}_{0}} M_{v} \tag{4.23}
\end{align*}
$$

where $M_{v}$ denotes the multiplication operator $M_{v}(f):=v f$. To this end, we write

$$
\Theta_{t}^{i}\left(G_{t}^{1,0}\right)=\left(\Theta_{t}^{i}-\gamma_{t}^{i} A_{t}\right) G_{t}^{1,0}+\gamma_{t}^{i} A_{t} G_{t}^{1,0}
$$

where $\gamma_{t}^{i}(x)$ is the principal part of the operator family $\Theta_{t}^{i}$ as in Definition 4.2. Using the principal part approximation Lemma 4.4 and Carleson's lemma 4.6 we obtain the estimate

$$
\left\|\Theta_{t}^{i}\left(G_{t}^{1,0}\right)\right\|\|=\| \Theta_{t}^{i}\left\|_{\text {off }}\right\|\left\|t \nabla\left(G_{t}^{1,0}\right)\right\|\|+\| \gamma_{t}^{i}\left\|_{C}\right\| N_{*}\left(A_{t} G_{t}^{1,0}\right) \| .
$$

By Proposition 2.25 and Lemma 2.26, we have $\left\|\Theta_{t}^{i}\right\|_{\text {off }} \lesssim\|\mathcal{E}\|_{\infty}$. Furthermore, by Lemma 4.7 and Theorem 4.8 we have

$$
\left\|\gamma_{t}^{i}\right\|_{C} \lesssim\| \| \Theta_{t}^{i}\left\|_{\mathrm{op}}+\right\| \Theta_{t}^{i}\left\|_{\mathrm{off}} \lesssim\left(\left\|\widetilde{Q}_{t}\right\|_{\mathrm{op}}+1\right)\right\| \mathcal{E} \|_{\infty}
$$

For $\Theta_{t}^{2}$ we have used that $\Theta_{t}^{2}=\widetilde{Q}_{t} \mathcal{E}\left(I-t^{2} \Gamma \Gamma_{\hat{B}_{0}}^{*} P_{t}^{\hat{B}_{0}}-t^{2} \Gamma_{\hat{B}_{0}}^{*} \Gamma P_{t}^{\hat{B}_{0}}\right) M_{v}$.

To bound $\left\|\left\|t \nabla\left(G_{t}^{1,0}\right)\right\|\right.$, we note that the vector part of $\underline{d} G_{t}$ is

$$
\underline{d} G_{t}^{1, \perp}=i m d\left(e_{0} G_{t}^{1,0}\right)=-i d\left(G_{t}^{1,0}\right)=-i \nabla\left(G_{t}^{1,0}\right)
$$

Thus, similar to the proof of (4.20), it follows that

$$
\left\|t \nabla\left(G_{t}^{1,0}\right)\right\|\left\|\left\|t \underline{d} G_{t}\right\| \lesssim\right\| t T_{B_{0}} P_{t}^{B_{0}} f\|\lesssim\| f \|
$$

Finally, to bound $\left\|N_{*}\left(A_{t} G_{t}^{1,0}\right)\right\|$ we write $G_{t}=\frac{1}{2}\left(H_{t}+H_{-t}\right)$, where $H_{t}:=\left(1+i t T_{B_{0}}\right)^{-1} f$. We now observe that the divergence form equation and estimate for $H_{t}^{1,0}$ in Lemma 2.57 in fact holds for all $f \in \mathcal{H}$, by inspection of the proof. We get

$$
\begin{aligned}
N_{*}\left(A_{t} H_{t}^{1,0}\right)(x) & =\sup _{|y-x|<t}\left|\int_{Q(y, t)} H_{t}^{1,0}(z) d z\right| \\
& \lesssim \sup _{t>0}\left(f_{B\left(x, r_{0} t\right)}\left|H_{t}^{1,0}(z)\right|^{q} d z\right)^{1 / q} \lesssim\left(M\left(|f|^{p}\right)(x)\right)^{1 / p},
\end{aligned}
$$

where $Q(y, t)$ denotes the dyadic cube $Q \in \Delta_{t}$ which contains $y$. Using the boundedness of the Hardy-Littlewood maximal function on $L_{2 / p}\left(\mathbf{R}^{n}\right)$, we obtain

$$
\left\|N_{*}\left(A_{t} H_{t}^{1,0}\right)\right\|_{2} \lesssim\left\|\left(M\left(|f|^{p}\right)\right)^{1 / p}\right\|_{2}=\left\|M\left(|f|^{p}\right)\right\|_{2 / p}^{1 / p} \lesssim\left\||f|^{p}\right\|_{2 / p}^{1 / p}=\|f\|_{2}
$$

A similar argument shows that $\left\|N_{*}\left(A_{t} H_{-t}^{1,0}\right)\right\| \lesssim\|f\|$, and thus $\left\|N_{*}\left(A_{t} G_{t}^{1,0}\right)\right\| \lesssim\|f\|$. We have proved that $\left\|\Theta_{t}^{i}\left(G_{t}^{1,0}\right)\right\| \lesssim\left(\left\|\widetilde{Q}_{t}\right\|_{\mathrm{op}}+1\right)\|\mathcal{E}\|_{\infty}\|f\|$, and therefore Corollary 4.15.

### 4.3. Proof of main theorems

We are now in position to prove Theorems 1.1, 1.3 and 1.4 stated in the introduction.
Proof of Theorem 1.4. Given a perturbation $B^{k} \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\bigwedge^{k}\right)\right)$ of the unperturbed coefficients $B_{0}^{k} \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\bigwedge^{k}\right)\right)$, we introduce $B:=I \oplus \cdots \oplus B^{k} \oplus \cdots \oplus I \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}(\bigwedge)\right)$ and $B_{0}:=I \oplus \cdots \oplus B_{0}^{k} \oplus \cdots \oplus I \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}(\bigwedge)\right)$ acting in all $\mathcal{H}$. By Theorem 3.3(i), $T_{B_{0}}$ satisfies quadratic estimates and by Lemma 4.1 and Theorem 4.8 there exists $\varepsilon>0$ such that we have quadratic estimates

$$
\left\|Q_{t}^{B} f\right\| \approx\|f\|, \quad \text { whenever }\left\|B-B_{0}\right\|_{\infty} \leqslant \varepsilon, f \in \mathcal{H}
$$

Therefore, by Propositions 2.32 and 2.42 we have when $\left\|B-B_{0}\right\|_{\infty} \leqslant \varepsilon / 2$ well-defined and bounded operators $E_{B}=\operatorname{sgn}\left(T_{B}\right)$ which depend Lipschitz continuously on $B$, i.e.

$$
\left\|E_{B_{2}}-E_{B_{1}}\right\| \leqslant C\left\|B_{2}-B_{1}\right\|_{\infty}, \quad \text { when }\left\|B_{i}-B_{0}\right\|<\varepsilon / 2, i=1,2
$$

To prove that $\left(\operatorname{Tr}-B^{k} \alpha^{ \pm}\right)$is well posed, note that by Lemma 2.51 it suffices to show that $\lambda$ $E_{B} N_{B}$ is invertible since then in particular $\lambda-E_{B^{k}} N_{B^{k}}=\left.\left(\lambda-E_{B} N_{B}\right)\right|_{\hat{\mathcal{H}}_{B}^{k}}$ is invertible. Here the
spectral parameter is $\lambda=(\alpha+1) /(\alpha-1)$ and $\alpha:=\alpha^{+} / \alpha^{-}$. By Theorem 3.3(ii), the unperturbed operator $\lambda-E_{B_{0}} N_{B_{0}}$ is invertible when $\lambda^{2}+1 \neq 0$. For the perturbed operator we write

$$
\lambda-E_{B} N_{B}=\left(\lambda-E_{B_{0}} N_{B_{0}}\right)\left(I+\frac{1}{\lambda^{2}+1}\left(\lambda-N_{B_{0}} E_{B_{0}}\right)\left(E_{B_{0}} N_{B_{0}}-E_{B} N_{B}\right)\right) .
$$

Here $\left\|E_{B_{0}} N_{B_{0}}-E_{B} N_{B}\right\| \lesssim\left\|B-B_{0}\right\|_{\infty}$ since we clearly have Lipschitz continuity $\| N_{B_{2}}-$ $N_{B_{1}}\|\lesssim\| B_{2}-B_{1} \|_{\infty}$. It follows that $\lambda-E_{B} N_{B}$ is invertible when $\left\|B-B_{0}\right\|_{\infty} \leqslant C\left|\lambda^{2}+1\right|$. Under the assumption $\left\|B-B_{0}\right\|_{\infty}<\varepsilon / 2$ with $\varepsilon$ small we can replace $\lambda^{2}+1=2\left(\alpha^{2}+1\right) /$ $(\alpha-1)^{2}$ with $\alpha^{2}+1$.

This proves that for each boundary function $g$ there exists a unique solution $f=f^{+}+f^{-}$satisfying the jump conditions in $\left(\operatorname{Tr}-B^{k} \alpha^{ \pm}\right)$and $\left\|f^{+}\right\|+\left\|f^{-}\right\| \approx\|f\| \approx\|g\|$. The norm estimates for $F^{ \pm}$in Theorem 1.4 now follows from Lemma 2.49. A formula for the solution is

$$
F^{ \pm}(t, x)=2 e^{\mp t\left|T_{B}\right|} E_{B}^{ \pm}\left(\left(\alpha^{+}+\alpha^{-}\right) E_{B}-\left(\alpha^{+}-\alpha^{-}\right) N_{B}\right)^{-1} g(x)
$$

This completes the proof of Theorem 1.4.
Before turning to the proofs of Theorems 1.1 and 1.3, we note some corollaries of Theorem 1.4. First, if we let $k=1$ in Theorem 1.4, then it proves that the following Neumannregularity transmission problem is well posed for small $L_{\infty}$ perturbations $A$ of a block matrix $A_{0}$.

Transmission problem ( $\mathbf{T r}-\boldsymbol{A} \boldsymbol{\alpha}^{ \pm}$). Let $\alpha^{ \pm} \in \mathbf{C}$ be given jump parameters. Given scalar functions $\psi, \phi: \mathbf{R}^{n} \rightarrow \mathbf{C}$ with $\nabla_{x} \psi \in L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}^{n}\right)$ and $\phi \in L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}\right)$, find gradient vector fields $F^{ \pm}(t, x)=\nabla_{t, x} U^{ \pm}(t, x)$ in $\mathbf{R}_{ \pm}^{n+1}$ such that $F_{t}^{ \pm} \in C^{1}\left(\mathbf{R}_{ \pm} ; L_{2}\left(\mathbf{R}^{n} ; \mathbf{C}^{n+1}\right)\right)$ and $F^{ \pm}$satisfies (1.2) for $\pm t>0$, and furthermore $\lim _{t \rightarrow \pm \infty} F_{t}^{ \pm}=0$ and $\lim _{t \rightarrow 0^{ \pm}} F_{t}^{ \pm}=f^{ \pm}$in $L_{2}$ norm, where the traces $f^{ \pm}$satisfy the jump conditions

$$
\left\{\begin{array}{l}
\alpha^{-} \nabla_{x} U^{+}(0, x)-\alpha^{+} \nabla_{x} U^{-}(0, x)=\nabla_{x} \psi(x) \\
\alpha^{+} \frac{\partial U^{+}}{\partial v_{A}}(0, x)-\alpha^{-} \frac{\partial U^{-}}{\partial v_{A}}(0, x)=\phi(x)
\end{array}\right.
$$

where $\nabla_{x} U^{ \pm}(0, x)=f_{\|}^{ \pm}(x)$ and $\frac{\partial U^{ \pm}}{\partial v_{A}}=\left(A f^{ \pm}, e_{0}\right)$ denotes the conormal derivative.
Secondly, Theorem 1.4 give perturbation results for the following boundary value problems for $k$-vector fields.

Normal BVP (Nor- $\boldsymbol{B}^{\boldsymbol{k}}$ ). Given a $k$-vector field $g \in \hat{\mathcal{H}}_{B}^{k}$, find a $k$-vector field $F(t, x)$ in $\mathbf{R}_{+}^{n+1}$ such that $F_{t} \in C^{1}\left(\mathbf{R}_{+} ; L_{2}\left(\mathbf{R}^{n} ; \bigwedge^{k}\right)\right)$ and $F$ satisfies (1.8) for $t>0$, and furthermore $\lim _{t \rightarrow \infty} F_{t}=0$ and $\lim _{t \rightarrow 0} F_{t}=f$ in $L_{2}$ norm, where $f$ satisfies

$$
\left.\left.e_{0}\right\lrcorner\left(B^{k} f\right)=e_{0}\right\lrcorner\left(B^{k} g\right) \quad \text { on } \mathbf{R}^{n}=\partial \mathbf{R}_{+}^{n+1} .
$$

Tangential BVP (Tan- $\boldsymbol{B}^{k}$ ). Given a $k$-vector field $g \in \hat{\mathcal{H}}_{B}^{k}$, find a $k$-vector field $F(t, x)$ in $\mathbf{R}_{+}^{n+1}$ such that $F_{t} \in C^{1}\left(\mathbf{R}_{+} ; L_{2}\left(\mathbf{R}^{n} ; \bigwedge^{k}\right)\right)$ and $F$ satisfies (1.8) for $t>0$, and furthermore $\lim _{t \rightarrow \infty} F_{t}=0$ and $\lim _{t \rightarrow 0} F_{t}=f$ in $L_{2}$ norm, where $f$ satisfies

$$
e_{0} \wedge f=e_{0} \wedge g \quad \text { on } \mathbf{R}^{n}=\partial \mathbf{R}_{+}^{n+1}
$$

Corollary 4.17. Let $B_{0}^{k}=B_{0}^{k}(x) \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\bigwedge^{k}\right)\right)$ be accretive and assume that $B_{0}^{k}$ is a block matrix. Then there exists $\varepsilon>0$ depending only on the constants $\left\|B_{0}^{k}\right\|_{\infty}$ and $\kappa_{B_{0}^{k}}$ and dimension $n$, such that if $B^{k} \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\bigwedge^{k}\right)\right)$ satisfies $\left\|B^{k}-B_{0}^{k}\right\|_{\infty}<\varepsilon$, then Normal and Tangential boundary value problems (Nor- $\left.B^{k}\right)$ and $\left(\operatorname{Tan}-B^{k}\right)$ above are well posed.

Proof. (i) For (Nor- $\left.B^{k}\right)$ in $\mathbf{R}_{+}^{n+1}$, we let $\alpha^{-}=0$ and $\alpha^{+}=1$ in ( $\left.\operatorname{Tr}-B^{k} \alpha^{ \pm}\right)$. Then we obtain two decoupled jump conditions

$$
\left\{\begin{array}{l}
-e_{0} \wedge f^{-}=e_{0} \wedge g \\
\left.\left.e_{0}\right\lrcorner\left(B^{k} f^{+}\right)=e_{0}\right\lrcorner\left(B^{k} g\right) .
\end{array}\right.
$$

Discarding the solution $F^{-}$, we obtain a unique solution $F=F^{+}$to (Nor- $B^{k}$ ).
(ii) For (Tan- $B^{k}$ ) in $\mathbf{R}_{+}^{n+1}$, we let $\alpha^{-}=1$ and $\alpha^{+}=0$ in $\left(\operatorname{Tr}-B^{k} \alpha^{ \pm}\right)$. Then we obtain two decoupled jump conditions

$$
\left\{\begin{array}{l}
e_{0} \wedge f^{+}=e_{0} \wedge g \\
\left.\left.-e_{0}\right\lrcorner\left(B^{k} f^{-}\right)=e_{0}\right\lrcorner\left(B^{k} g\right) .
\end{array}\right.
$$

Discarding the solution $F^{-}$, we obtain a unique solution $F=F^{+}$to $\left(\operatorname{Tan}-B^{k}\right)$.
Note that well-posedness of $(\operatorname{Neu}-A)$ and $(\operatorname{Reg}-A)$ for Theorem 1.1(b) is the special case $k=1$ of Corollary 4.17, as well as the special cases $\left(\alpha^{+}, \alpha^{-}\right)=(1,0)$ and $(0,1)$, respectively, of $\left(\operatorname{Tr}-A \alpha^{ \pm}\right)$.

Proof of Theorem 1.3. Let $A_{0}$ be such that $T_{B_{0}}$ has quadratic estimates in $\mathcal{H}$, where $B_{0}=$ $I \oplus A_{0} \oplus I \oplus \cdots \oplus I$. Thus by Lemma 4.1 and Theorem 4.12 there exists $\varepsilon>0$ such that we have quadratic estimates

$$
\left\|Q_{t}^{B} f\right\| \approx\|f\|, \quad \text { whenever }\left\|B-B_{0}\right\|_{\infty}<\varepsilon, f \in \mathcal{H}
$$

Therefore, by Proposition 2.32, we have when $\left\|B-B_{0}\right\|_{\infty}<\varepsilon$ well defined and bounded operators $E_{B}=\operatorname{sgn}\left(T_{B}\right)$. With Lemma 2.46, these restricts to bounded operators $E_{A}$ in $\hat{\mathcal{H}}^{1}$. In particular $f \in \hat{\mathcal{H}}^{1}$ can be decomposed as $f=f^{+}+f^{-}$, where $f^{ \pm}:=E_{A}^{ \pm} f$ and $\|f\| \approx\left\|f^{+}\right\|+\left\|f^{-}\right\|$. Moreover, Lemma 2.49 and Proposition 2.56 proves the stated norm equivalences for $F^{ \pm}$.

Proof of Theorem 1.1. Given a perturbation $A \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\bigwedge^{1}\right)\right)$ of the unperturbed coefficients $A_{0} \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}\left(\bigwedge^{1}\right)\right)$, which we assume are either of block form, real symmetric or constant, we introduce $B:=I \oplus A \oplus I \oplus \cdots \oplus I \in L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}(\bigwedge)\right)$ and $B_{0}:=I \oplus A_{0} \oplus I \oplus \cdots \oplus I \in$ $L_{\infty}\left(\mathbf{R}^{n} ; \mathcal{L}(\bigwedge)\right)$ acting in all $\mathcal{H}$. That $T_{B_{0}}$ satisfies quadratic estimates follows from Theorem 3.3(i), Theorem 3.12 and Proposition 3.5, respectively. Theorem 1.3 now shows that we have quadratic estimates for $T_{B}$ when $\left\|B-B_{0}\right\|_{\infty}<\varepsilon$. In case $A_{0}$ is a block matrix, we note that this result follows already from Theorem 4.8. By Propositions 2.32 and 2.42 we have when $\left\|B-B_{0}\right\|_{\infty}<\varepsilon / 2$ well-defined and bounded operators $E_{B}=\operatorname{sgn}\left(T_{B}\right)$ which depend Lipschitz continuously on $B$, so that

$$
\left\|E_{B_{2}}-E_{B_{1}}\right\| \leqslant C\left\|B_{2}-B_{1}\right\|_{\infty}, \quad \text { when }\left\|B_{i}-B_{0}\right\|<\varepsilon / 2, i=1,2 .
$$

To prove that (Neu- $A$ ) and ( $\operatorname{Reg}-A$ ) are well posed, note that by Lemma 2.53 it suffices to show that $I \pm E_{A} N_{A}: \hat{\mathcal{H}}^{1} \rightarrow \hat{\mathcal{H}}^{1}$ are invertible, where $E_{A}=\left.E_{B}\right|_{\hat{\mathcal{H}}^{1}}$ and $N_{A}=\left.N_{B}\right|_{\hat{\mathcal{H}}^{1}}$. By Theorems 3.3 (ii), 3.12 and 3.6, respectively, the unperturbed operators $I \pm E_{A_{0}} N_{A_{0}}: \hat{\mathcal{H}}^{1} \rightarrow \hat{\mathcal{H}}^{1}$ are invertible. For the perturbed operator we write

$$
I \pm E_{A} N_{A}=\left(I \pm E_{A_{0}} N_{A_{0}}\right)\left(I \pm\left(I \pm E_{A_{0}} N_{A_{0}}\right)^{-1}\left(E_{A} N_{A}-E_{A_{0}} N_{A_{0}}\right)\right)
$$

Here $\left\|E_{A} N_{A}-E_{A_{0}} N_{A_{0}}\right\| \lesssim\left\|A-A_{0}\right\|_{\infty}$ since we clearly have Lipschitz continuity $\| N_{A_{2}}-$ $N_{A_{1}}\|\lesssim\| A_{2}-A_{1} \|_{\infty}$. It follows that $I \pm E_{A} N_{A}$ are invertible when $\left\|A-A_{0}\right\|_{\infty}<\varepsilon^{\prime}$.

The well-posedness of $\left(\mathrm{Neu}^{\perp}-A\right)$ is a consequence of Proposition 2.52 , since our hypothesis is stable when taking adjoints $A \mapsto A^{*}$. Alternatively, we can replace $N_{A_{0}}$ with $N$ above, proving that $I \pm E_{A} N$ is an isomorphism, using Theorem 3.3(ii), Remarks 3.13 and 3.7. This proves that for (Neu- $A$ ), (Reg- $A$ ) and ( $\mathrm{Neu}^{\perp}-A$ ) and each boundary function $g$ (being $\phi$ and $\nabla_{x} \psi$, respectively), there exists a unique solution $f=E_{A}^{+} \hat{\mathcal{H}}^{1}$ satisfying the boundary condition and $\|f\| \approx\|g\|$. Lemma 2.49 and Proposition 2.56 proves that the stated norms of $f$ and $F(t, x):=$ $\left(e^{-t\left|T_{A}\right|} f\right)(x)$ are equivalent.

The well-posedness of ( $\operatorname{Dir}-A$ ), as well as the first three norm estimates, follows from Lemma 2.55 and $\left(\mathrm{Neu}^{\perp}-A\right)$. To show that $\left\|t \nabla_{x} U_{t}\right\|\left\|\left\|t \partial_{t} U_{t}\right\|\right\|$, we consider the gradient vector field $G_{t}=\nabla_{t, x} U_{t}$ as in the proof of the lemma. From (Reg-A) and ( $\mathrm{Neu}^{\perp}-A$ ), it follows that for all $t>0$, we have $\left\|\partial_{t} U_{t}\right\|=\left\|N^{-} G_{t}\right\| \approx\left\|N^{+} G_{t}\right\|=\left\|\nabla_{x} U_{t}\right\|$, from which the square function estimate for $\nabla_{x} U_{t}$ follows. To show that $\|u\| \approx\left\|\widetilde{N}_{*}(U)\right\|$, we consider the vector field $F_{t}$ of conjugate functions from the proof of the lemma. Proposition 2.56 shows that $\|u\| \approx\|f\| \approx\left\|\widetilde{N}_{*}(F)\right\| \gtrsim\left\|\widetilde{N}_{*}(U)\right\|$. Moreover, the proof of the reverse estimate $\|u\| \lesssim\left\|\widetilde{N}_{*}(U)\right\|$ is similar to the proof of $\|f\| \lesssim\left\|\widetilde{N}_{*}(F)\right\|$ in Proposition 2.56 , using the uniform boundedness of $\mathcal{P}_{t}$.

Finally we note that the solution operators for $(\mathrm{Neu}-A),\left(\mathrm{Neu}^{\perp}-A\right),(\operatorname{Reg}-A)$ and $(\operatorname{Dir}-A)$ are

$$
\begin{array}{ll}
F_{t}=2 e^{-t\left|T_{A}\right|}\left(E_{A}-N_{A}\right)^{-1}\left(a_{00}^{-1} \phi e_{0}\right), & F_{t}=2 e^{-t\left|T_{A}\right|}\left(E_{A}-N\right)^{-1}\left(\phi e_{0}\right), \\
F_{t}=2 e^{-t\left|T_{A}\right|}\left(E_{A}+N\right)^{-1}\left(\nabla_{x} \psi\right), & U_{t}=2\left(e^{-t\left|T_{A}\right|}\left(E_{A}-N\right)^{-1}\left(u e_{0}\right), e_{0}\right)
\end{array}
$$

The Lipschitz continuity of $u \mapsto U$ is a consequence of the corresponding result for $\left(\mathrm{Neu}^{\perp}-A\right)$. For the norms $\sup _{t>0}\left\|F_{t}\right\|$ and $\left\|t \partial_{t} F_{t}\right\| \|$, Lipschitz continuity for the solution operators follows from Proposition 2.42. It remains to show Lipschitz continuity for the norm $\|F\|_{\mathcal{X}}=\left\|\widetilde{N}_{*}(F)\right\|_{2}$. To this end, we consider

$$
f(x) \mapsto F_{A_{z}}=\left(e^{-t\left|T_{A_{z}}\right|} E_{A_{z}}^{+} f\right)(x): \hat{\mathcal{H}}^{1} \rightarrow \mathcal{X}
$$

and the truncations $f(x) \mapsto F_{A_{z}}^{k}(t, x)=\chi_{k}(t) F_{A_{z}}(t, x)$, where $\chi_{k}$ denotes the characteristic function for $(1 / k, k)$ as in the proof of Lemma 2.41(iii). We claim that it suffices to show that, for each fixed $k$, the operator $f(x) \mapsto F_{A_{z}}^{k}(t, x): \hat{\mathcal{H}}^{1} \rightarrow \mathcal{X}$ depends holomorphically on $z$. Indeed, using Schwarz' lemma as in Proposition 2.42, we obtain the Lipschitz estimate

$$
\left\|F_{A_{2}}^{k}-F_{A_{1}}^{k}\right\|_{\mathcal{X}} \leqslant C\left\|A_{2}-A_{1}\right\|_{\infty}\|f\|_{2}
$$

uniformly for all $k$, since $\left\|F^{k}\right\| \mathcal{X} \leqslant\|F\| \mathcal{X} \leqslant C\|f\|$ by Proposition 2.56. Furthermore, by the monotone convergence theorem we have

$$
\left\|F_{A_{2}}^{k}-F_{A_{1}}^{k}\right\|_{\mathcal{X}} \nearrow\left\|F_{A_{2}}-F_{A_{1}}\right\|_{\mathcal{X}}, \quad k \rightarrow \infty
$$

so the desired Lipschitz continuity follows after taking limits.
To prove that $f(x) \mapsto F_{A_{z}}^{k}(t, x)$ is holomorphic, we note that

$$
\begin{aligned}
\left\|F^{k}\right\|_{\mathcal{X}}^{2} & \lesssim \int_{\mathbf{R}^{n}}\left(\left.\sup _{t>0} \int_{|s-t|<c_{0} t}\left|F_{|y-x|<c_{1} s /\left(1-c_{0}\right)}\right| F^{k}(s, y)\right|^{2} \frac{d s d y}{s^{n+1}}\right) d x \\
& \lesssim \int_{\mathbf{R}^{n}}\left(\int_{1 / k}^{k} \int_{|y-x|<c_{1} s /\left(1-c_{0}\right)}\left|F^{k}(s, y)\right|^{2} \frac{d s d y}{s^{n+1}}\right) d x \approx \int_{1 / k}^{k}\left\|F_{s}^{k}\right\|_{2}^{2} d s
\end{aligned}
$$

for fixed $k$, where in the last step we use that $1 / k \leqslant s^{-1} \leqslant k$. Since $f(x) \mapsto F^{k}(t, x): \hat{\mathcal{H}}^{1} \rightarrow$ $L_{2}\left(\mathbf{R}^{n} \times(1 / k, k)\right)$ is holomorphic by Lemma 2.41(ii) and the embedding $L_{2}\left(\mathbf{R}^{n} \times(1 / k, k)\right) \hookrightarrow$ $\mathcal{X}$ is continuous and independent of $z$, it follows that $f(x) \mapsto F^{k}(t, x): \hat{\mathcal{H}}^{1} \rightarrow \mathcal{X}$ is holomorphic for each fixed $k$. This completes the proof of Theorem 1.1.

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