FUNDAMENTAL RESULTS FOR
PSEUDO-DIFFERENTIAL OPERATORS OF TYPE 1,1

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ABSTRACT. This paper develops some deeper consequences of an extended definition, proposed previously by the author, of pseudo-differential operators that are of type 1,1 in Hörmander’s sense. Thus, it contributes to the long-standing problem of creating a systematic theory of such operators. It is shown that type 1,1-operators are defined and continuous on the full space of temperate distributions, if they fulfil Hörmander’s twisted diagonal condition, or more generally if they belong to the self-adjoint subclass; and that they are always defined on the temperate smooth functions. As a main tool the paradifferential decomposition is derived for type 1,1-operators, and to confirm a natural hypothesis the symmetric term is shown to cause the domain restrictions; whereas the other terms are shown to define nice type 1,1-operators fulfilling the twisted diagonal condition. The decomposition is analysed in the type 1,1-context by combining the Spectral Support Rule and the factorisation inequality, which gives pointwise estimates of pseudo-differential operators in terms of maximal functions.

1. INTRODUCTION

Pseudo-differential operators $a(x,D)$ of type 1,1 have long been known to have peculiar properties, almost since their invention by Hörmander [16]. This is due to initial investigations in 1972 in the thesis of Ching [10] and in lecture notes of Stein (made available in [42, Ch. VII§1.3]); and again in 1978 by Parenti and Rodino [39].

The understanding of their unusual theory, and of the applications of these linear operators to non-linear problems in partial differential equations, grew crucially in the 1980’s through works of Meyer [35, 36], Bony [2], Bourdaud [5, 6, 8, 7], Hörmander [18, 19]. Cf. also the expositions of Hörmander [20, Ch. 9] and Taylor [44].

However, the first formal definition of general type 1,1-operators was put forward in 2008 by the author [26]. It would not be unjust to view this as an axiomatization of the type 1,1-theory, for whereas the previous contributions did not attempt to crystallise what a type 1,1-operator is or how it can be characterised in general, the definition from [26] has been a fruitful framework for raising questions and seeking answers about type 1,1-operators.

Indeed, being based on an operator theoretical approach, mimicking unbounded operators in Hilbert space, the definition gave from the outset a rigorous discussion of, e.g., unclosability, pseudo-locality, non-preservation of wavefront sets and the Spectral Support Rule [26]. This was
followed up with a systematic $L_p$-theory of type 1,1-operators in [29], where a main theorem relied on a symbol analysis proved in full detail in the present paper.

Meanwhile, Métivier also treated type 1,1-operators in 2008 in Chapter 4 of [34], but took recourse to the space dependent extensions of Stein [42, Ch. VII § 1.3]. Type 1,1-operators have also been investigated, or played a role, in works of e.g. Torres [46], Marschall [33], Grafakos and Torres [14], Taylor [45], Hérau [15], Lannes [32], Johnsen [25], Hounie and dos Santos Kapp [21]; and for bilinear operators in Bernicot and Torres [1]. Implicitly type 1,1-operators also enter many works treating partial differential equations with Bony’s paradifferential calculus; but this would lead too far to recall here.

The present paper goes into a deeper, systematic study of type 1,1-operators on $\mathcal{S}'(\mathbb{R}^n)$ and its subspaces. Indeed, the definition in [26] is shown here to give operators always defined on the maximal smooth subspace $C^\infty \cap \mathcal{S}'$, generalising results of Bourdaud [8] and David and Journé [12] — and shown to be defined on the entire $\mathcal{S}'$ if they belong to the self-adjoint subclass, by an extension of Hörmander’s analysis of this class [18, 19]. Moreover, the pointwise estimates in [28] are applied to the paradifferential decompositions, which are analysed in the type 1,1-context here. The decomposition gives 3 other type 1,1-operators, of which the so-called symmetric term is responsible for the possible domain restrictions, which occur when its infinite series diverges.

Altogether this should bring the theory of type 1,1-operators to a rather more mature level.

1.1. **Background.** Recall that the symbol $a(x, \eta)$ of a type 1,1-operator of order $d \in \mathbb{R}$ fulfills

$$|D_\eta^\alpha D_\xi^\beta a(x, \eta)| \leq C_{\alpha, \beta} (1 + |\eta|)^{d - |\alpha| + |\beta|} \quad \text{for} \quad x, \eta \in \mathbb{R}^n. \quad (1.1)$$

Classical pseudo-differential operators are, e.g., partial differential operators $\sum_{|\alpha| \leq d} a_\alpha(x)D^\beta$, having such symbols simply with $d - |\alpha|$ as exponents. The presence of $|\beta|$ allows for a higher growth with respect to $\eta$, which has attracted attention for a number of reasons.

The operator corresponding to (1.1) is for Schwartz functions $u(x)$, i.e. $u \in \mathcal{S}(\mathbb{R}^n)$,

$$a(x,D)u = (2\pi)^{-n} \int e^{ix\cdot \eta} a(x, \eta) \hat{u}(\eta) \, d\eta. \quad (1.2)$$

But for $u \in \mathcal{S}'(\mathbb{R}^n) \setminus \mathcal{S}$ it requires another definition to settle whether $u$ belongs to the domain of $a(x,D)$ or not. This is indeed a main subject of the present paper, which exploits the general definition of $a(x,D)$ presented in [26]; it is recalled in (1.10) below.

The non-triviality of the above task was discovered already by Ching [10], who showed unboundedness on $L_2$ for certain $a_\theta(x,D)$ with $d = 0$; cf. Example 2.5 below. As the adjoint $a_\theta(x,D)^*$ of Ching’s operator does not leave $\mathcal{S}$ invariant, as can be seen explicitly e.g. from the proof of [26, Lem. 3.1], the usual extension to $\mathcal{S}'$ by duality is not possible for $\text{OP}(S_{1,1})$.

In general the pathologies of type 1,1-operators are without doubt reflecting that, most interestingly, this operator class has important applications to non-linear problems:

This was first described around 1980 by Meyer [35, 36], who discovered that a composition operator $u \mapsto F \circ u = F(u)$ with $F \in C^\infty$, $F(0) = 0$, can be decomposed in its action on functions $u \in \bigcup_{k \geq n/p} H^k_p(\mathbb{R}^n)$, by means of a specific $u$-dependent type 1,1 symbol $a_u(x, \eta) \in S_{1,1}^0$, as

$$F(u(x)) = a_u(x,D)u(x). \quad (1.3)$$
He also showed that $a_u(x,D)$ extends to a bounded operator on $H^s_t$ for $t > 0$, so the fact that $u \mapsto F(u)$ sends $H^s_p$ into itself can be seen from (1.3) by taking $t = s$ and $r = p$ — indeed, this proof method is particularly elegant for non-integer $s > n/p$. It was carried over rigorously to the present type 1, 1-framework in [26, Sec. 9], with continuity of $u \mapsto F \circ u$ as a corollary. Some applications of (1.3) were explained by Taylor [44, Ch. 3].

Secondly, it was shown in [36] that type 1, 1-operators play a main role in the paradifferential calculus of Bony [2] and the microlocal inversion of nonlinear partial differential equations of the form

$G(x, (D^\alpha_x u(x))_{|\alpha| \leq m}) = 0.$  \hfill (1.4)

This was explicated by Hörmander, who devoted [20, Ch. 10] to the subject. The resulting setup was used e.g. by Hérau [15] in a study of hypoellipticity of (1.4). Moreover, it was used for propagation of singularities in [20, Ch. 11], with special emphasis on non-linear hyperbolic equations. Recently paradifferential operators, and thus type 1, 1-operators, were also exploited for non-linear Schrödinger operators in constructions of solutions, parametrices and propagation of singularities in global wave front sets; cf. works of e.g. Tataru [43], Delort [13], Nicola and Rodino [38].

Thirdly, both type 1, 1-theory as such and Bony’s paradifferential techniques played a crucial role in the author’s work on semi-linear elliptic boundary problems [25].

Because of the relative novelty of this application, a sketch is given using a typical example. In a bounded $C^\infty$-region $\Omega \subset \mathbb{R}^n$ with normal derivatives $\gamma_j u = (\vec{n} \cdot \nabla)^j u$ at the boundary $\partial \Omega$, and $\Delta := \partial^2_{x_1} + \cdots + \partial^2_{x_n}$, let $u(x)$ solve the perturbed $\ell$-harmonic Dirichlet problem

$(-\Delta)^\ell u + u^2 = f$ in $\Omega$, \quad $\gamma_j u = \varphi_j$ on $\partial \Omega$ for $0 \leq j < \ell$. \hfill (1.5)

Without $u^2$, the linear problem has a well-known solution $u_0 = R_\ell f + K_0 \varphi_0 + \cdots + K_{\ell-1} \varphi_{\ell-1}$, with operators belonging to the pseudo-differential boundary operator class of Boutet de Monvel [9]. For the non-linear problem in (1.5), the parametrix construction of [25] yields the solution formula

$u = P_u^{(N)}(R_\ell f + K_0 \varphi_0 + \cdots + K_{\ell-1} \varphi_{\ell-1}) + (R_\ell L_u)^N u$, \hfill (1.6)

where the parametrix $P_u^{(N)}$ is the linear map given by the finite Neumann series

$P_u^{(N)} = I + R_\ell L_u + \cdots + (R_\ell L_u)^{N-1}$. \hfill (1.7)

in terms of the exact paralinearisation $L_u$ of $u^2$ with the sign convention $-L_u(u) = u^2$; cf. [25].

One merit of (1.6) is to show why $u$’s regularity is unchanged by the non-linear term $u^2$: each parametrix $P_u^{(N)}$ is of order 0, hence does not change Sobolev regularity when applied to $u_0$; while in (1.6) the remainder $(R_\ell L_u)^N u$ will be in $C^k(\overline{\Omega})$ for every fixed $k$ if $N$ is taken large enough. Indeed, $R_\ell L_u$ has a fixed negative order if $u$ is given with just the weak a priori regularity necessary to make sense of the boundary condition and make $u^2$ defined and a priori more regular than $(-\Delta)^\ell u$.

Type 1, 1-operators are important for the fact that (1.6) easily implies that extra regularity properties of $f$ in subregions $\Xi \subset \Omega$ carry over to $u$; e.g. if $f|_\Xi$ is $C^\infty$ so is $u|_\Xi$. Indeed, such
implications boil down to the fact that the exact paralinearisation $L_{\Omega}$ factors through an operator $A_u$ of type 1, 1, that is, if $r_{\Omega}$ denotes restriction to $\Omega$ and $\ell_{\Omega}$ is a linear extension to $\mathbb{R}^n \setminus \Omega$,

$$L_{\Omega} = r_{\Omega}A_u\ell_{\Omega}, \quad A_u \in \text{OP}(S^{\alpha}_{1,1}).$$

(1.8)

Now, by inserting (1.8) into (1.7)–(1.6) for a large $N$ and using cut-off functions supported in $\Xi$ in a well-known way, cf. [25, Thm. 7.8], the regularity of $u$ locally in $\Xi$ is at once improved to the extent permitted by the data $f$ by using the pseudo-local property of $A_u$:

$$\text{sing supp} Au \subset \text{sing supp} u \quad \text{for} \quad u \in D(A).$$

(1.9)

However, the pseudo-local property of general type 1, 1-operators was only proved recently in [26], inspired by the application below (1.8). Yet, pseudo-locality was anticipated more than three decades ago by Parenti and Rodino [39], who gave an inspiring but incomplete indication, as they relied on the future to bring a specific meaning to $a(x,D)u$ for $u \in \mathcal{S}' \setminus C^\infty_0$.

A rigorous definition of general type 1, 1 operators was first given in [26]. In a way the definition abandons Fourier analysis (temporarily) and mimicks the theory of unbounded operators in Hilbert spaces. This is by viewing a type 1 operators by an operator theoretical construction of $a$.

$$a(x,D) := \lim_{m \to \infty} (2\pi)^{-n} \int_{\mathbb{R}^n} e^{ix\cdot\eta} \psi(2^{-m}D_x) a(x,\eta) \psi(2^{-m}\eta) \hat{u}(\eta) d\eta$$

(1.10)

whenever this limit does exist in $\mathcal{D}'(\mathbb{R}^n)$ for every $\psi \in C^\infty_0(\mathbb{R}^n)$ with $\psi = 1$ in a neighbourhood of the origin, and does not depend on such $\psi$.

In passing it is noted that, beyond the definition, operator theory is also felt in the rules of calculus, since as shown in Proposition 2.4 below the well-known commutator identity is replaced for type 1, 1-operators by an operator theoretical inclusion,

$$a(x,D)D_j + [D_j, a(x,D)] \subset D_j a(x,D).$$

(1.11)

The unconventional definition in (1.10), by vanishing frequency modulation, is a rewriting of the usual one, which is suitable for the present general symbols: clearly (1.10) gives back the integral in (1.2) if $u \in \mathcal{S}$. In case $a \in S^{\alpha}_{1,0}$ this identification extends further to $u \in \mathcal{S}'$ by duality and the calculus of classical pseudo-differential operators. Note that the above integral should be interpreted as the operator $\text{OP}(\psi(2^{-m}D_x)a(x,\eta)\psi(2^{-m}\eta))$ in $\text{OP}(S^{-\infty})$ acting on $u$.

Clearly (1.10) is reminiscent of oscillatory integrals, now with the addition that $u \in D(a(x,D))$ when the regularisation yields a limit independent of the integration factor. Of course it is not a conventional integration factor that is used here, but rather the Fourier multiplier $\psi(2^{-m}D_x)$ that modifies the frequencies of $a(\cdot, \eta)$. While the necessity of this modification was amply elucidated in [26], it is moreover beneficial because the use of $\psi(2^{-m}D_x)$ gives easy access to Littlewood–Paley analysis of $a(x,D)$.

The definition (1.10) was investigated in [26] from several other perspectives, of which some will be needed below. But mentioned briefly (1.10) was proved to be maximal among the definitions of $A = a(x,D)$ that gives back the usual operators in $\text{OP}(S^{-\infty})$ and is stable under the limit in (1.10); $A$ is always defined on $\mathcal{F}^{-1}\mathcal{D}'$; it is pseudo-local but does change wavefront sets.
in certain cases (even if $A$ is defined on $\bigcup H^s$); and $A$ transports supports via the distribution kernel, i.e. $\text{supp} A u \subset \text{supp} K \circ \text{supp} u$ when $u \in D(A) \cap \mathcal{E}'$, with a similar spectral support rule for $\text{supp} u$; cf. (1.24) below and Appendix B.

For the Weyl calculus, Hörmander [18] noted that type 1, 1-operators do not fit well, as Ching’s operator can have discontinuous Weyl-symbol. Conversely Boulkhemair [3, 4] showed that the Weyl operator $\int e^{i(x-y) \cdot \eta} a(x, y) u(y) dy \eta / (2\pi)^n$ may give peculiar properties by insertion of $a(x, \eta)$ from $S^1_{1,1}$. E.g., already for Ching’s symbol with $d = 0$, the real or imaginary part gives a Weyl operator that is unbounded on $H^s$ for every $s \in \mathbb{R}$.

For more remarks on the subject’s historic development the reader may refer to Section 2; or consult the review in the introduction of [26].

1.2. Outline of Results. The purpose of this paper is to continue the foundational study in [26] and support the definition in (1.10) with further consequences.

First of all this concerns the hitherto untreated question: under which conditions is a given type 1, 1-operator $a(x, D)$ an everywhere defined and continuous map

$$a(x, D) : \mathcal{S}'(\mathbb{R}^n) \to \mathcal{S}'(\mathbb{R}^n)?$$

(1.12)

For this it is shown in Proposition 4.2 and Theorem 6.3 below to be sufficient that $a(x, \eta)$ fulfils Hörmander’s twisted diagonal condition, i.e. the partially Fourier transformed symbol

$$\hat{a}(\xi, \eta) = \mathfrak{F}_{x \to \xi} a(x, \eta)$$

(1.13)

should vanish in a conical neighbourhood of a non-compact part of the twisted diagonal $\mathcal{T}$ given by $\xi + \eta = 0$ in $\mathbb{R}^n \times \mathbb{R}^n$. More precisely this means that for some $B \geq 1$

$$\hat{a}(\xi, \eta) \neq 0 \quad \text{only if} \quad |\xi + \eta + 1| \geq |\eta|/B.$$  

(1.14)

It should perhaps be noted that it is natural to consider $\hat{a}(\xi, \eta)$, as it is related (cf. [26, Prop. 4.2]) both to the kernel $K$ of $a(x, D)$ and to the kernel $\mathcal{K}$ of $\mathcal{F}^{-1} a(x, D) \mathcal{F}$,

$$(2\pi)^n \mathcal{K}(\xi, \eta) = \hat{a}(\xi - \eta, \eta) = \mathcal{F}(x, y) \rightarrow (\xi, \eta) K(x, -y).$$

(1.15)

More generally the $\mathcal{S}'$-continuity (1.12) is obtained in Theorems 4.6 and 6.5 below for the $a(x, \eta)$ in $S^d_{1,1}$ that merely satisfy Hörmander’s twisted diagonal condition of order $\sigma$ for all $\sigma \in \mathbb{R}$. These are the symbols which for some $c_{a, \sigma}$ and $0 < \varepsilon < 1$ fulfil

$$\sup_{x \in \mathbb{R}^n} R^{\alpha - d/2} \left( \int_{|R| \leq |\eta| \leq 2R} |D^\alpha \eta a_{\chi, \varepsilon}(x, \eta)|^2 \frac{d \eta}{R^n} \right)^{1/2} \leq c_{a, \sigma} \varepsilon^\sigma n^{2 - |\alpha|}.$$  

(1.16)

In this asymptotic formula $a_{\chi, \varepsilon}$ denotes a specific localisation of $\hat{a}(x, \eta)$ to the conical neighbourhood $|\xi + \eta| + 1 \leq 2 \varepsilon |\eta|$ of the twisted diagonal $\mathcal{T}$.

Details on the cut-off function $\chi$ in (1.16) are recalled in Section 2.3, in connection with an account of Hörmander’s fundamental result that validity of (1.16) for all $\sigma \in \mathbb{R}$ is equivalent to extendability of $a(x, D)$ to a bounded map $H^{s+d} \rightarrow H^s$ for all $s \in \mathbb{R}$, as well as equivalent to the adjoint $a(x, D)^*$ being of type 1, 1.

Of course these results of Hörmander make it natural to expect that the above two conditions (namely (1.14) and (1.16) for all $\sigma$) are sufficient for the $\mathcal{S}'$-continuity in (1.12), but this has not
been addressed explicitly in the literature before. As mentioned they are verified in Theorem 6.3,
respectively in Theorem 4.6 by duality and in Theorem 6.5 by exploiting (1.16) directly.

In the realm of smooth functions the situation is fundamentally different. Here there is a
commutative diagram for every type 1,1-operator \( a(x,D) \):

\[
\begin{array}{cccc}
\mathcal{S} & \longrightarrow & \mathcal{S} + \mathcal{S}^{-1} & \longrightarrow \mathcal{O}_M & \longrightarrow \mathcal{C}^\infty \cap \mathcal{S}' \\
\downarrow a(x,D) & & \downarrow a(x,D) & & \downarrow a(x,D) \\
\mathcal{S} & \longrightarrow & \mathcal{O}_M & \longrightarrow \mathcal{O}_M & \longrightarrow \mathcal{C}^\infty
\end{array}
\]

(1.17)

The first column is just the integral (1.2); the second an extension from \([24,26]\). Column three
is an improvement given below of the early contribution of Bourdaud \([8]\) that
\( a(x,D) \) extends to a map \( \mathcal{O}_M \rightarrow \mathcal{D}' \), whereby \( \mathcal{O}_M \) denotes Schwartz’ space of slowly increasing smooth functions.

However, the fourth column restates the full result that a type 1,1-operator is always defined
on the maximal space of smooth functions \( \mathcal{C}^\infty \cap \mathcal{S}' \). More precisely, according to Theorem 2.8
below, it restricts to a strongly continuous map

\[
a(x,D) : \mathcal{C}^\infty(\mathbb{R}^n) \cap \mathcal{S}'(\mathbb{R}^n) \rightarrow \mathcal{C}^\infty(\mathbb{R}^n).
\]

(1.18)

It is noteworthy that this holds without any of the conditions (1.14) and (1.16). Another point is
that, since \( \mathcal{C}^\infty \not\subset \mathcal{S}' \), it was necessary to ask for a limit in the topology of \( \mathcal{D}' \) in (1.10).

Perhaps it could seem surprising that the described results on (1.12) and (1.18) have not been
established in their full generality before. However, it should be emphasised that these properties
are valid for the operator defined in (1.10), so they go much beyond the mere extendability
discussed by Meyer \([36]\), Bourdaud \([7]\), Hörmander \([18,19,20]\), Torres \([46]\), Stein \([42]\).

The definition in (1.10) is also useful because it easily adapts to Littlewood–Paley analysis
of type 1,1-operators. Here the systematic point of departure is the well-known paradifferential
splitting based on dyadic coronas (cf. Section 5 for details), as used by e.g. Bony \([2]\),
Yamazaki \([48]\), Marschall \([33]\):

\[
a(x,D) = a_{\psi}^{(1)}(x,D) + a_{\psi}^{(2)}(x,D) + a_{\psi}^{(3)}(x,D).
\]

(1.19)

Since the 1980’s splittings like (1.19) have been used in microlocal analysis of (1.4) as well as in
numerous proofs of continuity of \( a(x,D) \) in Sobolev spaces \( H^s_p \) and Hölder–Zygmund spaces
\( C^s_x \), or the more general Besov and Lizorkin–Triebel scales \( B^s_{p,q} \) and \( F^s_{p,q} \). For type 1,1 operators
(1.19) was used by Bourdaud \([5,6,7]\), Marschall \([33]\), Runst \([41]\), and the author in \([26,29]\),
and in \([23,24]\) where the Lizorkin–Triebel spaces \( F^s_{p,1} \) were shown to be optimal substitutes for
the Sobolev spaces \( H^s_p \) at the borderline \( s = d \) for the domains of operators in \( \text{OP}(S^d_{1,1}) \).

It is known that the decomposition (1.19) follows from the bilinear way \( \psi \) enters (1.10), and
that one finds at once the three infinite series in (5.8)–(5.10) below, which define the \( a_{\psi}^{(j)}(x,D) \).
But it is a main point of Sections 5 and 6 to verify that each of these series gives an operator
\( a_{\psi}^{(j)}(x,D) \) also belonging to \( \text{OP}(S^d_{1,1}) \); which is non-trivial because of the modulation function \( \psi \)
in (1.10).
As general properties of the type 1, 1-operators $a^{(1)}_\psi(x, D)$ and $a^{(3)}_\psi(x, D)$, they are shown here to satisfy the twisted diagonal condition (1.14), so (1.19) can be seen as a main source of such operators. Consequently these terms are harmless as they are defined on $\mathcal{S}'$ because of (1.12) ff.

Therefore, it is the so-called symmetric term $a^{(2)}_\psi(x, D)$ which may cause $a(x, D)u$ to be undefined, as was previously known e.g. for functions $u$ in a Sobolev space; cf. [24]. This delicate situation is clarified in Theorem 6.4 with a natural identification of type 1, 1-domains, namely

$$D(a(x, D)) = D(a^{(2)}_\psi(x, D)).$$

(1.20)

This might seem obvious at first glance, but really is without meaning before the $a^{(2)}_\psi$-series has been shown to define a type 1, 1-operator. Hence (1.20) is a corollary to the cumbersome book-keeping needed for this identification of $a^{(2)}_\psi(x, D)$. In fact, the real meaning of (1.20) is that both domains consist of the $u \in \mathcal{S}'$ for which the $a^{(2)}_\psi$-series converges; cf. Theorem 6.4.

In comparison, convergence of the series for $a^{(1)}_\psi(x, D)u$ and $a^{(3)}_\psi(x, D)u$ in Theorem 6.1 verified explicitly for all $u \in \mathcal{S}'$, $a \in S^0_{1,1}$, and these operators are proved to be of type 1, 1. Thus (1.19) is an identity among type 1, 1-operators. It was exploited for estimates of arbitrary $a \in S^0_{1,1}$ in e.g. Sobolev spaces $H^s_p$ and Hölder–Zygmund spaces $C^s_\ast$ in [29], by giving full proofs (i.e. the first based on (1.10)) of the boundedness for all $s > 0$, $1 < p < \infty$,

$$a(x, D) : H^{r+d}_p(\mathbb{R}^n) \to H^r_p(\mathbb{R}^n), \quad a(x, D) : C^{s+d}_{\ast}(\mathbb{R}^n) \to C^s_{\ast}(\mathbb{R}^n).$$

(1.21)

This was generalised in [29] to all $s \in \mathbb{R}$ when $a$ fulfills the twisted diagonal condition of order $\sigma$ in (1.16) for all $\sigma \in \mathbb{R}$. This sufficient condition extends results for $p = 2$ of Hörmander [18, 19] to $1 < p < \infty$, $s \in \mathbb{R}$. The special case $s = 0 = d$ was considered recently in [21].

The present results on $a^{(j)}_\psi(x, D)$ are of course natural, but they do rely on two techniques introduced rather recently in works of the author. One ingredient is a pointwise estimate

$$|a(x, D)u(x)| \leq c u^s(x), \quad x \in \mathbb{R}^n,$$

(1.22)

cf. Section 3 and [28], in terms of the Peetre–Fefferman–Stein maximal function

$$u^s(x) = \sup_{x \in \mathbb{R}^n} \frac{|u(x - y)|}{(1 + |x - y|)^s}, \quad \text{when } \sup \overset{\wedge}{u} \subset \overline{B}(0, R).$$

(1.23)

Although $u \mapsto u^s$ is non-linear, it is useful for convergence of series: e.g. in $H^s_p$ since it is $L^p$-bounded, and as shown here also in $\mathcal{S}'$ since it has polynomial bounds $u^s(x) \leq c(1 + |x|)^N$.

The second ingredient is the Spectral Support Rule from [26]; cf. also [23, 24]. It provides control of $\sup \overset{\wedge}{\mathcal{F}}(a(x, D)u)$ in terms of the supports of $\overset{\wedge}{u}$ and $\mathcal{K}(\xi, \eta)$ in (1.15),

$$\sup \overset{\wedge}{\mathcal{F}}(a(x, D)u) \subset \sup \mathcal{K} \circ \sup \overset{\wedge}{\mathcal{F}}u = \{ \xi + \eta \mid (\xi, \eta) \in \sup \overset{\wedge}{u}, \eta \in \sup \overset{\wedge}{u} \}^\circ.$$  (1.24)

The simple case in which $u \in \mathcal{S}$ was covered by Metivier [34, Prop. 4.2.8]. A review of (1.24) is given in Appendix B, including an equally easy proof for arbitrary $\overset{\wedge}{u} \in \mathcal{S}'$ and $a \in S^d_{1,1}$.

A main purpose of (1.24) is to avoid a cumbersome approximation by elementary symbols. These were introduced by Coifman and Meyer [11] to reduce the task of bounding the support of $\mathcal{F}(a(x, D)u)$: indeed, elementary symbols have the form $a(x, \eta) = \sum m_j(x) \Phi_j(\eta)$ for multipliers...
In general, analysis of type 1 \cite[Ch. 9]{20}; with consequences derived from the present operator definition. Littlewood–Paley

2.1. \textquoteleft the paradifferential splitting (1.19) is further analysed, especially concerning their continuity on

the self-adjoint part of extended to a version for frequency modulated operators. Section 4 gives a precise analysis of

smooth functions and their duals \cite{23, 24, 26}, but it suffices to recall just a few facts.

\begin{equation}
\sum_j \Phi_j \ast (\hat{u} \mathcal{F}^\dagger) \text{ is a finite sum when } \hat{u} \in \mathcal{S}' \text{; whence the rule for convolutions yields (1.24) for such symbols.}
\end{equation}

However, approximation by elementary symbols is not just technically redundant because of (1.24), it would also be particularly cumbersome to use for a type 1,1-symbol, as (1.10) would then have to be replaced by a double-limit procedure. Moreover, in the proof of (1.19), as well as in the $L_p$-theory based on it in \cite{29}, (1.24) also yields a significant simplification.

\begin{remark}
The Spectral Support Rule (1.24) shows clearly that Hörmander’s twisted diagonal condition (1.14) ensures that $a(x, D)$ cannot change (large) frequencies in $\text{supp} \hat{u}$ to 0. In fact, the support condition in (1.14) implies that $\xi$ cannot be close to $-\eta$ when $(\xi, \eta) \in \text{supp} \hat{a}$, which by (1.24) means that $\eta \in \text{supp} \hat{u}$ will be changed by $a(x, D)$ to the frequency $\xi + \eta \neq 0$.
\end{remark}

\textbf{Contents.} Notation is settled in Section 2 along with basics on operators of type 1,1 and the $C^\infty$-results in (1.17) ff. In Section 3 some pointwise estimates are recalled from \cite{28} and then extended to a version for frequency modulated operators. Section 4 gives a precise analysis of the self-adjoint part of $S^1_{1,1}$, relying on the results and methods from Hörmander’s lecture notes \cite[Ch. 9]{20}; with consequences derived from the present operator definition. Littlewood–Paley analysis of type 1,1-operators is developed in Section 5. In Section 6 the operators resulting from the paradifferential splitting (1.19) is further analysed, especially concerning their continuity on $\mathcal{S}'(\mathbb{R}^n)$ and the domain relation (1.20). Section 7 contains a few final remarks.

\section{Preliminaries on Type 1,1-Operators}

Notation and notions from Schwartz’ distribution theory, such as the spaces $C^\infty_0$, $\mathcal{S}$, $C^\infty$ of smooth functions and their duals $\mathcal{S}'$, $\mathcal{S}'$, $\mathcal{E}'$ of distributions, and the Fourier transformation $\mathcal{F}$, will be as in Hörmander’s book \cite{17} with these exceptions: $\langle u, \varphi \rangle$ denotes the value of a distribution $u$ on a test function $\varphi$. The Sobolev space of order $s \in \mathbb{R}$ based on $L_p$ is written $H^s_p$, and $H^s = H^s_2$. The space $\mathcal{O}_M(\mathbb{R}^n)$ consists of the slowly increasing $f \in C^\infty(\mathbb{R}^n)$, i.e. the $f$ that for each multiindex $\alpha$ and some $N > 0$ fulfils $|D^\alpha f(x)| \leq c(1 + |x|)^N$.

As usual $t_+ = \max(0, t)$ is the positive part of $t \in \mathbb{R}$ whilst $[t]$ denotes the greatest integer $\leq t$. In general, $c$ will denote positive constants, specific to the place of occurrence.

\subsection{The General Definition of Type 1,1-Operators.}

For type 1,1-operators the reader may consult \cite{26} for an overview of previous results. The present paper is partly a continuation of \cite{23, 24, 26}, but it suffices to recall just a few facts.

By standard quantization, each operator $a(x, D)$ is defined on the Schwartz space $\mathcal{S}(\mathbb{R}^n)$ by

\begin{equation}
a(x, D)u = \text{OP}(a)u(x) = (2\pi)^{-n} \int e^{ix \cdot \eta} a(x, \eta) \mathcal{F}u(\eta) \, d\eta, \quad u \in \mathcal{S}(\mathbb{R}^n).
\end{equation}

Hereby its symbol $a(x, \eta)$ is required to be in $C^\infty(\mathbb{R}^n \times \mathbb{R}^n)$, of order $d \in \mathbb{R}$ and type 1,1, which means that for all multiindices $\alpha$, $\beta \in \mathbb{N}^n_0$ it fulfils (1.1), or more precisely has finite seminorms:

\begin{equation}
p_{\alpha, \beta}(a) := \sup_{x, \eta \in \mathbb{R}^n} (1 + |\eta|)^{-|d| - |\alpha| - |\beta|} |D^\alpha_\eta D^\beta_x a(x, \eta)| < \infty.
\end{equation}

The Fréchet space of such symbols is denoted by \( S_{1,1}^d(\mathbb{R}^n \times \mathbb{R}^n) \), or just \( S_{1,1}^d \) for brevity, while as usual \( S^{-\infty} = \cap_d S_{1,1}^d \). Basic estimates yield that the bilinear map \( (a,u) \mapsto a(x,D)u \) is continuous
\[
S_{1,1}^d \times \mathcal{S} \to \mathcal{S}.
\] (2.3)

The distribution kernel \( K(x,y) = \mathcal{F}^{-1}_{\eta \to x}a(x,\eta)|_{\eta = x-y} \) is well known to be \( C^\infty \) for \( x \neq y \) also in the type 1, 1 context; cf. [26, Lem. 4.3]. It fulfills \( \langle a(x,D)u, \varphi \rangle = \langle K, \varphi \otimes u \rangle \) for all \( u, \varphi \in \mathcal{S} \).

For arbitrary \( u \in \mathcal{S'} \setminus \mathcal{S} \) it is a delicate question whether or not \( a(x,D)u \) is defined. The general definition of type 1, 1-operators in [26] uses a symbol modification, exploited throughout below, namely \( b(x,\eta) = \psi(2^{-m}D_x)a(x,\eta) \), or more precisely
\[
\hat{b}(\xi,\eta) := \mathcal{F}_{x \to \xi}b(x,\eta) = \psi(2^{-m}\xi)\hat{a}(\xi,\eta).
\] (2.4)

**Definition 2.1.** If a symbol \( a(x,\eta) \) is in \( S_{1,1}^d(\mathbb{R}^n \times \mathbb{R}^n) \) and \( u \in \mathcal{S'}(\mathbb{R}^n) \) whilst \( \psi \in C_0^\infty(\mathbb{R}^n) \) is an arbitrary cut-off function equal to 1 in a neighbourhood of the origin, let
\[
a_{\psi}(x,D)u := \lim_{m \to \infty} \text{OP}(\psi(2^{-m}D_x)a(x,\eta)\psi(2^{-m}\eta))u.
\] (2.5)

When for each such \( \psi \) the limit \( a_{\psi}(x,D)u \) exists in \( \mathcal{D}'(\mathbb{R}^n) \) and moreover is independent of \( \psi \), then \( u \) belongs to the domain \( D(a(x,D)) \) by definition and
\[
a(x,D)u = a_{\psi}(x,D)u.
\] (2.6)

This way \( a(x,D) \) is a linear map \( \mathcal{S}'(\mathbb{R}^n) \to \mathcal{D}'(\mathbb{R}^n) \) with dense domain, as by (2.3) it contains \( \mathcal{D}(\mathbb{R}^n) \). (Use of \( D(\cdot) \) for the domain should not be confounded with derivatives, such as \( D^\alpha \) or \( a(x,D) \).)

This was called definition by *vanishing frequency modulation* in [26], because the removal of high frequencies in \( x \) and \( \eta \) achieved by \( \psi(2^{-m}D_x) \) and \( \psi(2^{-m}\eta) \) disappears for \( m \to \infty \). Note that the action on \( u \) is well defined for each \( m \) in (2.5) as the modified symbol is in \( S^{-\infty} \).

Occasionally the function \( \psi \) will be referred to as a *modulation function*.

The frequency modulated operator \( \text{OP}(\psi(2^{-m}D_x)a(x,\eta)\psi(2^{-m}\eta)) \) has, by the comparison made in [26, Prop. 5.11], its kernel \( K_m(x,y) \) conveniently given as a convolution, up to conjugation by the involution \( M: (x,y) \mapsto (x,x-y) \),
\[
K_m(x,y) = 4^{mn}(\mathcal{F}^{-1}_{x \to \xi}(\varphi(2^{m}\xi)) \otimes \mathcal{F}^{-1}_{y \to \zeta}(\psi(2^{m}\zeta))) \ast (K \circ M)(x,x-y).
\] (2.7)

**Remark 2.2.** It is used below that when \( \varphi, \chi \in C_0^\infty(\mathbb{R}^n) \) are such that \( \chi \equiv 1 \) on a neighbourhood of \( \text{supp} \varphi \), then since \( \text{supp} \varphi \otimes (1-\chi) \) is disjoint from the diagonal and bounded in the \( x \)-direction, there is convergence in the topology of \( \mathcal{D}'(\mathbb{R}^n \times \mathbb{R}^n) \):
\[
\varphi(x)(1-\chi(y))K_m(x,y) \quad \text{converges as } m \to \infty \quad \text{to} \quad \varphi(x)(1-\chi(y))K(x,y).
\] (2.8)

However, this requires verification because the commutator of the convolution (2.7) and pointwise multiplication by \( \varphi \otimes (1-\chi) \) is a nontrivial pseudo-differential, hence non-local operator. A proof of (2.8) based on the Regular Convergence Lemma can be found in [26, Prop. 6.3].

In general the calculus of type 1, 1-operators is delicate, cf. [18, 19, 20], but the following result from [29] is just an exercise (cf. the proof there). It is restated here for convenience.
Proposition 2.3. When \( a(x, \eta) \) is in \( S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n) \) and a symbol with constant coefficients \( b(\eta) \) belongs to \( S^{d_2}_{1,0}(\mathbb{R}^n \times \mathbb{R}^n) \), then \( c(x, \eta) := a(x, \eta)b(\eta) \) is in \( S^{d_1+d_2}_{1,1}(\mathbb{R}^n \times \mathbb{R}^n) \) and

\[
c(x, D)u = a(x, D)b(D)u. \tag{2.9}
\]

In particular \( D(c(x, D)) = D(a(x, D)b(D)) \), so the two sides are simultaneously defined.

This result applies especially to differential operators, say \( b(D) = D_f \) for simplicity. But as a minor novelty, the classical commutator identity needs an atypical substitute:

Proposition 2.4. For \( a \in S^d_{1,1} \) the commutator

\[
[D_j, a(x, D)] = D_ja(x, D) - a(x, D)D_j \tag{2.10}
\]

equals \( \text{OP}(D_x, a(x, \eta)) \) on the Schwartz space \( \mathcal{S}(\mathbb{R}^n) \), whilst in \( \mathcal{S}'(\mathbb{R}^n) \) there is an inclusion in the operator theoretical sense,

\[
a(x, D)D_j + [D_j, a(x, D)] \subset D_ja(x, D). \tag{2.11}
\]

The commutator symbol \( D_x, a(x, \eta) \) is in \( S^{d+|\beta|}_{1,1} \).

Proof. By classical calculations, any modulation function \( \psi \) gives the following formula for \( u \in \mathcal{S}' \), hence for all \( u \in \mathcal{S}' \) as the symbols are in \( S^{-\infty} \),

\[
\text{OP}(\psi(2^{-m}D_x)a(x, \eta)\psi(2^{-m}\eta)\eta_j)u + \text{OP}(\psi(2^{-m}D_x)D_xa(x, \eta)\psi(2^{-m}\eta))u \\
= D_j\text{OP}(\psi(2^{-m}D_x)a(x, \eta)\psi(2^{-m}\eta))u. \tag{2.12}
\]

When both terms on the left have \( \psi \)-independent limits for \( m \to \infty \), so has the right-hand side. As the first term then is \( a(x, D)D_j u \), cf. Proposition 2.3, this entails that the common domain \( D(a(x, D)D_j) \cap D([D_j, a(x, D)]) \) is contained in that of \( D_ja(x, D) \), with the same actions. \( \square \)

The inclusion (2.11) is strict in some cases, for the domains are not always invariant under differentiation. This is a well-known consequence of the classical counterexamples, which are recalled below for the reader’s convenience:

Example 2.5. The classical example of a symbol of type 1,1 results from an auxiliary function \( A \in C^\infty_0(\mathbb{R}^n) \), say with \( \text{supp} A \subset \{ \eta \mid \frac{1}{4} \leq |\eta| \leq \frac{3}{4} \} \), and a fixed vector \( \theta \in \mathbb{R}^n \),

\[
a_\theta(x, \eta) = \sum_{j=0}^{\infty} 2^{jd} e^{-i2^{j}(x \cdot \theta)} A(2^{-j}\eta). \tag{2.13}
\]

Here \( a_\theta \in C^\infty(\mathbb{R}^n \times \mathbb{R}^n) \), since the terms are disjointly supported, and it clearly belongs to \( S^d_{1,1} \).

These symbols were used both by Ching [10] and Bourdaud [7] to show \( L_2 \)-unboundedness for \( d = 0, |\theta| = 1 \). Refining this, Hörmander [18] established that continuity \( \mathcal{H}^s \to \mathcal{D}' \) with \( s > -r \) holds if and only if \( \theta \) is a zero of \( A \) of order \( r \in \mathbb{N}_0 \). [26] gave an extension to \( d \in \mathbb{R}, \theta \neq 0 \).

The non-preservation of wavefront sets discovered by Parenti and Rodino [39] also relied on \( a_\theta(x, \eta) \). Their ideas were extended to all \( n \geq 1 \), \( d \in \mathbb{R} \) in [26, Sect. 3.2] and refined by applying \( a_{2\theta}(x, D) \) to a product \( \nu(x)f(x \cdot \theta) \), where \( \nu \in \mathcal{F}^{-1}C^\infty_0 \) is an analytic function controlling
the spectrum, whilst the highly oscillating \( f \) is Weierstrass’ nowhere differentiable function for orders \( d \in [0, 1] \), in a complex version with its wavefront set along a half-line. (Nowhere differentiability was shown with a small microlocalisation argument, explored in [27].)

Moreover, \( a_\theta(x,D) \) is unclosable in \( \mathcal{S}' \) when \( A \) is supported in a small ball around \( \theta \), as in [26, Lem. 3.2]. Hence Definition 2.1 cannot in general be simplified to a closure of the graph in \( \mathcal{S}' \times \mathcal{D}' \).

As a basic result, it was shown in [26, Sec. 4] that the \( C^\infty \)-subspace \( \mathcal{S}(\mathbb{R}^n) + \mathcal{F}^{-1}\mathcal{E}'(\mathbb{R}^n) \) always is contained in the domain of \( a(x,D) \) and that

\[
a(x,D): \mathcal{S}(\mathbb{R}^n) + \mathcal{F}^{-1}\mathcal{E}'(\mathbb{R}^n) \to \mathcal{O}_M(\mathbb{R}^n).
\]

(2.14)

In fact, if \( u = v + v' \) is any splitting with \( v \in \mathcal{S} \) and \( v' \in \mathcal{F}^{-1}\mathcal{E}' \), then

\[
a(x,D)u = a(x,D)v + \text{OP}(a(1 \otimes \chi))v',
\]

(2.15)

whereby \( a(1 \otimes \chi)(x,\eta) = a(x,\eta)\chi(\eta) \) and \( \chi \in C^\infty_{0}(\mathbb{R}^n) \) is arbitrarily chosen so that \( \chi = 1 \) holds in a neighbourhood of \( \text{supp } \mathcal{F}v' \subseteq \mathbb{R}^n \). Here \( a(x,\eta)\chi(\eta) \) is in \( \mathcal{S}^{-\infty} \) so that \( \text{OP}(a(1 \otimes \chi)) \) is defined on \( \mathcal{S}' \). Hence \( a(x,D)(\mathcal{F}^{-1}\mathcal{E}') \subseteq \mathcal{O}_M(\mathbb{R}^n) \).

It is a virtue of (2.14) that \( a(x,D) \) is compatible with for example \( \text{OP}(\mathcal{S}^\infty_{1,0}) \); cf. [26] for other compatibility questions. Therefore some well-known facts extend to type 1,1-operators:

**Example 2.6.** Each \( a(x,D) \) of type 1,1 is defined on all polynomials and

\[
a(x,D)(x^\alpha) = D^\alpha_{\eta}(e^{ix\cdot\eta}a(x,\eta)) \mid_{\eta = 0}.
\]

(2.16)

In fact, \( f(x) = x^\alpha \) has \( \hat{f}(\eta) = (2\pi)^n(-D\eta)^\alpha \delta_0(\eta) \) with \( \text{supp } \{ 0 \} \), so it is seen for \( v = 0 \) in (2.15) that \( a(x,D)f(x) = (\hat{f}, (2\pi)^{-n}e^{ix\cdot\eta}a(x,\cdot)\chi(\cdot)) \text{ where } \chi = 1 \text{ around } 0; \) thence (2.16).

**Example 2.7.** Also when \( A = a(x,D) \) of type 1,1, one can recover its symbol from the formula

\[
a(x,\xi) = e^{-ix\cdot\xi}A(e^{ix\cdot\xi}).
\]

(2.17)

Here \( \mathcal{F}e^{ix\cdot\xi} = (2\pi)^n\delta_\xi(\eta) \) has compact support, so again it follows from (2.14) that (via a suppressed cut-off) one has \( A(e^{ix\cdot\xi}) = \langle \delta_\xi, e^{ix\cdot\eta}a(x,\cdot) \rangle = e^{ix\cdot\xi}a(x,\xi) \).

2.2. General Smooth Functions. To go beyond the smooth functions in (2.14), it is shown in this subsection how one can extend a remark by Bourdaud [8] on singular integral operators, which shows that every type 1,1 symbol \( a(x,\eta) \) of order \( d = 0 \) induces a map \( \hat{A}: \mathcal{O}_M \to \mathcal{D}' \).

Indeed, Bourdaud defined \( \hat{A}f \) for \( f \in \mathcal{O}_M(\mathbb{R}^n) \) as the distribution that on \( \phi \in C^\infty_{0}(\mathbb{R}^n) \) is given by the following, using the distribution kernel \( K \) of \( a(x,D) \) and an auxiliary function \( \chi \in C^\infty_{0}(\mathbb{R}^n) \) equal to 1 on a neighbourhood of \( \text{supp } \phi \),

\[
\langle \hat{A}f, \phi \rangle = \langle a(x,D)(\chi f), \phi \rangle + \iint K(x,y)(1 - \chi(y))f(y)\phi(x)dydx.
\]

(2.18)

However, to free the discussion from the slow growth in \( \mathcal{O}_M \), one may restate this in terms of the tensor product \( 1 \otimes f \) in \( \mathcal{S}'(\mathbb{R}^n \times \mathbb{R}^n) \) acting on \( (\phi \otimes (1 - \chi))K \in \mathcal{S}(\mathbb{R}^n \times \mathbb{R}^n) \), i.e.

\[
\langle \hat{A}f, \phi \rangle = \langle a(x,D)(\chi f), \phi \rangle + \langle 1 \otimes f, (\phi \otimes (1 - \chi))K \rangle.
\]

(2.19)
One advantage here is that both terms obviously make sense as long as \( f \) is smooth and temperate, i.e. for every \( f \in C^\infty(\mathbb{R}^n) \cap \mathcal{S}'(\mathbb{R}^n) \).

Moreover, for \( \varphi \) with support in the interior \( \mathcal{C}^\circ \) of a compact set \( \mathcal{C} \subset \mathbb{R}^n \) and \( \chi = 1 \) on a neighbourhood of \( \mathcal{C} \), the right-hand side of (2.19) gives the same value for any \( \tilde{\chi} \in C^\infty_0 \) equal to 1 around \( \mathcal{C} \), for in the difference of the right-hand sides equals 0 since \( \langle a(x,D)(\chi - \tilde{\chi})f, \varphi \rangle \) is seen from the kernel relation to equal \(-\langle 1 \otimes f, (\varphi(\tilde{\chi} - \chi))K \rangle \).

Crude estimates of (2.19) now show that \( \tilde{\varphi} \) follows, for the first term is in \( C^\infty \), the recollement de morceaux theorem yields that a distribution \( \tilde{\varphi} \in \mathcal{D}'(\mathbb{R}^n) \) is defined by (2.19).

There is also a more explicit formula for \( \tilde{\varphi} \): when \( \tilde{\varphi} \in C^\infty_0 \) is chosen so that \( \tilde{\varphi} \equiv 1 \) around \( \mathcal{C} \) while \( \text{supp} \tilde{\varphi} \) has a neighbourhood where \( \chi = 1 \), then \( \varphi = \tilde{\varphi} \varphi \) in (2.19) gives, for \( x \in \mathcal{C}^\prime \),

\[
\tilde{\varphi}(x) = a(x,D)(\chi f)(x) + \langle f, (\tilde{\varphi}(x)(1 - \chi(\cdot)))K(x,\cdot) \rangle.
\] (2.20)

Now \( \tilde{\varphi} \in C^\infty \) follows, for the first term is in \( \mathcal{S} \), and the second coincides in \( \mathcal{C}^\circ \) with a function in \( \mathcal{S} \), as a corollary to the construction of \( g \otimes f \in \mathcal{S}'(\mathbb{R}^n \times \mathbb{R}^n) \) for \( f, g \in \mathcal{S}'(\mathbb{R}^n) \).

Post festum, it is seen in (2.19) that when \( f \to 0 \) in both \( C^\infty \) and in \( \mathcal{S}' \), then \( \chi f \to 0 \) in \( \mathcal{S} \) while \( 1 \otimes f \to 0 \) in \( \mathcal{S}' \). Therefore \( \tilde{\varphi} \to 0 \) in \( \mathcal{S}' \), which is a basic continuity property of \( \tilde{\varphi} \).

By setting \( \tilde{\varphi} \) in relation to Definition 2.1, the above gives the new result that \( a(x,D) \) always is a map defined on the maximal set of smooth functions, i.e. on \( C^\infty \cap \mathcal{S}' \):

**Theorem 2.8.** Every \( a(x,D) \in \text{OP}(S^1_{1,1}(\mathbb{R}^n \times \mathbb{R}^n)) \) restricts to a map

\[
a(x,D) : C^\infty(\mathbb{R}^n) \cap \mathcal{S}'(\mathbb{R}^n) \to C^\infty(\mathbb{R}^n),
\] (2.21)

which locally is given by formula (2.20). The map (2.21) is continuous when \( C^\infty(\mathbb{R}^n) \) has the usual Fréchet space structure and \( \mathcal{S}'(\mathbb{R}^n) \) has the strong dual topology.

The intersection \( C^\infty \cap \mathcal{S}' \) is topologised by enlarging the set of seminorms on \( C^\infty \) by those on \( \mathcal{S}' \). Here the latter have the form \( f \mapsto \sup_{\psi \in \mathcal{B}} |\langle f, \psi \rangle| \) for an arbitrary bounded set \( \mathcal{B} \subset \mathcal{S} \).

**Proof.** Let \( A_m = \text{OP}(\psi(2^{-m}D_x)a(x,\eta)\psi(2^{-m}\eta)) \) so that \( a(x,D)u = \lim_m A_m u \) when \( u \) belongs to \( D(a(x,D)) \). With \( f \in C^\infty \cap \mathcal{S}' \) and \( \varphi, \chi \) as above, this is the case for \( u = \chi f \in C^\infty_0 \).

Exploiting the convergence in Remark 2.2 in (2.19), it is seen that

\[
\langle \tilde{\varphi}, \varphi \rangle = \lim_m \langle A_m(\chi f), \varphi \rangle + \lim_m \iint K_m(x,y)(1 - \chi(y))f(x)\varphi(x) \, dy \, dx.
\] (2.22)

Here the integral equals \( \langle A_m(f - \chi f), \varphi \rangle \) by the kernel relation, for \( A_m \in \text{OP}(S^{-\infty}) \) and \( f \) may as an element of \( \mathcal{S}' \) be approached from \( C^\infty_0 \). So (2.22) yields

\[
\langle \tilde{\varphi}, \varphi \rangle = \lim_m \langle A_m(\chi f), \varphi \rangle + \lim_m \langle A_m(f - \chi f), \varphi \rangle = \lim_m \langle A_m f, \varphi \rangle.
\] (2.23)

Thus \( A_m f \to \tilde{\varphi} \), which by (2.19) is independent of \( \psi \). Hence \( \tilde{\varphi} \subset a(x,D) \) as desired.

That \( a(x,D)(C^\infty \cap \mathcal{S}') \) is contained in \( C^\infty \) now follows from the remarks to (2.20).

When \( f \to 0 \) in \( C^\infty \), then clearly \( D^\infty a(x,D)(f \chi) \to 0 \) in \( \mathcal{S} \), hence uniformly on \( \mathbb{R}^n \). It is also straightforward to see that \( (\tilde{\varphi}(x)(1 - \chi(\cdot)))K(x,\cdot) \) stays in a bounded set in \( \mathcal{S}'(\mathbb{R}^n) \) as \( x \) runs through \( \mathcal{C} \). Therefore, when \( f \to 0 \) also in the strong dual topology on \( \mathcal{S}' \), then the second term
in (2.20) tends to 0 uniformly with respect to \( x \in \mathcal{C} \). As \( x \)-derivatives may fall on \( K(x, \cdot) \), the same argument gives that \( \sup_{\mathcal{C}} |D^{\alpha} \tilde{A}f| \to 0 \). Hence \( f \to 0 \) in \( C^\infty(\mathcal{S}', \mathcal{S}) \) implies \( D^{\alpha} \tilde{A}f \to 0 \) in \( C^\infty(\mathbb{R}^n) \), which gives the stated continuity property. \( \square \)

**Remark 2.9.** It is hardly a drawback that continuity in Theorem 2.8 holds for the strong dual topology on \( \mathcal{S}' \), as for sequences weak and strong convergence are equivalent (a well-known consequence of the fact that \( \mathcal{S} \) is a Montel space).

In view of Theorem 2.8, the difficulties for type 1,1-operators do not stem from growth at infinity for \( C^\infty \)-functions. Obviously the codomain \( C^\infty \) is not contained in \( \mathcal{S}' \), but this is not just made possible by the use of \( \mathcal{S}' \) in Definition 2.1, it is indeed decisive for the above construction.

In the proof above, the fact that \( \tilde{A}f \in C^\infty \) also follows from the pseudo-local property of \( a(x, D) \); cf [26, Thm. 6.4]. The direct argument above is rather short, though. In addition to the smoothness, the properties of \( a(x, D)f \) can be further sharpened by slow growth of \( f \):

**Corollary 2.10.** Every type 1,1-operator \( a(x, D) \) leaves \( \mathcal{O}_M(\mathbb{R}^n) \) invariant.

**Proof.** If \( f \in \mathcal{O}_M \), then it follows that \((1 + |x|)^{-2N}D^{\alpha} \tilde{A}f(x)\) is bounded for sufficiently large \( N \), since in the second contribution to (2.20) clearly \((1 + |y|)^{-2N}f(y)\) is in \( L_1 \) for large \( N \): the resulting factor \((1 + |y|)^{-2N}\) may be absorbed by \( K \), using that \( r = \text{dist}(\text{supp} \Phi, \text{supp}(1 - \chi)) > 0 \), since for \( x \in \text{supp} \Phi, y \in \text{supp}(1 - \chi) \),

\[
(1 + |y|)^{-2N}|D^{\alpha}K(x, y)| \leq (1 + |x|)^{-2N} \max(1, 1/r)^{-2N}(r + |x - y|)^{-2N}|D^{\alpha}K(x, y)| \leq c(1 + |x|)^{-2N} \sup_{x \in \mathbb{R}^n} \int |(4 \Delta \bar{\eta})^N((\eta + D_x)\alpha a(x, \eta))| d\eta, \tag{2.24}
\]

where the supremum is finite for \( 2N > d + |\alpha| + n \) (by induction \((\eta + D_x)\alpha : \mathcal{S}'_{1,1} \to \mathcal{S}'_{d + |\alpha|}\)). Moreover, \( c = \max(1, 1/r)^{-2N}/(2\pi)^n \) can be chosen uniformly for \( x \in \mathbb{R}^n \) as it suffices to have (2.20) with \( 0 \leq \chi \leq 1 \) and \( \mathcal{C} = \mathcal{B}(0, j) \), \( \text{supp} \Phi = \overline{B}(0, j + 1) \) and \( \chi^{-1}(\{1\}) = \overline{B}(0, j + 2) \) for an arbitrarily large \( j \in \mathbb{N} \), which yields \( r = 1, c \leq 1 \). Thus \((1 + |x|)^{-2N}|D^{\alpha} \tilde{A}f(x)| \) is less than \( s_{\alpha, N} \int_{\mathbb{R}^n} (1 + |y|)^{-2N}|f(y)| dy \) for all \( x \in \mathbb{R}^n \), \( s_{\alpha, N} \) as the sup in (2.24). Hence \( \tilde{A}f \in \mathcal{O}_M \). \( \square \)

**Example 2.11.** The space \( C^\infty(\mathbb{R}^n) \cap \mathcal{S}'(\mathbb{R}^n) \) clearly contains functions of non-slow growth, e.g.

\[
f(x) = e^{x_1 + \cdots + x_n} \cos(e^{x_1 + \cdots + x_n}). \tag{2.25}
\]

Recall that \( f \in \mathcal{S}' \) because \( f = iD_1 g \) for \( g(x) = \sin(e^{x_1 + \cdots + x_n}) \), which is in \( L_\infty \subset \mathcal{S}' \). But \( g \notin \mathcal{O}_M \), so already for \( a(x, D) = iD_1 \) the space \( \mathcal{O}_M \) cannot contain the range in Theorem 2.8.

**Remark 2.12.** Prior to the T1-theorem, David and Journé explained in [12] how a few properties of the distribution kernel of a continuous map \( T : C^\infty_0(\mathbb{R}^n) \to \mathcal{S}'(\mathbb{R}^n) \) makes \( T(1) \) well defined modulo constants; in particular if \( T \in \text{Op}(\mathcal{S}'_{1,1}) \). Bourdaud [8] used this construction for \( A \), so by Theorem 2.8 this extension of \( T \in \text{Op}(\mathcal{S}'_{1,1}) \) from [12] is contained in Definition 2.1.

### 2.3. Conditions along the Twisted Diagonal.

As the first explicit condition formulated for the symbol of a type 1,1-operator, Hörmander [18] proved that \( a(x, D) \) has an extension by continuity

\[
H^{s + d} \to H^{\infty} \quad \text{for every } s \in \mathbb{R} \tag{2.26}
\]
whenever \( a \in S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n) \) fulfils the twisted diagonal condition: for some \( B \geq 1 \)
\[
\hat{a}(\xi, \eta) = 0 \quad \text{where} \quad B(1 + |\xi + \eta|) < |\eta|.
\] (2.27)

In detail this means that the partially Fourier transformed symbol \( \hat{a}(\xi, \eta) := \mathcal{F}_{\xi \rightarrow \hat{\xi}}a(x, \eta) \) is trivial in a conical neighbourhood of a non-compact part of the twisted diagonal
\[
\mathcal{F} = \{ (\hat{\xi}, \eta) \in \mathbb{R}^n \times \mathbb{R}^n \mid \hat{\xi} + \eta = 0 \}.
\] (2.28)

Localisations to certain conical neighbourhoods of \( \mathcal{F} \) were also used by Hörmander [18, 19, 20] as
\[
\hat{a}_{\chi, \varepsilon}(\xi, \eta) = \hat{a}(\xi, \eta)\chi(\xi + \varepsilon\eta),
\] (2.29)

whereby the cut-off function \( \chi \in C^\infty(\mathbb{R}^n \times \mathbb{R}^n) \) is chosen to satisfy
\[
\chi(t\xi, t\eta) = \chi(\xi, \eta) \quad \text{for} \quad t \geq 1, |\eta| \geq 2
\] (2.30)
\[
\supp \chi \subset \{ (\xi, \eta) \mid 1 \leq |\eta|, |\xi| \leq |\eta| \}
\] (2.31)
\[
\chi = 1 \quad \text{in} \quad \{ (\xi, \eta) \mid 2 \leq |\eta|, 2|\xi| \leq |\eta| \}.
\] (2.32)

Using this, Hörmander introduced and analysed a milder condition than the strict vanishing in (2.27). Namely, for some \( \sigma \in \mathbb{R} \) the symbol should satisfy an estimate, for all multiindices \( \alpha \) and \( 0 < \varepsilon < 1 \),
\[
N_{\chi, \varepsilon, \alpha}(a) := \sup_{R > 0, \eta \in \mathbb{R}^n} R^{-d} \left( \int_{R \leq |\eta| \leq 2R} |R^{[\alpha]}D_{\eta}^\alpha a_{\chi, \varepsilon}(x, \eta)|^2 \frac{d\eta}{R^n} \right)^{1/2} \leq c_{\alpha, \sigma} \varepsilon^{\sigma + n/2 - |\alpha|}.
\] (2.33)

This is an asymptotic formula for \( \varepsilon \to 0 \). It always holds for \( \sigma = 0 \), cf. [20, Lem. 9.3.2]:

**Lemma 2.13.** When \( a \in S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n) \) and \( 0 < \varepsilon \leq 1 \), then \( a_{\chi, \varepsilon} \in C^\infty \) and
\[
|D_{\eta}^\beta D_{\xi}^\alpha a_{\chi, \varepsilon}(x, \eta)| \leq C_{\alpha, \beta}(a) \varepsilon^{-|\alpha|}(1 + |\eta|)^{|d| - |\alpha| + |\beta|}
\] (2.34)
\[
\left( \int_{R \leq |\eta| \leq 2R} |D_{\eta}^\alpha a_{\chi, \varepsilon}(x, \eta)|^2 d\eta \right)^{1/2} \leq C_{\alpha} R^{d}(\varepsilon R)^{n/2 - |\alpha|}.
\] (2.35)

The map \( a \mapsto a_{\chi, \varepsilon} \) is continuous in \( S^d_{1,1} \).

The last remark on continuity has been inserted here for later reference. It is easily verified by observing in the proof of [20, Lem. 9.3.2] (to which the reader is referred) that the constant \( C_{\alpha, \beta}(a) \) is a continuous seminorm in \( S^d_{1,1} \).

In case \( \sigma > 0 \) there is a faster convergence to 0 in (2.33). In [19] this was proved to imply that \( a(x, D) \) is bounded as a densely defined map
\[
H^{s+d}((\mathbb{R}^n) \to H^s(\mathbb{R}^n) \quad \text{for} \quad s > -\sigma.
\] (2.36)

The reader may consult [20, Thm. 9.3.5] for this (whilst [20, Thm. 9.3.7] gives four pages of proof of necessity of \( s \geq -\sup \sigma \), with supremum over all \( \sigma \) satisfying (2.33)).

Consequently, if \( \hat{a}(\xi, \eta) \) is so small along \( \mathcal{F} \) that (2.33) holds for all \( \sigma \in \mathbb{R} \), there is boundedness \( H^{s+d} \to H^s \) for all \( s \in \mathbb{R} \). E.g. this is the case when (2.27) holds, for since
\[
\supp \hat{a}_{\chi, \varepsilon} \subset \{ (\xi, \eta) \mid 1 + |\xi + \eta| \leq 2\varepsilon|\eta| \},
\] (2.37)
clearly $a_{\varepsilon} \equiv 0$ for $2\varepsilon < 1/B$ then.

**Example 2.14.** For the present paper it is useful to exploit Ching’s symbol (2.13) to show the existence of symbols fulfilling (2.33) for a given $\sigma$, at least for $\sigma \in \mathbb{N}$. To do so one may fix $|\theta| = 1$ and take some $A(\eta)$ in $C_0^\infty(\{ \eta \mid \frac{3}{4} < |\eta| < \frac{5}{4} \})$ with a zero of order $\sigma$ at $\theta$, so that Taylor’s formula gives $|A(\eta)| \leq c|\eta - \theta|^\sigma$ in a neighbourhood of $\theta$:

Indeed, as $\hat{a}(x, \eta) = (2\pi)^n \sum_{j=0}^\infty 2^j (\xi + 2^j \theta) A(2^{-j} \eta)$, clearly

$$a_{\varepsilon, \chi, \varepsilon}(x, \eta) = \sum_{j=0}^\infty 2^j e^{-ix \cdot 2^j \eta} \varepsilon^j \chi(2^j \eta, \varepsilon \eta) A(2^{-j} \eta).$$

Because $[R, 2R]$ is contained in $[\frac{3}{4} 2^{j-1}, \frac{3}{4} 2^{j-1}] \cup [\frac{3}{4} 2^{j}, \frac{3}{4} 2^{j}]$ for some $j \in \mathbb{Z}$, it suffices to estimate the integral in (2.33) only for $R = 3 \cdot 2^{j-2}$ with $j \geq 1$. Then it involves only the $j$th term, i.e.

$$\int_{R \leq |\eta| \leq 2R} |a_{\varepsilon, \chi, \varepsilon}(x, \eta)|^2 \, d\eta = \int_{R \leq |\eta| \leq 2R} R^{d/2} |A(\eta/R)|^2 |\chi(\eta/R, \varepsilon \eta)|^2 \, d\eta.$$  

By the choice of $\chi$, the integrand is 0 unless $|\eta - R \theta| \leq \varepsilon |\eta| \leq 2\varepsilon R$ and $1 \leq \varepsilon R$, so for small $\varepsilon$,

$$\int_{R \leq |\eta| \leq 2R} |a_{\varepsilon, \chi, \varepsilon}(x, \eta)|^2 \, d\eta \leq \|\chi\|_2^2 R^{d+2} \int_{|\xi - \theta| \leq 2\varepsilon} c|\xi - \theta|^{2\sigma} \, d\xi \leq c \varepsilon^{2\sigma + n} R^{n+2d}.$$  

Applying $(RD\eta)^\alpha$ before integration, $(RD\eta)^\gamma$ may fall on $A(\eta/R)$, which lowers the degree and yields (at most) $\varepsilon^{n/2 + \sigma - |\gamma|}$. In the factor $(RD\eta)^{\alpha - |\gamma|} \chi(\eta - R \theta, \varepsilon \eta)$ the homogeneity of degree $-|\alpha - \gamma|$ applies for $\varepsilon R \geq 2$ and yields a bound in terms of finite suprema over $B(\theta, 2) \times B(0, 2)$, hence is $O(1)$; else $\varepsilon R < 2$ so the factor is $O(R^{\alpha - |\gamma|}) = O(\varepsilon^{n - |\alpha|})$ when non-zero, as both entries are in norm less than 4 then. Altogether this verifies (2.33).

A lower bound of (2.39) by $e \varepsilon^{2\sigma + n} R^{n+2d}$ is similar (cf. [20, Ex. 9.3.3] for $\sigma = 0$) when $|A(\eta)| \geq c_0 |\eta - \theta|^\sigma$, which is obtained by taking $A$ as a localisation of $|\eta - \theta|^\sigma$ for even $\sigma$ (so $A \in C_0^\infty$). This implies that (2.33) does not hold for larger values of $\sigma$ for this $a_{\varepsilon}(x, \eta)$.

3. **Pointwise Estimates**

A crucial technique in this paper will be to estimate $|a(x, D)u(x)|$ at an arbitrary point $x \in \mathbb{R}^n$. Some of recent results on this by the author [28] are recalled here and further elaborated in Section 3.2 with an estimate of frequency modulated operators.

3.1. **The Factorisation Inequality.** First of all, by [28, Thm. 4.1], when $\hat{a} \subset \overline{B}(0, R)$, the action on $u$ by the operator $a(x, D)$ can be separated from $u$ at the cost of an estimate, which is the factorisation inequality

$$|a(x, D)u(x)| \leq F_0(N, R; x) u^*(N, R; x).$$  

(3.1)

Hereby $u^*(x) = u^*(N, R; x)$ denotes the maximal function of Peetre–Fefferman–Stein type,

$$u^*(N, R; x) = \sup_{y \in \mathbb{R}^n} \frac{|u(x - y)|}{(1 + R|y|)^N} = \sup_{y \in \mathbb{R}^n} \frac{|u(y)|}{(1 + R|x - y|)^N}.  $$

(3.2)

The parameter $N$ is often chosen to satisfy $N \geq \text{order } \hat{a}$. 
The $a$-factor $F_a$, also called the symbol factor, only depends on $u$ in a vague way, i.e. only through the $N$ and $R$ in (3.2). It is related to the distribution kernel of $a(x,D)$ as

$$F_a(N,R;x) = \int_{\mathbb{R}^n} (1 + R y)^N |\mathcal{F}^{-1}_{\eta \to y}(a(x,\eta)\chi(\eta))| dy, \quad (3.3)$$

where $\chi \in C_0^\infty(\mathbb{R}^n)$ should equal 1 in a neighbourhood of supp $\hat{u}$, or of $\bigcup_x \text{supp} a(x,\cdot) \hat{u}(\cdot)$.

In (3.1) both factors are easily controlled. For one thing the non-linear map $u \mapsto u^*$ has long been known to have bounds with respect to the $L_p$-norm; cf. [28, Thm. 2.6] for an elementary proof. But in the present paper it is more important that $u^*(x)$ is polynomially bounded thus: $|u(y)| \leq c(1 + |y|)^N \leq c(1 + R|y - x|)^N (1 + |x|)^N$ holds according to the Paley–Wiener–Schwartz Theorem if $N \geq \text{order } \hat{u}$ and $R \geq 1$, which by (3.2) implies

$$u^*(N,R;x) \leq c(1 + |x|)^N, \quad x \in \mathbb{R}^n. \quad (3.4)$$

Here it is first recalled that every $u \in \mathcal{S}'$ has finite order since, for $\psi \in \mathcal{S}$,

$$|\langle u, \psi \rangle| \leq c p_N(\psi), \quad (3.5)$$

$$p_N(\psi) = \sup \{ (1 + |x|)^N |D^\alpha u(x)| : x \in \mathbb{R}^n, |\alpha| \leq N \}. \quad (3.6)$$

Indeed, since $(1 + |x|)^N$ is finite on supp $\psi$ for $\psi \in C_0^\infty$, $u$ is of order $N$. To avoid a discussion of the converse, it will throughout be convenient to call the least integer $N$ fulfilling (3.5) the temperate order of $u$, written $N = \text{order } \mathcal{S}'(u)$.

Returning to (3.4), when the compact spectrum of $u$ results from Fourier multiplication, then the below $O(2^{kn})$-information on the constant will be used repeatedly in the present paper.

**Lemma 3.1.** Let $u \in \mathcal{S}'(\mathbb{R}^n)$ be arbitrary and $N \geq \text{order } \mathcal{S}'(\hat{u}).$ When $\psi \in C_0^\infty(\mathbb{R}^n)$ has support in $\mathcal{B}(0,R)$, then $w = \psi(2^{-k}D) u$ fulfils

$$w^*(N,R2^k;x) \leq C 2^{kN}(1 + |x|)^N, \quad k \in \mathbb{N}_0, \quad (3.7)$$

for a constant $C$ independent of $k$.

**Proof.** As $\psi(2^{-k}D) u(x) = \langle \hat{u}, \psi(2^{-k} \cdot) e^{i(x,\cdot)}(2\pi)^{-n} \rangle$, continuity of $\hat{u} : \mathcal{S} \rightarrow \mathbb{C}$ yields

$$|w(x)| \leq c \sup \{ (1 + |\xi|)^N |D_\xi^\alpha \psi(2^{-k} \xi) e^{i(x,\xi)}| : \xi \in \mathbb{R}^n, |\alpha| \leq N \}. \quad (3.8)$$

Since $(1 + |\xi|)^N |D_\xi^\alpha \psi(2^{-k} \xi)| \leq c'2^{kN-|\alpha|}$, Leibniz’ rule gives that $|w(x)| \leq c'' 2^{kN}(1 + |x|)^N$. Proceeding as before the lemma, the claim follows with $C = c'' \max(1,R^{-N})$. \hfill $\square$

Secondly, for the $a$-factor in (3.3) one has $F_a \in C(\mathbb{R}^n) \cap L_\infty(\mathbb{R}^n)$ and an estimate highly reminiscent of the Mihlin–Hörmander conditions for Fourier multipliers:

**Theorem 3.2.** Assume $a(x,\eta)$ is in $S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n)$ and let $F_a(N,R;x)$ be given by (3.3) for parameters $R,N > 0$, with the auxiliary function taken as $\chi = \psi(R^{-1} \cdot)$ for $\psi \in C_0^\infty(\mathbb{R}^n)$ equalling 1 in a set with non-empty interior. Then one has for all $x \in \mathbb{R}^n$ that

$$0 \leq F_a(x) \leq c_{n,N} \sum_{|\alpha| \leq [N+\frac{n}{2}] + 1} \left( \int_{\text{supp } \psi} |R^{|\alpha|} D_\eta^\alpha a(x,\eta)|^2 \frac{d\eta}{R^n} \right)^{1/2}. \quad (3.9)$$
Corollary 3.3. Assume \( a \in S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n) \) and let \( N, R \) and \( \psi \) be as in Theorem 3.2. When \( R \geq 1 \) there is a seminorm \( p \) on \( S^d_{1,1} \) and a constant \( c > 0 \), that depends only on \( n, N \) and \( \psi \), such that

\[
0 \leq F_a(x) \leq c_1 p(a) R^{\max(d,\lceil N+n/2 \rceil)+1} \quad \text{for all} \quad x \in \mathbb{R}^n.
\] (3.10)

If \( \psi(\eta) \neq 0 \) only holds in a corona \( 0 < \theta_0 \leq |\eta| \leq \Theta_0, \) and \( \psi(\eta) = 1 \) holds for \( \theta_1 \leq |\eta| \leq \Theta_1, \) then

\[
0 \leq F_a(x) \leq c_1 p(a) R^d \quad \text{for all} \quad x \in \mathbb{R}^n,
\] (3.11)

whereby \( c_0 = c_1 \max(1, \theta_0^{d-N-\lfloor n/2 \rfloor-1}, \theta_0^d) \).

The above asymptotics is \( O(R^d) \) for \( R \to \infty \) if \( d \) is large. This can be improved when \( a(x, \eta) \) is modified by removing the low frequencies in the \( x \)-variable (cf. the \( a^{(3)} \)-term in Section 5 below). In fact, with a second spectral quantity \( Q > 0 \), the following is contained in \([28, \text{Cor. 4.4}]\):

Corollary 3.4. When \( a_Q(x, \eta) = \varphi(Q^{-1}D_x) a(x, \eta) \) for some \( a \in S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n) \) and \( \varphi \in C_0^\infty(\mathbb{R}^n) \) with \( \varphi = 0 \) in a neighbourhood of \( \xi = 0 \), then there is a seminorm \( p \) on \( S^d_{1,1} \) and constants \( c_M \), depending only on \( M, n, N, \) and \( \psi \), such that

\[
0 \leq F_{a_Q}(N, R; x) \leq c_M p(a) Q^{-M} R^{\max(d+M,\lceil N+n/2 \rceil)+1} \quad \text{for} \quad M, Q, R > 0.
\] (3.12)

Here \( d+M \) can replace \( \max(d+M,\lceil N+n/2 \rceil+1) \) when the auxiliary function \( \psi \) in \( F_{a_Q} \) fulfils the corona condition in Corollary 3.3.

Remark 3.5. By the proofs in \([28]\), the seminorms in Corollaries 3.3 and 3.4 may be chosen in the same way for all \( d \), namely \( p(a) = \sum |a| \leq \lceil N+n/2 \rceil+1 p(a,0)(a) \); cf. (2.2).

3.2. Estimates of Frequency Modulated Operators. The results in the previous section easily give the following, which later in Sections 5 and 6 will be used repeatedly.

Proposition 3.6. For \( a(x, \eta) \) in \( S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n) \), \( v \in \mathcal{S}'(\mathbb{R}^n) \) and arbitrary \( \Phi, \Psi \in C_0^\infty(\mathbb{R}^n) \), for which \( \Psi \) is constant in a neighbourhood of the origin and with its support in \( \overline{B}(0, R) \) for some \( R \geq 1 \), there is for \( k \in \mathbb{N}_0 \) and \( N \geq \text{order}_{\mathcal{F}}(\mathcal{F} v) \), cf. (3.6), a polynomial bound

\[
\left| \text{OP}(\Phi(2^{-k}D_x)a(x, \eta)\Psi(2^{-k}\eta))v(x) \right| \leq C(k)(1+|x|)^N,
\] (3.13)

whereby

\[
C(k) = \begin{cases} 
2^{2k(N+d)+} & \text{for } N+d \neq 0, \\
c_k & \text{for } N+d = 0.
\end{cases}
\] (3.14)

For \( 0 \notin \text{supp} \Psi \) this may be sharpened to \( C(k) = 2^{2k(N+d)} \) for all values of \( N+d \).

Proof. In this proof it is convenient to let \( a^k(x, \eta) = \Phi(2^{-k}D_x)a(x, \eta) \) and \( \psi^k = \Psi(2^{-k}D) \). Then the factorisation inequality (3.1) gives

\[
|a^k(x, D)\psi^k(x)| \leq F_{a^k}(N, R^2^k; x) \cdot (\psi^k)^*(N, R^2^k; x).
\] (3.15)

Since \( N \geq \text{order}_{\mathcal{F}}(\psi) \), Lemma 3.1 gives \( (\psi^k)^*(N, R^2^k; x) \leq C(2^{kn}(1+|x|)^N, x \in \mathbb{R}^n.\)
In case \( 0 \notin \text{supp } \Psi \), the auxiliary function \( \chi = \psi(\cdot/(R2^k)) \) used in \( F_{\alpha} \), cf. Theorem 3.2, can be so chosen that it fulfils the corona condition in Corollary 3.3; e.g. it is possible to have \( \Theta_1 = 1 \) and \( \theta_1 = r/R \) when \( \Psi \equiv 0 \) on \( B(0, r) \). Since Remark 3.5 implies \( p(a^k) \leq p(a) \int |\mathcal{F}^{-1}\Phi(y)| \, dy \),
\[
0 \leq F_{\alpha}(N, R2^j; x) \leq c_0 \|\mathcal{F}^{-1}\Phi\|_1 \|p(a)R^d2^{kd}\tag{3.16}
\]
When combined with the above, this inequality yields the claim in case 0

In the general case one has \( v^k = v_k + v_{k-1} + \cdots + v_1 + v^0 \), whereby \( v_j \) denotes the difference \( v^j - v^{j-1} = \Psi(2^{-j}D)v - \Psi(2^{-j+1}D)v \). Via (3.1) this gives the starting point
\[
|a^k(x, D)v^0(x)| \leq |a^k(x, D)v^0(x)| + \sum_{j=1}^k F_{\alpha}(N, R2^j; x)v^j_s(N, R2^j; x)\tag{3.17}
\]
As \( \Psi = \Psi - \Psi(2^j) \) does not have 0 in its support, the above shows that for \( j = 1, \ldots, k \) one has \( F_{\alpha}(N, R2^j; x) \leq c_0 \|\mathcal{F}^{-1}\Phi\|_1 \|p(a)R^d2^{jd}\). Lemma 3.1 yields polynomial bounds of \( v^*_s \), say with a constant \( C' \), so the sum on the right-hand side of (3.17) is estimated, for \( d + N \neq 0 \), by
\[
\sum_{j=1}^k c_0C'R^d|p(a)(1 + |x|)^N| \leq \frac{c_0C'R^d}{2^{d+N+1}}|p(a)(1 + |x|)^N2^{(k+1)(N+d)}\tag{3.18}
\]
In case \( d + N = 0 \) the \( k \) bounds are equal.

The remainder in (3.17) fulfils \( |a^k(x, D)v^0(x)| \leq c_1R^N(1 + |x|)^N \) for a large \( N' \); cf. the first part of Corollary 3.3 and Lemma 3.1. Altogether \( |a^k(x, D)v^k(x)| \leq C(k)(1 + |x|)^N \[4\]

4. ADJOINTS OF TYPE 1,1-OPERATORS

For classical pseudo-differential operators \( a(x, D) : \mathcal{S} \to \mathcal{S} \) it is well known that the adjoint \( a(x, D)^* : \mathcal{S} \to \mathcal{S} \) has symbol \( a_s^*(x, \eta) = e^{iDx - D^\eta}a(x, \eta) \), and that \( a \mapsto a^* \) sends e.g. \( S^d_{1,0} \) into itself.

4.1. The Basic Lemma. In order to show that the twisted diagonal condition (2.27) also implies continuity \( a(x, D) : \mathcal{S} \to \mathcal{S} \), a basic result on the adjoint symbols is recalled from [18] and [20, Lem. 9.4.1]:

\textbf{Lemma 4.1.} When \( a(x, \eta) \) is in \( S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n) \) and for some \( B \geq 1 \) satisfies the twisted diagonal condition (2.27), then the adjoint symbol \( b(x, \eta) = e^{iDx - D^\eta}a(x, \eta) \) is in \( S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n) \) in this case and
\[
\wedge b(\xi, \eta) = 0 \quad \text{when} \quad |\xi + \eta| > B(|\eta| + 1). \tag{4.1}
\]
Moreover,
\[
|D^\alpha_\eta D^\beta_\xi b(x, \eta)| \leq C_{\alpha \beta}(a)B(1 + B^{-|\alpha|+|\beta|})(1 + |\eta|)^{d-|\alpha|+|\beta|}, \tag{4.2}
\]
for certain seminorms \( C_{\alpha \beta} \) that are continuous on \( S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n) \) and do not depend on \( B \).

The twisted diagonal condition (2.27) implies that \( a^*(x, D) = b(x, D) \) is a map \( \mathcal{S} \to \mathcal{S} \), as it is of type 1,1 by Lemma 4.1, so then \( a(x, D) \) has the continuous linear extension \( b(x, D)^* : \mathcal{S} \to \mathcal{S} \). It is natural to expect that this coincides with the definition of \( a(x, D) \) by vanishing frequency modulation:
Proposition 4.2. If \( a(x, \eta) \in S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n) \) fulfils (2.27), then \( a(x, D) \) is a continuous linear map \( \mathcal{S}'(\mathbb{R}^n) \rightarrow \mathcal{S}'(\mathbb{R}^n) \) that equals the adjoint of \( b(x, D) : \mathcal{S}(\mathbb{R}^n) \rightarrow \mathcal{S}(\mathbb{R}^n) \), when \( b(x, \eta) \) is the adjoint symbol as in Lemma 4.1.

Proof. When \( \psi \in C_0^\infty(\mathbb{R}^n) \) is such that \( \psi = 1 \) in a neighbourhood of the origin, a simple convolution estimate (cf. [26, Lem. 2.1]) gives that in the topology of \( S^{d+1}_{1,1} \),
\[
\psi(2^{-m}D_x)a(x, \eta)\psi(2^{-m}\eta) \rightarrow a(x, \eta) \quad \text{for} \quad m \to \infty.
\] (4.3)
Since the supports of the partially Fourier transformed symbols
\[
\psi(2^{-m}\xi)\mathcal{F}_{x \rightarrow \xi}a(\xi, \eta)\psi(2^{-m}\eta), \quad m \in \mathbb{N},
\] (4.4)
are contained in \( \text{supp} \mathcal{F}_{x \rightarrow \xi}a(\xi, \eta) \), clearly this sequence also fulfils (2.27) for the same \( B \). As the passage to adjoint symbols by (4.2) is continuous from the metric subspace of \( S^d_{1,1} \) fulfilling (2.27) to \( S^{d+1}_{1,1} \), one therefore has that
\[
b_m(x, \eta) := e^{iD_xD_\eta} \left( \psi(2^{-m}D_x)a(x, \eta)\psi(2^{-m}\eta) \right) \xrightarrow{m \to \infty} e^{iD_xD_\eta}a(x, \eta) =: b(x, \eta).
\] (4.5)
Combining this with the fact that \( b(x, D) \) as an operator on the Schwartz space depends continuously on the symbol, one has for \( u \in \mathcal{S}'(\mathbb{R}^n) \), \( \varphi \in \mathcal{S}(\mathbb{R}^n) \),
\[
(b(x, D)^*u|\varphi) = (u|\lim_{m \to \infty} \text{OP}(b_m(x, \eta))\varphi)
= \lim_{m \to \infty} (\text{OP}(\psi(2^{-m}D_x)a(x, \eta)\psi(2^{-m}\eta))u|\varphi).
\] (4.6)
As the left-hand side is independent of \( \psi \) the limit in (2.5) is so, hence the definition of \( a(x, D) \) gives that every \( u \in \mathcal{S}'(\mathbb{R}^n) \) is in \( D(a(x, D)) \) and \( a(x, D)u = b(x, D)^*u \) as claimed. \( \square \)

The mere extendability to \( \mathcal{S}' \) under the twisted diagonal condition (2.27) could have been observed already in [18, 20], but the above result seems to be the first sufficient condition for a type 1,1-operator to be defined on the entire \( \mathcal{S}'(\mathbb{R}^n) \).

4.2. The Self-Adjoint Subclass \( \tilde{S}^d_{1,1} \). Proposition 4.2 shows that (2.27) suffices for \( D(a(x, D)) = \mathcal{S}' \). But (2.27) is too strong to be necessary; a vanishing to infinite order along \( \mathcal{S} \) should suffice.

In this section, the purpose is to prove that \( a(x, D) : \mathcal{S}' \rightarrow \mathcal{S}' \) is continuous if more generally the twisted diagonal condition of order \( \sigma \), that is (2.33), holds for all \( \sigma \in \mathbb{R} \).

This will supplement Hörmander’s investigation in [18, 19, 20], from where the main ingredients are recalled. Using (2.29) and \( \mathcal{F}_{x \rightarrow \xi} \) one has in \( \mathcal{S}'(\mathbb{R}^n \times \mathbb{R}^n) \),
\[
a(x, \eta) = (a(x, \eta) - a_{\chi,1}(x, \eta)) + \sum_{v=0}^\infty (a_{\chi,2-v}(x, \eta) - a_{\chi,2-v-1}(x, \eta)).
\] (4.7)
Here the first term \( a(x, \eta) - a_{\chi,1}(x, \eta) \) fulfils (2.27) for \( B = 1 \), so Proposition 4.2 applies to it. Introducing \( e_x(x, D) \) like in [20, Sect. 9.3] as
\[
\hat{e}_x(x, \eta) = \hat{a}_{\chi,\varepsilon}(\xi, \eta) - \hat{a}_{\chi,\varepsilon/2}(\xi, \eta) = (\chi(\xi + \eta, \varepsilon\eta) - \chi(\xi + \eta, \varepsilon\eta/2))\hat{a}(x, \eta),
\] (4.8)
it is useful to infer from the choice of $\chi$ that
\[
\text{supp } \hat{e}_\epsilon \subset \{ (\xi, \eta) \mid \frac{\xi}{\epsilon} |\eta| \leq \max(1, |\xi + \eta|) \leq |\eta| \}.
\] (4.9)

In particular this yields that $e_\epsilon = 0$ when $1 + |\xi + \eta| < |\eta|/\epsilon$, so $e_\epsilon$ fulfills (2.27) for $B = 4/\epsilon$. Hence the terms $e_{2^{-v}}$ in (4.7) do so for $B = 2^{v+2}$.

The next result characterises the $a \in S^d_{1,1}$ for which the adjoint symbol $a^*$ is again in $S^d_{1,1}$; cf. the below condition (i). As adjoining is an involution, these symbols constitute the class
\[
\mathcal{S}^d_{1,1} := S^d_{1,1} \cap (S^d_{1,1})^*.
\] (4.10)

**Theorem 4.3.** When $a(x, \eta)$ is a symbol in $S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n)$ the following properties are equivalent:

(i) The adjoint symbol $a^*(x, \eta)$ is also in $S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n)$.
(ii) For arbitrary $N > 0$ and $\alpha, \beta$ there is some constant $C_{\alpha, \beta, N}$ such that
\[
|D^\alpha_{\eta} D^\beta_x a(x, \eta)| \leq C_{\alpha, \beta, N} e^N (1 + |\eta|)^{d-|\alpha|+|\beta|} \quad \text{for } 0 < \epsilon < 1.
\] (4.11)
(iii) For all $\sigma \in \mathbb{R}$ there is a constant $c_{\alpha, \sigma}$ such that for $0 < \epsilon < 1$
\[
\sup_{R > 0, \lambda \in \mathbb{R}^n} R^{\sigma-d} \left( \int_{|\eta| \leq 2R} |D^\alpha_{\eta} D^\beta_x a(x, \eta)|^2 \frac{d\eta}{R^n} \right)^{1/2} \leq c_{\alpha, \sigma} e^\sigma + \frac{\beta-|\alpha|}{2}.
\] (4.12)

In the affirmative case $a \in \mathcal{S}^d_{1,1}$, cf. (4.10), and $a^*$ fulfills an estimate
\[
|D^\alpha_{\eta} D^\beta_x a^*(x, \eta)| \leq (C_{\alpha, \beta}(a) + C_{\alpha, \beta, N})(1 + |\eta|)^{d-|\alpha|+|\beta|}
\] (4.13)
for a certain continuous seminorm $C_{\alpha, \beta}$ on $S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n)$ and some finite sum $C_{\alpha, \beta, N}$ of constants fulfilling the inequalities in (ii).

It should be observed from (i) that $a(x, \eta)$ fulfills condition (ii) or (iii) if and only if $a^*(x, \eta)$ does so — whereas neither (ii) nor (iii) make this obvious. But (ii) immediately gives the (expected) inclusion $S^d_{1,1} \subset \mathcal{S}^d_{1,1}$ for $d' > d$. Condition (iii) is close in spirit to the Mihlin–Hörmander multiplier theorem, and it is useful for the estimates shown later in Section 6.

**Remark 4.4.** Conditions (ii), (iii) both hold either for all $\chi$ satisfying (2.33) or for none, for (i) does not depend on $\chi$. It suffices to verify (ii) or (iii) for $0 < \epsilon < \epsilon_0$ for some convenient $\epsilon_0 \in [0, 1]$. This is implied by Lemma 2.13 since every power $\epsilon^p$ is bounded on the interval $[\epsilon_0, 1]$.

Theorem 4.3 was undoubtedly known to Hörmander, for he stated the equivalence of (i) and (ii) explicitly in [18, Thm. 4.2] and [20, Thm. 9.4.2] and gave brief remarks on (iii) in the latter. Equivalence with continuous extensions $H^{s+d} \to H^s$ for all $s \in \mathbb{R}$ was also shown. However, the exposition there left a considerable burden of verification to the reader.

Moreover, Theorem 4.3 was used without proof in a main $L_p$-theorem in [29], and below in Section 4.2.4 a corollary to the proof will follow and decisively enter the first proof of $\mathcal{S}'$-continuity. Hence full details on the main result in Theorem 4.3 should be in order here:
4.2.1. Equivalence of (ii) and (iii). That (ii) implies (iii) is seen at once by insertion, taking \( \beta = 0 \) and \( N = \sigma + \frac{d}{2} - |\alpha| \).

Conversely, note first that \( |\xi + \eta| \leq \varepsilon|\eta| \) in the spectrum of \( a_{X,\varepsilon}(\cdot, \eta) \). That is, \( |\xi| \leq (1 + \varepsilon)|\eta| \) so Bernstein’s inequality gives

\[
|D^\beta D^\alpha a_{X,\varepsilon}(x, \eta)| \leq ((1 + \varepsilon)|\eta|)^{|\beta|} \sup_{x \in \mathbb{R}} |D^\alpha a_{X,\varepsilon}(x, \eta)|. \tag{4.14}
\]

Hence \( C_{\alpha,\beta,N} = 2^{|\beta|} C_{\alpha,0,N} \) is possible, so it suffices to prove (iii) \( \Rightarrow \) (ii) only for \( \beta = 0 \).

For the corona \( 1 \leq |\xi| \leq 2 \) Sobolev’s lemma gives for \( f \in C^\infty(\mathbb{R}^n) \),

\[
|f(\xi)| \leq c_1 \left( \sum_{|\beta| \leq |n/2| + 1} \int_{1 \leq |\zeta| \leq 2} |D^\beta f(\zeta)|^2 d\zeta \right)^{1/2}. \tag{4.15}
\]

Substituting \( D^\alpha a_{X,\varepsilon}(x, R\zeta) \) and \( \zeta = \eta/R \), whereby \( R \leq |\eta| \leq 2R, R > 0 \), yields

\[
|D^\alpha a_{X,\varepsilon}(x, \eta)| \leq c_1 \left( \sum_{|\beta| \leq |n/2| + 1} R^{2|\beta|} \left( \int_{R \leq |\eta| \leq 2R} |D^{\alpha+\beta} a_{X,\varepsilon}(x, \eta)|^2 \frac{d\eta}{R^n} \right)^{1/2} \right)
\leq c_1 \left( \sum_{|\beta| \leq |n/2| + 1} R^{2d-2|\alpha|} C_{\alpha+\beta,\sigma}^2 \epsilon^{2(\sigma + \frac{n}{2} - |\alpha| - |\beta|)} \right)^{1/2}
\leq c_1 \left( \sum_{|\beta| \leq |n/2| + 1} C_{\alpha+\beta,\sigma}^2 \epsilon^{\sigma - 1 - |\alpha|} R^{d - |\alpha|} \right). \tag{4.16}
\]

Here \( R^{d - |\alpha|} \leq (1 + |\eta|)^{d - |\alpha|} \) for \( d \geq |\alpha| \), that leads to (ii) as \( \sigma \in \mathbb{R} \) can be arbitrary.

For \( |\alpha| > d \) it is first noted that, by the support condition on \( X \), one has \( a_{X,\varepsilon}(x, \eta) \neq 0 \) only for \( 2R \geq |\eta| \geq \varepsilon^{-1} > 1 \). But \( R \geq 1/2 \) yields \( R^{d - |\alpha|} \leq \left( \frac{1}{2} + 2R \right)^{d - |\alpha|} \leq 6|\alpha|^{-d} (1 + |\eta|)^{d - |\alpha|} \), so (ii) follows from the above.

4.2.2. The implication (ii) \( \Rightarrow \) (i) and the estimate. The condition (ii) is exploited for each term in the decomposition (4.7). Setting \( b_v(x, \eta) = e_{2-v}^*(x, \eta) \) it follows from Lemma 4.1 that \( b_v \) is in \( S_{1,1}^d \) by the remarks after (4.9), cf (4.7) ff, and (4.2) gives

\[
|D^\alpha D^\beta b_v(x, \eta)| \leq C_{\alpha,\beta}(e_{2-v})^{2v+2} (1 + 2^{(v+2)(d - |\alpha| + |\beta|)}) (1 + |\eta|)^{d - |\alpha| + |\beta|}. \tag{4.17}
\]

Now (ii) implies that \( C_{\alpha,\beta}(a_{X,2-v}) \leq C_{\alpha,\beta,N}^2 \) for all \( N > 0 \) (with other constants \( C_{\alpha,\beta,N}^2 \) as the seminorms \( C_{\alpha,\beta} \) may contain derivatives of higher order than \( |\alpha| \) and \( |\beta| \)). Hence \( C_{\alpha,\beta}(e_{2-v}) \leq C_{\alpha,\beta,N}^2 \). It follows from this that \( \sum b_v \) converges to some \( b \) in \( S_{1,1}^d \) (in the Fréchet topology of this space), so that \( a^+(x, \eta) = b(x, \eta) \) is in \( S_{1,1}^d \). More precisely, (4.2) and the above yields for
\[ N = 2 + (d - |\alpha| + |\beta|) + \]
\[ \frac{|D_{\eta}^\alpha p_{\delta}(x, \eta)|}{(1 + |\eta|)d - |\alpha| + |\beta|} \leq 2^N C_{\alpha, \beta}(a - a_{\chi, 1}) + \sum_{v=0}^{\infty} C_{\alpha, \beta}(e_{2^{-v}})2^{v+2}(1 + 2^{(v+2)(d - |\alpha| + |\beta|)} + ) \]
\[ \leq 2^N C_{\alpha, \beta}(a - a_{\chi, 1}) + \sum_{v=0}^{\infty} 16 C_{\alpha, \beta, N} 2^{-v(N-1)}2^{(v+2)(d - |\alpha| + |\beta|)} \]
\[ \leq 2^N C_{\alpha, \beta}(a - a_{\chi, 1}) + 4^{N+2} C_{\alpha, \beta, N}. \]

Invoking continuity from Lemma 2.13 in the first term, the last part of the theorem follows.

4.2.3. Verification of (ii) \( \implies \) (i). It suffices to derive another decomposition
\[ a = A + \sum_{v=0}^{\infty} a_v, \]
(4.19)
in which \( A \in S^{-\infty} \) and each \( a_v \in S^d_{1,1} \) with \( \hat{a}_v(\xi, \eta) = 0 \) for \( 2^{v+1}|\xi + \eta| < |\xi| \) and seminorms
\[ C_{\alpha, \beta}(a_v) = O(2^{-v}) \]
for each \( N > 0 \).

Indeed, when \( \hat{\chi}(\xi + \eta, \xi + \eta) \neq 0 \) the triangle inequality gives \( |\xi + \eta| \leq |\xi + \eta| + \frac{1}{2} |\xi| \), whence \( |\xi + \eta|/(1 - \varepsilon) \leq |\xi| \), so that for one thing
\[ \hat{a}_v(\xi, \eta) = \chi(\xi + \eta, \xi + \eta)\hat{A}(\xi, \eta) + \sum_{2^{v+1} > (1 - \varepsilon)} \chi(\xi + \eta, \xi + \eta)\hat{a}_v(\xi, \eta). \]
(4.20)

Secondly, for each seminorm \( C_{\alpha, \beta} \) in \( S^d_{1,1} \) one has \( C_{\alpha, \beta}(a_{\chi, \varepsilon}) \leq e^{-|\alpha|} C_{\alpha, \beta}(a_v) \) by Lemma 2.13, so by estimating the geometric series by its first term, the above formula entails that
\[ C_{\alpha, \beta}(a_{\chi, \varepsilon}) \leq C_{\alpha, \beta}(A_{\chi, \varepsilon}) + \sum_{\varepsilon \geq 2^{v+1} > (1 - \varepsilon)} \frac{C_{2^{v+1} > (1 - \varepsilon)}}{e^{2^{v+1} > (1 - \varepsilon)}2^{(v+2)(d + 1 - |\alpha|)}} \leq C_{\alpha, \beta}(A_{\chi, \varepsilon}) + \frac{C_{\alpha, \beta}(A_{\chi, \varepsilon})}{e^{2^{v+1} > (1 - \varepsilon)}2^{(v+2)(d + 1 - |\alpha|)}}. \]
(4.21)

This gives the factor \( e^N \) in (ii) \( \forall \varepsilon < 1/2 \). For \( 1/2 < \varepsilon < 1 \) the series is \( O(\varepsilon^{-|\alpha|}) \) because \( 2^{-v} \leq 1 < 2\varepsilon/(1 - \varepsilon) \) for all \( v \). However, \( 1 < 2\varepsilon/(1 - \varepsilon) \) for all \( \varepsilon \in [0, 1] \). (It is seen directly that \( |A_{\chi, \varepsilon}| \leq e^N (1 + |\eta|)^d \) etc, for only the case \( \varepsilon \leq 1 \) is non-trivial, and then \( e^{-N} \leq (1 + |\eta|)^N \) while \( A \in S^{-\infty} \).

In the deduction of (4.19) one can use a Littlewood–Paley partition of unity, say \( 1 = \sum \Phi_{v} \) with dilated functions \( \Phi_{v}(\eta) = \Phi(2^{-v}\eta) \neq 0 \) only for \( \frac{11}{2^v} \leq |\eta| \leq \frac{13}{2^v} \) if \( v \geq 1 \). Beginning with a trivial split \( a_{\chi} = A_0 + A_1 \) into two terms for which \( A_0 \in S^{-\infty} \) and \( A_1 \in S^d_{1,1} \) such that
\[ A_1(x, \eta) = 0 \quad \text{for} \quad |\eta| < 1/2, \]
this gives
\[ \hat{a}_{\chi, \varepsilon}(\xi, \eta) = A_0(\xi, \eta) + \sum_{v=0}^{\infty} \Phi_{v}(\xi/|\eta|)A_1(\xi, \eta). \]
(4.22)

This yields the desired \( a_{\chi, \varepsilon}(x, \eta) \) by taking the adjoint of \( \mathcal{F}_{x}^{-1}(\Phi_{v}(\frac{\xi}{|\eta|})A_1(\xi, \eta)) \), that is, of the symbol \( \int |2^v\eta|^{n}\Phi(2^v|\eta|)|_{y}A_1(x-y, \eta) \) dy. Indeed, it follows directly from [18, Prop. 3.3] (where the proof uses Taylor expansion and vanishing moments of \( \Phi \) for \( \varepsilon \geq 1 \)) that \( a^*_{\chi, \varepsilon} \) belongs to \( S^d_{1,1} \).
with \((2^{Nv}a^*_v)_{v \in \mathbb{N}}\) bounded in \(S^d_{1,1}\) for all \(N > 0\). Therefore (4.22) gives (4.19) by inverse Fourier transformation. Moreover, \(\hat{a}^*_v(\xi, \eta)\) is for \(v \geq 1\) supported by the region
\[
\frac{11}{10}2^v|\eta| \leq |\xi| \leq \frac{13}{10}2^v|\eta|,
\]
where a fortiori \(1 + |\xi + \eta| \geq |\xi| - |\eta| \geq (\frac{11}{10}2^v - 1)|\eta| \geq \frac{1}{100}|\eta|\), so it is clear that
\[
\hat{a}^*_v(\xi, \eta) = 0 \quad \text{if} \quad 10(|\xi + \eta| + 1) < |\eta|.
\]
(4.24)
According to Lemma 4.1 this implies that \(a_v = a^*_v\) is also in \(S^d_{1,1}\) and that, because of the above boundedness in \(S^d_{1,1}\), there is a constant \(c\) independent of \(v\) such that
\[
|D^\alpha_{\eta}D^\beta_{\xi}a_v(x, \eta)| \leq C_{\alpha, \beta}(a_v)10(1 + 10^{d - |\alpha| + |\beta|})(1 + |\eta|)^{d - |\alpha| + |\beta|}
\]
\[
\leq c2^{-Nv}(1 + |\eta|)^{d - |\alpha| + |\beta|}.
\]
(4.25)
Therefore, the \(a_v\) tend rapidly to 0, which completes the proof of Theorem 4.3.

4.2.4. Consequences for the Self-Adjoint Subclass. One can set Theorem 4.3 in relation to the definition by vanishing frequency modulation, simply by elaborating on the above proof:

**Corollary 4.5.** On \(\widehat{S}_{1,1}^d(\mathbb{R}^n \times \mathbb{R}^n)\) the adjoint operation is stable with respect to vanishing frequency modulation in the sense that, when \(a \in \widehat{S}_{1,1}^d, \psi \in C^\infty_0(\mathbb{R}^n)\) with \(\psi = 1\) around 0, then
\[
(\psi(2^{-m}D_x)a(x, \eta)\psi(2^{-m}\eta))^* \xrightarrow{m \to \infty} a(x, \eta)^*
\]
holds in the topology of \(S_{1,1}^{d+1}(\mathbb{R}^n \times \mathbb{R}^n)\).

**Proof.** For brevity \(b_m(x, \eta) = \psi(2^{-m}D_x)a(x, \eta)\psi(2^{-m}\eta)\) denotes the symbol that is frequency modulated in both variables. The proof consists in insertion of \(a(x, \eta) - b_m(x, \eta)\) into (4.18), where the first sum tends to 0 for \(m \to \infty\) by majorised convergence.

Note that for each \(v \geq 0\) in the first sum of (4.18) one must control \(C_{\alpha, \beta}(e^{m \xi}\psi)\) for \(m \to \infty\) and
\[
e^{m \xi} = (\chi(\xi, \eta, 2^{-v}\eta) - \chi(\xi, \eta, 2^{-v+1}\eta))(1 - \psi(2^{-m}\xi)\psi(2^{-m}\eta))\hat{a}(\xi, \eta).
\]
(4.27)
To do so, a convolution estimate first gives \(p_{\alpha, \beta}(b_m) \leq c\sum_{\gamma \leq \alpha}p_{\gamma, \beta}(a)\), whence \((b_m)_{m \in \mathbb{N}}\) is bounded in \(S_{1,1}^d\). Similar arguments yield that \(b_m \to a\) in \(S_{1,1}^{d+1}\) for \(m \to \infty\); cf. [26, Lem. 2.1]. Moreover, for each \(v \geq 0\), every seminorm \(p_{\alpha, \beta}\) now on \(S_{1,1}^{d+1}\), gives
\[
p_{\alpha, \beta}(e_{-\xi}^{m \xi}) \leq p_{\alpha, \beta}((a - b_m)\chi_{2^{-v-1}\eta}) + p_{\alpha, \beta}((a - b_m)\chi_{2^{-v-1}\eta}).
\]
(4.28)
Here both terms on the right-hand side tend to 0 for \(m \to \infty\), in view of the continuity of \(a \mapsto a\chi_{\varepsilon}\) on \(S_{1,1}^{d+1}\), cf. Lemma 2.13. Hence \(C_{\alpha, \beta}(e_{-\xi}^{m \xi}) \to 0\) for \(m \to \infty\).

It therefore suffices to replace \(d\) by \(d+1\) in (4.18) and majorise. However, \(a \mapsto a\chi_{\varepsilon}\) commutes with \(a \mapsto b_m\) as maps in \(S_{1,1}^d(\mathbb{R}^n \times \mathbb{R}^n)\), so since \(a \in \widehat{S}_{1,1}^d\), it follows from (ii) that
\[
p_{\alpha, \beta}((a - b_m)\chi_{\varepsilon}) \leq p_{\alpha, \beta}(a\chi_{\varepsilon}) + c\sum_{\gamma \leq \alpha}p_{\gamma, \beta}(a\chi_{\varepsilon}) \leq (1 + c)(\sum_{\gamma \leq \alpha}C_{\gamma, \beta, N})e^N \leq C_{\alpha, \beta, N}e^N.
\]
(4.29)
Using this in the previous inequality, $C_{a,b}(e^{t_{2,v}}) \leq C2^{-vN}$ is obtained for $C$ independent of $m \in \mathbb{N}$. Now it follows from (4.18) that $b_{m}(x, \eta)^{\ast} \rightarrow a(x, \eta)^{\ast}$ in $S_{1,1}^{d+1}$ as desired. \hfill \Box

Thus prepared, the proof of Proposition 4.2 can now be repeated from (4.5) onwards, which immediately gives the first main result of the paper:

**Theorem 4.6.** When a symbol $a(x, \eta)$ of type 1,1 belongs to the class $\tilde{S}_{1,1}^{d}(\mathbb{R}^{n} \times \mathbb{R}^{n})$, as characterised in Theorem 4.3, then $a(x, D)$ is everywhere defined and continuous

$$a(x, D) : \mathcal{S}'(\mathbb{R}^{n}) \rightarrow \mathcal{S}'(\mathbb{R}^{n})$$

It equals the adjoint of $\text{OP}(e^{iD_{x}D_{\eta}a(x, \eta)}) : \mathcal{S} \rightarrow \mathcal{S}$.

Like for Proposition 4.2, there seems to be no previous attempts in the literature to obtain this clarification (Theorem 4.6 was stated without proof in [29]). However, it seems to be open whether (4.30) conversely implies that $a \in \tilde{S}_{1,1}^{d}$.

5. **Dyadic Corona Decompositions**

This section adopts Littlewood–Paley techniques to provide a passage to auxiliary operators $a^{(j)}(x, D)$, $j = 1, 2, 3$, which may be easily analysed with the pointwise estimates of Section 3.

5.1. **The Paradifferential Splitting.** Recalling the definition of type 1,1-operators in (2.5) and (2.6), it is noted that to each modulation function $\psi$, i.e. $\psi \in C_{0}^{\infty}(\mathbb{R}^{n})$ with $\psi = 1$ in a neighbourhood of 0, there exist $R > r > 0$ with $R \geq 1$ satisfying

$$\psi(\xi) = 1 \quad \text{for} \quad |\xi| \leq r; \quad \psi(\xi) = 0 \quad \text{for} \quad |\xi| \geq R. \quad (5.1)$$

For fixed $\psi$ it is convenient to take an integer $h \geq 2$ so large that $2R < r2^{h}$.

To obtain a Littlewood–Paley decomposition from $\psi$, set $\varphi = \psi - \psi(2\cdot)$. Then a dilation of this function is supported in a corona,

$$\text{supp} \varphi(2^{-k}) \subset \{ \xi \mid r2^{k-1} \leq |\xi| \leq R2^{k} \}, \quad \text{for} \ k \geq 1. \quad (5.2)$$

The identity $1 = \psi(x) + \sum_{k=1}^{\infty} \varphi(2^{-k}x)$ follows by letting $m \to \infty$ in the telescopic sum,

$$\psi(2^{-m}x) = \psi(x) + \varphi(\xi/2) + \cdots + \varphi(\xi/2^{m}). \quad (5.3)$$

Using this, functions $u(x)$ and symbols $a(x, \eta)$ will be localised to frequencies $|\eta| \approx 2^{j}$ as

$$u_{j} = \varphi(2^{-j}D)u, \quad a_{j}(x, \eta) = \varphi(2^{-j}D_{x})a(x, \eta). \quad (5.4)$$

Localisation to balls given by $|\eta| \leq R2^{j}$ are written with upper indices,

$$u^{j} = \psi(2^{-j}D)u, \quad a^{j}(x, \eta) = \psi(2^{-j}D_{x})a(x, \eta). \quad (5.5)$$

In addition $u_{0} = u^{0}$ and $a_{0} = a^{0}$; as an index convention they are all taken $\equiv 0$ for $j < 0$. (To avoid having two different meanings of sub- and superscripts, the dilations $\psi(2^{-j} \cdot)$ are written as such, with the corresponding Fourier multiplier as $\psi(2^{-j}D)$, and similarly for $\varphi$). Note that the corresponding operators are $a^{k}(x, D) = \text{OP}(\psi(2^{-k}D_{x})a(x, \eta))$ etc.
Inserting the relation (5.3) twice in (2.5), bilinearity gives

$$\text{OP}(\psi(2^{-m} D_x) a(x, \eta) \psi(2^{-m} \eta)) u = \sum_{j,k=0}^{m} a_j(x, D) u_k. \quad (5.6)$$

Of course the sum may be split in three groups having \( j \leq k - h, \ |j - k| < h \) and \( k \leq j - h \). For \( m \to \infty \) this yields the well-known paradifferential decomposition

$$a_{\psi}(x, D) u = a_{\psi}^{(1)}(x, D) u + a_{\psi}^{(2)}(x, D) u + a_{\psi}^{(3)}(x, D) u, \quad (5.7)$$

whenever \( a \) and \( u \) fit together such that the three series below converge in \( \mathcal{E}'(\mathbb{R}^n) \):

$$a_{\psi}^{(1)}(x, D) u = \sum_{k=h}^\infty \sum_{j \leq k - h} a_j(x, D) u_k = \sum_{k=h}^\infty a^{k-h}(x, D) u_k \quad (5.8)$$

$$a_{\psi}^{(2)}(x, D) u = \sum_{k=0}^\infty (a_{k-h+1}(x, D) u_k + \cdots + a_{k-1}(x, D) u_k + a_k(x, D) u_k) \quad (5.9)$$

$$a_{\psi}^{(3)}(x, D) u = \sum_{j=h}^\infty \sum_{j \leq k-h} a_j(x, D) u_k = \sum_{j=h}^\infty a_j(x, D) u^{j-h} \quad (5.10)$$

Note the shorthand \( a^{k-h}(x, D) \) for \( \sum_{j \leq k-h} a_j(x, D) = \text{OP}(\psi(2^{h-k} D_x) a(x, \eta)) \) etc. Using this and the index convention, the so-called symmetric term in (5.9) has the brief form

$$a_{\psi}^{(2)}(x, D) u = \sum_{k=0}^\infty \left((a_k - a^{k-h})(x, D) u_k + a_k(x, D) (u^{k-1} - u^{k-h})\right). \quad (5.11)$$

In the following the subscript \( \psi \) is usually dropped because this auxiliary function will be fixed (\( \psi \) was left out already in \( a_j \) and \( a^j \); cf. (5.4)–(5.5)). Note also that the above \( a^{(j)}(x, D) \) for now is just a convenient notation for the infinite series. The full justification of this operator notation will first result from Theorems 6.1–6.4 below.

**Remark 5.1.** It was tacitly used in (5.6) and (5.8)–(5.10) that one has

$$a_j(x, D) u_k = \text{OP}(a_j(x, \eta) \varphi(2^{-k} \eta)) u. \quad (5.12)$$

This is because, with \( \chi \in C_0^\infty \) equalling 1 on \( \text{supp} \varphi(2^{-k} \cdot) \), both sides are equal to

$$\text{OP}(a_j(x, \eta) \chi(\eta)) u_k. \quad (5.13)$$

Indeed, while this is trivial for the right-hand side of (5.12), where the symbol is in \( S^{-\infty} \), it is for the type 1,1-operator on the left-hand side of (5.12) a fact that follows at once from (2.14). Thus the inclusion \( \mathcal{F}^{-1} \mathcal{E}' \subset D(a(x, D)) \) in (2.14) is crucial for the simple formulae in the present paper. Analogously Definition 2.1 may be rewritten briefly as \( a(x, D) u = \lim_m d^m(x, D) u^m \).

The importance of the decomposition in (5.8)–(5.10) lies in the fact that the summands have localised spectra. E.g. there is a dyadic corona property:
Proposition 5.2. If \( a \in S_{1,1}^{d}(\mathbb{R}^{n} \times \mathbb{R}^{n}) \) and \( u \in \mathcal{S}'(\mathbb{R}^{n}) \), and \( r, R \) are chosen as in (5.1) for each auxiliary function \( \psi \), then every \( h \in \mathbb{N} \) such that \( 2R < r2^{h} \) gives

\[
\text{supp} \mathcal{F}(a^{k-h}(x, D)u_k) \subset \{ \xi | R_h 2^{k} \leq |\xi| \leq \frac{5R}{4}2^{k} \} \quad (5.14)
\]

\[
\text{supp} \mathcal{F}(a_k(x, D)u^{k-h}) \subset \{ \xi | R_h 2^{k} \leq |\xi| \leq \frac{5R}{4}2^{k} \}, \quad (5.15)
\]

whereby \( R_h = \frac{r}{2} - R2^{-h} > 0 \).

Proof. By (5.2) and the Spectral Support Rule, cf. the last part of Theorem B.1,

\[
\text{supp} \mathcal{F}(a^{k-h}(x, D)u_k) \subset \{ \xi + \eta \mid (\xi, \eta) \in \text{supp}(\psi_{h-k} \otimes 1)\hat{\alpha}, \ r2^{k-1} \leq |\eta| \leq R2^{k} \}. \quad (5.16)
\]

So by the triangle inequality every \( \zeta = \xi + \eta \) in the support fulfils, as \( h \geq 2 \),

\[
r2^{k-1} - R2^{-h} \leq |\zeta| \leq R2^{k-1} + R2^{k} \leq \frac{5}{4}R2^{k}. \quad (5.17)
\]

This shows (5.14), and (5.15) follows analogously. \( \square \)

To achieve simpler constants one could take \( h \) so large that \( 4R \leq r2^{h} \), which instead of \( R_h \) would allow \( r/4 \) (and \( 9R/8 \)). But the present choice of \( h \) is preferred in order to reduce the number of terms in \( a^{(2)}(x, D)u \).

In comparison the terms in \( a^{(2)}(x, D)u \) only satisfy a dyadic ball condition. Previously this was observed e.g. for functions \( u \in \bigcup H^{s} \) in [24], as was the fact that when the twisted diagonal condition (2.27) holds, then the situation improves for large \( k \). This is true for arbitrary \( u \):

Proposition 5.3. When \( a \in S_{1,1}^{d}(\mathbb{R}^{n} \times \mathbb{R}^{n}) \), \( u \in \mathcal{S}'(\mathbb{R}^{n}) \), and \( r, R \) are chosen as in (5.1) for each auxiliary function \( \psi \), then every \( h \in \mathbb{N} \) such that \( 2R < r2^{h} \) gives

\[
\text{supp} \mathcal{F}(a_k(x, D)(u^{k-1} - u^{k-h})) \bigcup \text{supp} \mathcal{F}((a^{k} - a^{k-h})(x, D)u_k)
\]

\[
\subset \{ \xi \in \mathbb{R}^{n} | |\xi| \leq 2R2^{k} \} \quad (5.18)
\]

If (2.27) holds for some \( B \geq 1 \), then the support is contained in the annulus

\[
\{ \xi \mid \frac{r}{2h+1}B2^{k} \leq |\xi| \leq 2R2^{k} \} \quad \text{for all } k \geq h + 1 + \log_{2}(\frac{B}{r}). \quad (5.19)
\]

Proof. As in Proposition 5.2, \( \text{supp} \mathcal{F} a_k(x, D)(u^{k-1} - u^{k-h}) \) is seen to be a subset of

\[
\{ \xi + \eta \mid (\xi, \eta) \in \text{supp}(\varphi_k \otimes 1)\hat{\alpha}, \ r2^{k-h} \leq |\eta| \leq R2^{k-1} \}. \quad (5.20)
\]

Thence any \( \zeta \) in the support fulfils \( |\zeta| \leq R2^{k} + R2^{k-1} = (3R/2)2^{k} \). If (2.27) holds, then one has \( B(1 + |\xi + \eta|) \geq |\eta| \) on \( \text{supp} \mathcal{F} x^{-\xi}a \), so for all \( k \) larger than the given limit

\[
|\zeta| \geq \frac{1}{B}|\eta| - 1 \geq \frac{1}{B}r2^{k-h} - 1 \geq (\frac{r}{2B} - 2^{-k})2^{k} \geq \frac{r}{2B+1}2^{k}. \quad (5.21)
\]

The term \( (a^{k} - a^{k-h})(x, D)u_k \) is analogous but will cause \( 3R/2 \) to be replaced by \( 2R \). \( \square \)
Remark 5.4. The inclusions in Propositions 5.2 and 5.3 have been a main reason for the introduction of the paradifferential splitting (5.7) in the 1980’s, but they were then only derived for elementary symbols; cf. [6, 7, 48]. With the Spectral Support Rule, cf. Theorem B.1, this restriction is redundant; cf. also the remarks to (1.24) in the introduction.

5.2. Polynomial Bounds. In the treatment of \(a^{(1)}(x,D)u\) and \(a^{(3)}(x,D)u\) in (5.8) and (5.10) one may conveniently commence by observing that, according to Proposition 5.2, the terms in these series fulfil condition (A.1) in Lemma A.1 for \(\theta_0 = \theta_1 = 1\).

So to deduce their convergence from Lemma A.1, it remains to obtain the polynomial bounds in (A.2). For this it is natural to use the efficacy of the pointwise estimates in Section 3:

**Proposition 5.5.** If \(a(x,\eta)\) is in \(S^1_{1,1}(\mathbb{R}^n \times \mathbb{R}^n)\) and \(N \geq \text{order}_{\mathcal{F}}(\mathcal{P}u)\) fulfils \(d + N \neq 0\), then

\[
|a^{k-h}(x,D)u_k(x)| \leq c2^{k(N+d)}(1 + |x|)^N, \quad (5.22)
\]

\[
|a_k(x,D)u^{k-h}(x)| \leq c2^{k(N+d)}(1 + |x|)^N, \quad (5.23)
\]

\[
|(a^k - a^{k-h})(x,D)u_k(x)| \leq c2^{k(N+d)}(1 + |x|)^N, \quad (5.24)
\]

\[
|a_k(x,D)(u^{k-1} - u^{k-h})(x)| \leq c2^{k(N+d)}(1 + |x|)^N. \quad (5.25)
\]

**Proof.** The second inequality follows by taking the two cut-off functions in Proposition 3.6 as \(\Phi = \varphi\) and \(\Psi = \psi(2^{-h} \cdot)\). The first claim is seen by interchanging their roles, i.e. for \(\Phi = \psi(2^{-h} \cdot)\) and \(\Psi = \varphi\); the latter is 0 around the origin so \(N + d\) is obtained without the positive part.

Clearly similar estimates hold for the terms in \(a^{(2)}(x,D)u\). E.g., taking \(\psi - \psi(2^{-h} \cdot)\) and \(\varphi\), respectively, as the cut-off functions in Proposition 3.6, one finds for \(k \geq h\) the estimate in (5.24). Note that the positive part can be avoided for \(0 \leq k < h\) by using a sufficiently large constant. \(\square\)

The difference in the above estimates appears because \(u_k\) in (5.22) has spectrum in a corona. However, one should not confound this with spectral inclusions like (A.1) that one might obtain after application of \(a^{k-h}(x,D)\), for these are irrelevant for the pointwise estimates here.

5.3. Induced Paradifferential Operators. Although (5.8)–(5.10) yield a well-known splitting, the operator notation \(a^{(j)}(x,D)\) requires justification in case of type 1,1-operators.

Departing from the right hand sides of (5.8)–(5.10) one is via (5.12) led directly to the symbols

\[
a^{(1)}(x,\eta) = \sum_{k=h}^{\infty} a^{k-h}(x,\eta)\varphi(2^{-k}\eta) \quad (5.26)
\]

\[
a^{(3)}(x,\eta) = \sum_{j=h}^{\infty} a_j(x,\eta)\psi(2^{-(j-h)}\eta). \quad (5.27)
\]

In addition, letting \(\delta_{k \geq h}\) stand for 1 when \(k \geq h\) and for 0 in case \(k < h\),

\[
a^{(2)}(x,\eta) = \sum_{k=1}^{\infty} \left( (a^k(x,\eta) - a^{k-h}(x,\eta))\varphi(2^{-k}\eta) + a_k(x,\eta)(\psi(2^{-(k-1)}\eta) - \psi(2^{-(k-h)}\eta)\delta_{k \geq h}) \right) + a^0(x,\eta)\psi(\eta) \quad (5.28)
\]
These three series converge in the Fréchet space $S_1^{d+1}([\mathbb{R}^n \times \mathbb{R}^n])$, for the sums are locally finite. Therefore it is clear that
\[ a(x, \eta) = a^{(1)}(x, \eta) + a^{(2)}(x, \eta) + a^{(3)}(x, \eta), \] (5.29)
where some of the partially Fourier transformed symbols have conical supports,
\[ \hat{a}^{(1)}(\xi, \eta) \neq 0 \implies |\xi| \leq \frac{2R}{r_2^n}|\eta|, \quad \hat{a}^{(3)}(\xi, \eta) \neq 0 \implies |\eta| \leq \frac{2R}{r_2^n}|\xi|. \] (5.30)
This well-known fact follows from the supports of $\psi$ and $\varphi$. But a sharper exploitation gives

**Proposition 5.6.** For each $a \in S_1^d$ and every modulation function $\psi \in C_0^\infty(\mathbb{R}^n)$ in (5.1), the associated symbols $a^{(1)}_\psi(x, \eta)$ and $a^{(3)}_\psi(x, \eta)$ fulfill the twisted diagonal condition (2.27).

**Proof.** When $\hat{a}^{(3)}(\xi, \eta) \neq 0$ it follows from (5.30), which in particular yields $|\eta| < |\xi|$, that
\[ |\xi + \eta| \geq |\xi| - |\eta| \geq |\xi| (1 - \frac{2R}{r_2^n}) > |\eta| (1 - \frac{2R}{r_2^n}). \] (5.31)
Therefore $\hat{a}^{(3)}(\xi, \eta) = 0$ whenever $B_1|\xi + \eta| < |\eta|$ holds for $B_1 = (1 - \frac{2R}{r_2^n})^{-1}$; a fortiori (2.27) is fulfilled with $B = B_1 > 1$. The case of $a^{(1)}$ is a little simpler. □

To elucidate the role of the twisted diagonal, note that the lower bound in Proposition 5.2 reappears by using $|\xi| \geq r2^{k-1}$ in the middle of (5.31).

Anyhow, it is a natural programme to verify that $u \in \mathcal{S}'$ belongs to the domain of the operator $a^{(j)}(x, D)$ precisely when the previously introduced series denoted $a^{(j)}(x, D)u$ converges; cf. (5.8)–(5.10). In view of the definition by vanishing frequency modulation in (2.5) ff, this will necessarily be lengthy because a second modulation function $\Psi$ has to be introduced.

To indicate the details for $a^{(1)}(x, \eta)$, let $\psi, \Psi \in C_0^\infty(\mathbb{R}^n)$ be equal to 1 around the origin, and let $\psi$ be used as the fixed modulation function entering $a^{(1)}(x, D) = a^{(1)}_\psi(x, D)$ in (5.8); and set $\varphi = \psi - \psi(2 \cdot)$. The numbers $r, R$ and $h$ are then chosen in relation to $\psi$ as in (5.1).

Moreover, $\Psi$ is used for the frequency modulation made when Definition 2.1 is applied to $a^{(1)}_\psi(x, D)$. This gives the following identity in $S_1^d$, where prime indicates a finite sum,
\[
\Psi(2^{-m}D_x)a^{(1)}(x, \eta)\Psi(2^{-m}\eta) = \sum_{k=h}^{m+\mu} a^{k-h}(x, \eta)\varphi(2^{-k}\eta) \\
+ \sum_k' \Psi(2^{-m}D_x)a^{k-h}(x, \eta)\varphi(2^{-k}\eta)\Psi(2^{-m}\eta). \] (5.32)

Indeed, if $\lambda, \Lambda > 0$ fulfill that $\Psi(\eta) = 1$ for $|\eta| \leq \lambda$ while $\Psi = 0$ for $|\eta| \geq \Lambda$, the support of $\varphi(2^{-k}\eta)$ in (5.26) lies by (5.2) in one of the ‘harmless’ level sets $\Psi(2^{-m}\eta) = 1$ or $\Psi(2^{-m}\eta) = 0$ when, respectively,
\[
R2^k \leq \lambda 2^m \quad \text{or} \quad r2^{k-1} \geq \Lambda 2^m. \] (5.33)
That is, $\text{supp } \varphi(2^{-k} \cdot)$ is contained in these level sets unless $k$ fulfills
\[
m + \log_2(\lambda/R) < k < m + 1 + \log_2(\Lambda/r). \] (5.34)
Therefore the primed sum has at most 1 + log₂ \( \frac{RA}{\lambda r} \) terms, independently of the parameter \( m \); in addition \( \Psi(2^{-m}\eta) \) and \( \Psi(2^{-m}D_x) \) disappear from the other terms, as stated in (5.32).

Consequently, with \( \mu = \lfloor \log_2(\lambda/R) \rfloor \) and \( k = m + l \), for \( l \in \mathbb{Z} \), one has for \( u \in \mathcal{S}'(\mathbb{R}^n) \) that

\[
\text{OP}(\Psi(2^{-m}D_x)a^{(1)}(x, \eta)\Psi(2^{-m}\eta))u = \sum_{k=h}^{m+\mu} a^{k-h}(x,D)u_k
\]

\[+ \sum_{\mu < l < 1 + \log_2(\Lambda/\varrho)} \text{OP}(\Psi(2^{-m}D_x)\psi(2^{h-l-m}D_x)a(x, \eta)\varphi(2^{-m-l}\eta)\Psi(2^{-m}\eta))u. \tag{5.35} \]

A similar reasoning applies to \( a^{(3)}(x, \eta) \). The main difference is that the possible inclusion of \( \text{supp}\varphi(2^{-j} \cdot) \), into the level sets where \( \Psi(2^{-m} \cdot) \) equals 1 or 0, in this case applies to the symbol \( \Psi(2^{-m}D_x)a_j(x, \eta) = \mathcal{F}^{-1}_{\xi \rightarrow x}(\Psi(2^{-m} \xi)\varphi(2^{-j} \xi)a(\xi, \eta)) \). Therefore one has for the same \( \mu \),

\[
\text{OP}(\Psi(2^{-m}D_x)a^{(3)}(x, \eta)\Psi(2^{-m}\eta))u = \sum_{j=h}^{m+\mu} a_j(x,D)u^{j-h}
\]

\[+ \sum_{\mu < l < 1 + \log_2(\Lambda/\varrho)} \text{OP}(\Psi(2^{-m}D_x)\varphi(2^{-l-m}D_x)a(x, \eta)\psi(2^{h-m-l}\eta)\Psi(2^{-m}\eta))u. \tag{5.36} \]

Treating \( a^{(2)}(x, D) \) analogously, it is not difficult to see that once again the central issue is whether \( \text{supp}\varphi(2^{-k} \cdot) \) is contained in the set where \( \Psi(2^{-m} \cdot) = 1 \) or 0. So when \( m \geq h \) for simplicity, one has for the same \( \mu \), and with primed sums over the same integers \( l \) as above,

\[
\text{OP}(\Psi(2^{-m}D_x)a^{(2)}(x, \eta)\Psi(2^{-m}\eta))u = \sum_{k=0}^{m+\mu} (a^{k-h}(x,D)u_k + a_k(x,D)(a^{k-1} - a^{k-h}))
\]

\[+ \sum_{\mu < l < 1 + \log_2(\Lambda/\varrho)} \text{OP}(\Psi(2^{-m}D_x)(\varphi^{m+1}(x, \eta) - \varphi^{m+l-h}(x, \eta))\varphi(2^{-m-l}\eta)\Psi(2^{-m}\eta))u
\]

\[+ \sum_{\mu < l < 1 + \log_2(\Lambda/\varrho)} \text{OP}(\Psi(2^{-m}D_x)a^{m+1}(x, \eta)(\psi(2^{1-m-l}\eta) - \psi(2^{h-m-l}\eta))\Psi(2^{-m}\eta))u. \tag{5.37} \]

The programme introduced after Proposition 5.6 is now completed by letting \( m \rightarrow \infty \) in (5.35)–(5.37) and observing that the infinite series in (5.8)–(5.10) reappear in this way. Of course, this relies on the fact that the remainders in the primed sums over \( l \) can be safely ignored:

**Proposition 5.7.** When \( a(x, \eta) \) is given in \( S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n) \) and \( \Psi, \varphi \in C_0^\infty(\mathbb{R}^n) \) equal 1 in neighbourhoods of the origin, then it holds for every \( u \in \mathcal{S}'(\mathbb{R}^n) \) that each term (with \( l \) fixed) in the primed sums in (5.35)–(5.36) tends to 0 in \( \mathcal{S}'(\mathbb{R}^n) \) for \( m \rightarrow \infty \).

This is valid for (5.37) too, if \( a(x, \eta) \) in addition fulfils the twisted diagonal condition (2.27).

**Proof.** To show that each remainder term tends to 0 for \( m \rightarrow \infty \) and fixed \( l \), it suffices to verify (A.1) and (A.2) in view of Remark A.2.

For \( a^{(1)}(x, D) \), note that by repeating the proof of Proposition 5.2 (ignoring \( \Psi \)) each remainder in (5.35) has \( \xi \) in its spectrum only when \( (R_02^l)^2m \leq |\xi| \leq \frac{5^l}{2}R2^m \).
Moreover, each remainder term is \( \leq c2^{k(N+d)}(1+|x|)^N \) for \( N \geq \text{order } \nabla \phi(\hat{u}) \) according to Proposition 3.6, for with the cut-off functions \( \Psi, \psi(2^{k-1} \cdot) \) and \( \varphi(2^{-l} \cdot) \Psi \) the latter is 0 around the origin. So a crude estimate by \( c2^{k(N+d+N)}(1+|x|)^N \) shows that (A.2) is fulfilled.

Similarly for the primed sum in (5.36), where \( \psi(2^{h-l} \cdot) \Psi \) is 1 around the origin; which again results in the bound \( c2^{k(N+d)}(1+|x|)^N \) for \( N \neq -d \).

The procedure also works for (5.37), for (A.1) is verified as in Proposition 5.3, cf. (5.19), because the extra spectral localisations provided by \( \Psi(2^{-m} \cdot) \) cannot increase the spectra. For the pointwise estimates one may now use e.g. \( \Psi \varphi(2^{-l} \cdot) \) and \( (\psi(2^{1-l} \cdot) - \psi(2^{h-l} \cdot)) \Psi \) as the cut-off functions in the last part of (5.37). This yields the proof of Proposition 5.7. \( \square \)

An extension of the proposition’s remainder analysis to general \( a^{(2)}(x,D) \) without a condition on the behaviour along the twisted diagonal does not seem feasible. But such results will follow in Section 6 from a much deeper investigation of \( a(x,D) \) itself; cf. Theorem 6.4.

Remark 5.8. The type 1,1-operator \( a^{(1)}(x,D) \) induced by (5.26) is a paradifferential operator in the sense of Bony [2], as well as in Hörmander’s framework of residue classes in [20, Ch. 10]. The latter follows from (5.30), but will not be pursued here. \( a^{(2)}(x,D) \) and \( a^{(3)}(x,D) \) are also called paradifferential operators, following Yamazaki [48]. The decomposition (5.7)–(5.10) can be traced back to Kumano-go and Nagase, who used a variant of \( a^{(1)}(x,\eta) \) to smooth non-regular symbols, cf. [31, Thm 1.1]. It was exploited in continuity analysis of pseudo-differential operators in e.g. [48, 33, 24, 32].

Remark 5.9. For pointwise multiplication decompositions analogous to (5.7) were used implicitly by Peetre [40], Triebel [47]; and more explicitly in the paraproducts of Bony [2]. Moreover, for \( a = a(x) \) Definition 2.1 reduces to the product \( \pi(a,u) \) introduced formally in [22] as

\[
\pi(a,u) = \lim_{m \to \infty} a^m \cdot u^m. \tag{5.38}
\]

This was analysed in [22], including continuity properties deduced from (5.7), that essentially is a splitting of the generalised pointwise product \( \pi(\cdot,\cdot) \) into paraproducts. Partial associativity, i.e. \( f \pi(a,u) = \pi(fu,a) = \pi(a,fu) \) for \( f \in C^\infty \), was first obtained with the refined methods developed later in [26, Thm. 6.7], though.

6. Action on Temperate Distributions

6.1. Littlewood–Paley Analysis of Type 1,1-Operators. First the full set of conclusions is drawn for the operators \( a^{(j)}(x,D), \ j = 1,2,3 \) studied in Section 5.3. Of course none of them have anomalies if \( a(x,\eta) \) fulfils (2.27):

**Theorem 6.1.** When \( a(x,\eta) \) is a symbol in \( S^{d}_{1,1}(\mathbb{R}^n \times \mathbb{R}^n) \) for some \( d \in \mathbb{R} \) and \( \eta \in C^\infty_0(\mathbb{R}^n) \) equals 1 around the origin, then the associated type 1,1-operators \( a^{(1)}(x,D) \) and \( a^{(3)}(x,D) \) are everywhere defined continuous linear maps

\[
a^{(1)}(x,D), a^{(3)}(x,D) : \mathcal{S}'(\mathbb{R}^n) \to \mathcal{S}'(\mathbb{R}^n), \tag{6.1}
\]
that are given by formulae (5.8) and (5.10), where the infinite series converge rapidly in \( \mathcal{S}'(\mathbb{R}^n) \) for every \( u \in \mathcal{S}'(\mathbb{R}^n) \). The adjoints are also in \( \text{OP}(S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n)) \).

If furthermore \( a(x, \eta) \) fulfills the twisted diagonal condition (2.27), these conclusions are valid verbatim for the operator \( a^{(2)}_\varphi(x,D) \) given by the series in (5.9).

**Proof.** As the symbols \( a^{(1)}(x, \eta) \) and \( a^{(3)}(x, \eta) \) both belong to \( S^d_{1,1} \) and fulfill (2.27) by Proposition 5.6, the corresponding operators are defined and continuous on \( \mathcal{S}' \) by Proposition 4.2, with \( a^{(1)}(x,D)^* \) and \( a^{(3)}(x,D)^* \) both of type 1, 1.

Since \( \text{supp} \mathcal{F}_{x \to \xi} a^{(2)}(x,D) \subset \text{supp} \mathcal{F}_{x \to \xi} a \) it follows that \( a^{(2)}(x,D) \) satisfies (2.27), when \( a(x, \eta) \) does so. Hence the preceding argument also applies to \( a^{(2)}(x,D) \), so that it is continuous on \( \mathcal{S}' \) with its adjoint being of type 1, 1.

Moreover, the series \( \sum_{k=0}^\infty a^{k,h}(x,D)u_k \) in (5.8) converges rapidly in \( \mathcal{S}' \) for every \( u \in \mathcal{S}' \). This follows from \(^1\) of Lemma A.1, for the terms fulfil (A.1) and (A.2) by Proposition 5.2, cf. (5.14), and Proposition 5.5, respectively. (The latter gives a bound by \( 2^{k(N+d_+)}(1 + |x|)^{N+d_+} \))

Now the distribution \( \sum_{k=0}^\infty a^{k,h}(x,D)u_k \) equals \( a^{(1)}(x,D)u \) because of formula (5.35), since the primed sum there goes to 0 for \( m \to \infty \), as shown in Proposition 5.7.

Similarly Lemma A.1 yields convergence of the series (5.10) for \( a^{(3)}(x,D)u \) when \( u \in \mathcal{S}' \). By Propositions 5.3 and 5.5, convergence of the \( a^{(2)} \)-series in (5.11) also follows from Lemma A.1. The series identify with the operators in view of the remark made prior to Proposition 5.7.

It should be emphasized that duality methods and pointwise estimates contribute in two different ways in Theorem 6.1: once the symbol \( a^{(1)}(x, \eta) \) has been introduced, continuity on \( \mathcal{S}'(\mathbb{R}^n) \) of the associated type 1, 1-operator \( a^{(1)}(x,D) \) is obtained by duality through Proposition 4.2. However, the pointwise estimates in Section 3 yield (vanishing of the remainder terms, hence) the identification of \( a^{(1)}(x,D)u \) with the series in (5.8). Furthermore, the pointwise estimates also give an explicit proof of the fact that \( a^{(1)}(x,D) \) is defined on the entire \( \mathcal{S}'(\mathbb{R}^n) \), for the right-hand side of (5.8) does not depend on the modulation function \( \Psi \). Similar remarks apply to \( a^{(3)}(x,D) \). Thus duality methods and pointwise estimates together lead to a deeper analysis of type 1, 1-operators.

**Remark 6.2.** Theorem 6.1 and its proof generalise a result of Coifman and Meyer [37, Ch. 15] in three ways. They stated Lemma A.1 for \( \theta_0 = \theta_1 = 1 \) and derived a corresponding fact for pararmultiplication, though only with a treatment of the first and third term.

Going back to the given \( a(x,D) \), one derives from Theorem 6.1 and (5.7) the following

**Theorem 6.3.** When \( a \in S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n) \) fulfills the twisted diagonal condition (2.27), then the associated type 1, 1-operator \( a(x,D) \) defined by vanishing frequency modulation is an everywhere defined continuous linear map

\[
a(x,D) : \mathcal{S}'(\mathbb{R}^n) \rightarrow \mathcal{S}'(\mathbb{R}^n),
\]

with its adjoint \( a(x,D)^* \) also in \( \text{OP}(S^d_{1,1}(\mathbb{R}^n \times \mathbb{R}^n)) \). The operator fulfills

\[
a(x,D)u = a^{(1)}(x,D)u + a^{(2)}(x,D)u + a^{(3)}(x,D)u
\]
for every \( \Psi \in C^\infty_c(\mathbb{R}^n) \) equal to 1 in a neighbourhood of the origin, and the series in (5.8)–(5.10) converge rapidly in \( \mathcal{S}'(\mathbb{R}^n) \) for every \( u \in \mathcal{S}'(\mathbb{R}^n) \).

To extend the discussion to general \( a(x,D) \) without vanishing along the twisted diagonal, note that Theorem 6.1 at least shows that \( a^{(1)}(x,D)u \) and \( a^{(3)}(x,D)u \) are always defined and that (5.8) and (5.10) are operator identities.

It remains to justify the operator notation \( a^{(2)}(x,D) \) in (5.9) and to give its precise relation to \( a(x,D) \) itself. The point of departure is of course the symbol splitting (5.29); the corresponding type 1,1-operators are still denoted by \( a^{(j)}(x,D) \). However, to avoid ambiguity the series in (5.8)–(5.10) will now be temporarily written as \( A^{(j)}_\Psi u \), whence (5.7) amounts to

\[
a(x,D)u = A^{(1)}_\Psi u + A^{(2)}_\Psi u + A^{(3)}_\Psi u \quad \text{for } u \in \mathcal{S}'.
\]  

Here the left-hand side exists if and only if the series \( A^{(2)}_\Psi u \) converges, as \( A^{(1)}_\Psi u, A^{(3)}_\Psi u \) always converge by Theorem 6.1. This strongly indicates that (6.5) below is true. In fact, this main result of the analysis is obtained by frequency modulation:

**Theorem 6.4.** When \( a(x, \eta) \in \mathcal{S}'_{1,1}(\mathbb{R}^n \times \mathbb{R}^n) \) and \( a^{(2)}(x, \eta) \) denotes the type 1,1-symbol in (5.28), derived from the paradifferential decomposition (5.7), then

\[
D(a^{(2)}(x,D)) = D(a(x,D))
\]

and \( u \in \mathcal{S}'(\mathbb{R}^n) \) belongs to these domains if and only if the series (5.9), or equivalently (5.11), converges in \( \mathcal{S}'(\mathbb{R}^n) \) — in which case (also) formulae (5.9), (5.11) are operator identities.

**Proof.** A variant of (6.4) follows at once from (5.29), using a second modulation function \( \Psi \) and the brief notation from Remark 5.1,

\[
d^m(x,D)u^m = \sum_{l=1,2,3} \text{OP}(\Psi(2^{-m}D_x)a^{(l)}(x, \eta)\Psi(2^{-m}\eta))u.
\]

Here the terms with \( l = 1 \) and \( l = 3 \) always have \( \Psi \)-independent limits for \( m \to \infty \) according to Theorem 6.1, so it is clear from the definition by vanishing frequency modulation that \( u \in D^{(2)}(a(x,D)) \) is equivalent to \( u \in D(a(x,D)) \), hence to convergence of \( A^{(2)}_\Psi u \), cf. (6.4) ff.

As for the last claim, whenever \( u \in D(a(x,D)) \), then passage to the limit \( (m \to \infty) \) in the above equation yields the following, when (5.35), (5.37) and (5.36) are applied, now with the remainders in the primed sums there denoted by \( R^{(1)}_m u, R^{(2)}_m u, R^{(3)}_m u \) for brevity:

\[
a(x,D)u = \lim_{m \to \infty} \sum_{l=1,2,3} \text{OP}(\Psi(2^{-m}D_x)a^{(l)}(x, \eta)\Psi(2^{-m}\eta))u
\]

\[
= A^{(1)}_\Psi u + A^{(2)}_\Psi u + A^{(3)}_\Psi u + 0 + \lim_{m \to \infty} R^{(2)}_m u + 0.
\]

Note that convergence of \( R^{(2)}_m u \) follows from that of the other six terms; cf. Proposition 5.7. Compared to (6.4) this yields \( \lim_{m \to \infty} R^{(2)}_m u = 0 \), which via (5.37) gives that \( a^{(2)}(x,D)u = A^{(2)}_\Psi u \).
6.2. The Twisted Diagonal Condition of Arbitrary Order. When \( a(x, \eta) \) is in the self-adjoint subclass \( \tilde{S}^d_{1,1} \), then it follows Theorem 4.6 that the domains in (6.5) equal \( \mathcal{S}' \).

However, it is interesting to give an explicit proof that the domains in (6.5) equal \( \mathcal{S}' \) whenever \( a \in \tilde{S}^d_{1,1} \). This can be done in a natural way by extending the proof of Theorem 6.4, where the special estimates in (2.33) enter the convergence proof for \( a^{(2)}(x, D)u \) directly, because they are rather close to the symbol factors from the factorisation inequalities in Section 3. The full generality with \( \theta_0 < \theta_1 \) in the corona criterion Lemma A.1 is also needed now.

**Theorem 6.5.** Suppose \( a(x, \eta) \in \tilde{S}^d_{1,1} (\mathbb{R}^n \times \mathbb{R}^n) \), i.e. \( a(x, \eta) \) fulfils one of the equivalent conditions in Theorem 4.3. Then the conclusions of Theorems 6.3–6.4 remain valid for \( a(x, D); \) in particular \( D(a(x, D)) = \mathcal{S}'(\mathbb{R}^n) \).

**Proof.** The continuity on \( \mathcal{S}' \) is assured by Theorem 4.6. For the convergence of the series in the paradifferential splitting, it is convenient to write, in the notation of (2.33) ff,

\[
a(x, \eta) = (a(x, \eta) - a_{\chi, 1}(x, \eta)) + a_{\chi, 1}(x, \eta),
\]

where \( a-a_{\chi, 1} \) satisfies (2.27) for \( B = 1 \), so that Theorem 6.3 applies to it. As \( a_{\chi, 1} \) is in \( \tilde{S}^d_{1,1} \) like \( a \) and \( a-a_{\chi, 1} \) (the latter by Proposition 4.2), one may reduce to the case in which

\[
\hat{a}(x, \eta) \neq 0 \implies \max(1, |\xi + \eta|) \leq |\eta|. \tag{6.9}
\]

Continuing under this assumption, it is according to Theorems 6.1 and 6.4 enough to show for all \( u \in \mathcal{S}' \) that there is convergence of the two contributions to \( a^{(2)}(x, D)u \),

\[
\sum_{k=0}^{\infty} (a^k - a^{k-h})(x, D)u_k, \quad \sum_{k=1}^{\infty} a_k(x, D)(u^{k-1} - u^{k-h}). \tag{6.10}
\]

Since the terms here are functions of polynomial growth by Proposition 5.5, it suffices to improve the estimates there; and to do so for \( k \geq h \).

Using Hörmander’s localisation to a neighbourhood of \( \mathcal{S} \), cf. (2.30)–(2.32), one arrives at

\[
\hat{a}_{\chi, \varepsilon}(\xi, \eta) = \hat{a}(\xi, \eta)\varphi(2^{-k}\xi)\chi(\xi + \eta, \varepsilon \eta), \tag{6.11}
\]

This leaves the remainder \( b_k(x, \eta) = a_k(x, \eta) - a_{\chi, \varepsilon}(x, \eta) \), that applied to the above difference \( v_k = u^{k-1} - u^{k-h} = \mathcal{F}^{-1}((\varphi(2^{1-k}) - \varphi(2^{h-k}))\hat{u}) \) gives

\[
a_k(x, D)v_k = a_{\chi, \varepsilon}(x, D)v_k + b_k(x, D)v_k. \tag{6.12}
\]

To utilise the pointwise estimates, fix \( N \geq \text{ord}_{\mathcal{S}'}(\hat{u}) \) so that \( d + N \neq 0 \); and pick \( \Psi \in C^\infty_0(\mathbb{R}^n) \) equal to 1 in a neighbourhood of the corona \( 2^{-1-h} \leq |\eta| \leq 1 \) and equal to 0 outside the set with \( 2^{-2-h} \leq |\eta| \leq 2 \). Taking the dilated function \( \Psi(\eta/(R2^k)) \) as the auxiliary function in the symbol factor, the factorisation inequality (3.1) and Theorem 3.2 give

\[
|a_{\chi, \varepsilon}(x, D)v_k(x)|F_{a_{\chi, \varepsilon}}(N, R2^k; x)v_k^*(N, R2^k; x) \tag{6.13}
\]
which is estimated from above by
\[
.cv_k^*(x) \sum_{|a|=0}^{N+|a|/2+1} (\int_{R^{k-1}} |(R_{q}^{k})| d^2 \eta^{a}\phi_{k}(x, \eta)^2 d\eta)^{1/2}.
\]

(6.14)

Here the ratio of the limits is \(2R/(r_2^{-h-2}) > 32\), so with extension to \(R_2^{k+1-L} \leq |\eta| \leq R_2^{k+1}\), there is \(L \geq 6\) dyadic coronas. This gives an estimate by \(c(R_{q}^{k})^d L^{1/2} N_{X,e,a}(a_k)\). In addition, Minkowski’s inequality gives
\[
N_{X,e,a}(a_k) \leq \sup_{\rho > 0} |a|^d \int_{\mathbb{R}^n} (2^{kn} \phi(2^{k} |x|))(\int_{\rho \leq |\eta| \leq 2\rho} |D_{\eta} a(x, \eta)^2 d\eta|)^{1/2} dy
\]

(6.15)

so it follows from the above that
\[
|a_k(x, \eta)| \leq cN_{X,e,a}(a_k) (\sum_{|a| \leq N+|a|/2+1} c_{\alpha, \sigma}) (\sum_{|a| \leq N+|a|/2+1} c_{\alpha, \sigma}) L^{1/2} (R_{q}^{k})^d.
\]

(6.16)

Using Lemma 3.1 and taking \(\varepsilon = 2^{-k\theta}\), say for \(\theta = 1/2\) this gives
\[
|a_k(x, \eta)| \leq c(1 + |x|)^N 2^{-k(\sigma_1-2d-3N/2)}.
\]

(6.17)

Choosing \(\sigma > 3N + 2d + 1\), the series \(\sum b_k(x, \eta)\) converges rapidly for \(\phi \in \mathcal{F}\).

To treat \(\sum b_k(x, \eta)\) it is observed that \(\hat{a}_{k, e}(\xi, \eta) = \hat{a}(\xi, \eta)\) holds in the set where \(\chi(\xi, \eta, 2^{-k\theta} \eta) = 1\), that is, when \(2 \max(1, |\xi + \eta|) \leq 2^{-k\theta} |\eta|\), so by (6.9),
\[
\sup \hat{b} \subset \{ (\xi, \eta) \mid 2^{-1-k\theta} |\eta| \leq \max(1, |\xi + \eta|) \leq |\eta| \}
\]

(6.18)

This implies by Theorem B.1 that \(\xi = \xi + \eta\) is in \(\sup \mathcal{F} b_k(x, \eta)\) only if both
\[
|\xi| \leq |\eta| \leq R_2^k
\]

(6.19)
\[
\max(1, |\xi|) \geq 2^{-1-k\theta} |\eta| \geq r_2^{k(1-\theta)-h-2}.
\]

(6.20)

When \(k\) fulfils \(2^{k(1-\theta)} > 2^{h+2}/r\), so that the last right-hand side is \(> 1\), these inequalities give
\[
(r_2^{-h-2}) 2^{k(1-\theta)} \leq |\xi| \leq R_2^k.
\]

(6.21)

This shows that the corona condition (A.1) in Lemma A.1 is fulfilled for \(\theta_0 = 1 - \theta = 1/2\) and \(\theta_1 = 1\), and the growth condition (A.2) is easily checked since both \(a_k(x, \eta)\) and \(a_k(x, \eta)\) are estimated by \(2^k(1+|x|)^N\), as can be seen from (6.17) and Proposition 3.6, respectively. Hence \(\sum b_k(x, \eta)\) converges rapidly.

For the series \(\sum b_k(x, \eta)\) it is not complicated to modify the above. Indeed, the pointwise estimates of \(v_k^*\) are easily carried over to \(u_k^*\), for \(R_2^k\) was used as the outer spectral radius of \(v_k^*\); and \(R_2^{k-1}\) may also be used as the inner spectral radius of \(u_k^*\). In addition the symbol \(a_k - a_{k-h}^*\) can be treated by replacing \(\phi(2^{-k\xi})\) by \(\psi(2^{-k\xi}) - \psi(2^{h-k} \xi)\) in
(6.11) ff., for the use of Minkowski’s inequality will now give the factor \( \int |\psi - \psi(2^h \cdot)| \, dy \) in the constant. For the remainder

\[
\tilde{b}_k(x, D) u_k = (a_k - a_k^{-h})(x, D) u_k - (a_k - a_k^{-h}) \chi(x, D) u_k
\]

one can apply the treatment of \( b_k(x, D) v_k \) verbatim.

\[\Box\]

Remark 6.6. The analysis in Theorem 6.5 is also exploited in the \( L_\varphi \)-theory of type 1, 1-operators in \([29]\). Indeed, the main ideas of the above proof was used in \([29, \text{Sect. 5.3}]\) to derive certain continuity results in the Lizorkin–Triebel scale \( F_{p,q}^s(\mathbb{R}^n) \) for \( p < 1 \), which (except for a small loss of smoothness) generalise results of Hounie and dos Santos Kapp \([21]\).

7. Final Remarks

In view of the satisfying results on type 1, 1-operators in \( \mathscr{S}'(\mathbb{R}^n) \), cf. Section 6, and the continuity results in the scales \( H_p^s, C_p^s, F_{p,q}^s \) and \( B_{p,q}^s \) presented in \([29]\), their somewhat unusual definition by vanishing frequency modulation in Definition 2.1 should be well motivated.

As an open problem, it remains to characterise the type 1, 1-operators \( a(x, D) \) that are everywhere defined and continuous on \( \mathscr{S}'(\mathbb{R}^n) \). For this it was shown above to be sufficient that \( a(x, \eta) \) is in \( \tilde{S}_{1,1}^d(\mathbb{R}^n \times \mathbb{R}^n) \), and it could of course be conjectured that this is necessary as well.

Similarly, since the works of Bourdaud and Hörmander, cf. \([6, \text{Ch. IV}], [7], [18, 19] \) and also \([20]\), it has remained an open problem to determine

\[
\mathbb{B}(L_2(\mathbb{R}^n)) \cap \text{OP}(S^{0}_{1,1}) \quad (7.1)
\]

Indeed, this set was shown by Bourdaud to contain the self-adjoint subclass \( \text{OP}(\tilde{S}^{0}_{1,1}) \), and this sufficient condition has led some authors to a few unfortunate statements, for example that lack of \( L_2 \)-boundedness for \( \text{OP}(\tilde{S}^{0}_{1,1}) \) is “attributable to the lack of self adjointness”. But self-adjointness is not necessary, since already Bourdaud, by modification of Ching’s operator \((2.13)\), gave an example \([7, \text{p. 1069}]\) of an operator \( \sigma(x, D) \) in \( \mathbb{B}(L_2) \cap \text{OP}(S^{0}_{1,1} \setminus \tilde{S}^{0}_{1,1}) \); that is, this \( \sigma(x, D)^* \) is not of type 1, 1.

However, it could be observed that \( N_{\chi,\varepsilon,\alpha}(a) = \mathcal{O}(\varepsilon^{n/2-|\alpha|}) \) by Lemma 2.13 is valid for Ching’s symbol \( a_\theta \) and that this estimate is sharp for the \( L_2 \)-unbounded version of \( a_\theta(x, D) \), by the last part of Example 2.14. Therefore, the condition

\[
N_{\chi,\varepsilon,\alpha}(a) = \mathcal{O}(\varepsilon^{n/2-|\alpha|}) \quad \text{for} \ \varepsilon \to 0 \quad (7.2)
\]

is conjectured to be necessary for \( L_2 \)-continuity of a given \( a(x, D) \) in \( \text{OP}(S^{0}_{1,1}) \).

Appendix A. Dyadic Corona Criteria

Convergence of a series \( \sum_{j=0}^{\infty} u_j \) of temperate distributions follows if the \( u_j \) both fulfil a growth condition and have their spectra in suitable dyadic coronas. This is a special case of Lemma A.1, which for \( \theta_0 = \theta_1 = 1 \) was given by Coifman and Meyer \([37, \text{Ch. 15}]\) without arguments.

Extending the proof given in \([30]\), the refined version in Lemma A.1 allows the inner and outer radii of the spectra to grow at different exponential rates \( \theta_0 < \theta_1 \), even though the number of overlapping spectra increases with \( j \). This is crucial for Theorem 6.5, so a full proof is given.
Lemma A.1. 1° Let \((u_j)_{j \in \mathbb{N}_0}\) be a sequence in \(\mathcal{S}'(\mathbb{R}^n)\) fulfilling that there exist \(A > 1\) and \(\theta_1 \geq \theta_0 > 0\) such that \(\text{supp } \hat{u}_0 \subset \{ \xi \mid |\xi| \leq A \}\) while for \(j \geq 1\)

\[
\text{supp } \hat{u}_j \subset \{ \xi \mid \frac{1}{2}2^{j\theta_0} \leq |\xi| \leq 2A^{j\theta_1} \},
\]  
(A.1)

and that for suitable constants \(C \geq 0, N \geq 0\),

\[
|u_j(x)| \leq C2^{jN\theta_1}(1 + |x|)^N \text{ for all } j \geq 0.
\]  
(A.2)

Then \(\sum_{j=0}^{\infty} u_j\) converges rapidly in \(\mathcal{S}'(\mathbb{R}^n)\) to a distribution \(u\), for which \(\hat{u}\) is of order \(N\).

2° For every \(u \in \mathcal{S}'(\mathbb{R}^n)\) both (A.1) and (A.2) are fulfilled for \(\theta_0 = \theta_1 = 1\) by the functions \(u_0 = \Phi_0(D)u\) and \(u_j = \Phi(2^{-j}D)u\) when \(\Phi_0, \Phi \in C_0^\infty(\mathbb{R}^n)\) and \(0 \notin \text{supp } \Phi\). In particular this is the case for a Littlewood–Paley decomposition \(1 = \Phi_0 + \sum_{j=1}^{\infty} \Phi(2^{-j}\cdot)\).

Proof. In 2° it is clear that \(\Phi\) is supported in a corona, say \(\{ \xi \mid \frac{1}{A} \leq |\xi| \leq A \}\) for a large \(A > 0\); hence (A.1), (A.2) follows from the proof of Lemma 3.1.

The proof of 1° exploits a well-known construction of an auxiliary function: taking \(\psi \in C_0^\infty(\mathbb{R}^n)\) depending on \(|\xi|\) alone and so that \(0 \leq \psi_0 \leq 1\) with \(\psi_0(\xi) = 1\) for \(|\xi| \leq 1/(2A)\) while \(\psi_0(\xi) = 0\) for \(|\xi| \geq 1/A\), then

\[
\frac{d}{dt}\psi_0(\frac{\xi}{t}) = \psi(\frac{\xi}{t})\frac{1}{t} \quad \text{for} \quad \psi(\xi) = -\xi \cdot \nabla \psi_0(\xi),
\]  
(A.3)

which by integration for \(1 \leq t \leq \infty\) gives an uncountable partition of unity

\[
1 = \psi_0(\xi) + \int_{1}^{\infty} \psi(\frac{\xi}{t}) \frac{dt}{t}, \quad \xi \in \mathbb{R}^n.
\]  
(A.4)

Clearly the support of \(\psi(\xi/t)\) is compact and given by \(A|\xi| \leq t \leq 2A|\xi|\) when \(\xi\) is fixed. For \(j \geq 1\) this implies

\[
\hat{u}_j = \hat{u}_j \psi_0 + \hat{u}_j \int_{1}^{\infty} \psi(\frac{\xi}{t}) \frac{dt}{t} = \hat{u}_j \int_{2^{j\theta_0}}^{2^{(j+1)\theta_0}} \psi(\frac{\xi}{t}) \frac{dt}{t}.
\]  
(A.5)

Defining \(\psi_j \in C_0^\infty(\mathbb{R}^n)\) as the last integral here, \(\psi_j = 1\) on \(\text{supp } \hat{u}_j\); so if \(\varphi \in \mathcal{S}\),

\[
|\langle u_j, \varphi \rangle| \leq \left| \langle 1 + (1 + |x|^2)^{-\frac{N+n}{2}} u_j \rangle \right| \left| \left( 1 + |x|^2 \right)^{-\frac{N+n}{2}} \mathcal{F}^{-1}(\psi_j \varphi) \right|_2.
\]  
(A.6)

The first norm is \(O(2^{N\theta_1})\) by (A.2). For the second, note that

\[
\text{supp } \psi_j \subset \{ \xi \in \mathbb{R}^n \mid A^{-1}2^{j\theta_0} \leq |\xi| \leq 2A^{j\theta_1} \}.
\]  
(A.7)

and \(\|D^\alpha \psi_j\|_\infty \leq 2^{-j\theta_0} |\alpha| \|D^\alpha \psi_0\|_\infty / |\alpha|\) for \(\alpha \neq 0\) while \(\|\psi_j\|_\infty \leq \text{diam}(\psi_0(\mathbb{R}^n)) \leq 1\) by (A.3). In addition the identity \((1 + |x|^2)^{N+n} \mathcal{F}^{-1} = \mathcal{F}^{-1}(1 - \Delta)^{N+n}\) gives for arbitrary \(k > 0\),

\[
\left| (1 + |x|^2)^{N+n} \mathcal{F}^{-1}(\psi_j \varphi) \right|_2 \leq \sum_{|\alpha|, |\beta| \leq N+n} c_{\alpha, \beta} \|D^\alpha \psi_j\|_\infty \| (1 + |\xi|)^{k+n/2} D^\beta \varphi \|_\infty \left( \int_{2^{j\theta_0} \cdot / A}^\infty r^{-1-2k} dr \right)^{1/2}.
\]  
(A.8)

Here \(\|D^\alpha \psi_j\|_\infty = O(1)\), so because of the \(L_2\)-norm the above is \(O(2^{-j\theta_0})\) for every \(k > 0\).
Hence $\langle u_j, \varphi \rangle = \mathcal{O}(2^j(\theta_1N - \theta_0k))$, so $k > N\theta_1/\theta_0$ yields that $\sum_{j=0}^{\infty} \langle u_j, \varphi \rangle$ converges. \hfill \Box

**Remark A.2.** The above proof yields that the conjunction of (A.1) and (A.2) implies $\langle u_j, \varphi \rangle = \mathcal{O}(2^{jN})$ for all $N > 0$; hence there is rapid convergence of $u = \sum_{j=0}^{\infty} u_j$ in $\mathcal{S}'$ in the sense that

$$\langle u - \sum_{j<k} u_j, \varphi \rangle = \sum_{j \geq k} \langle u_j, \varphi \rangle = \mathcal{O}(2^{-kN})$$

for $N > 0$, $\varphi \in \mathcal{S}'$.

---

**APPENDIX B. THE SPECTRAL SUPPORT RULE**

To control the spectrum of $x \mapsto a(x,D)u$, i.e. the support of $\xi \mapsto \mathcal{F}a(x,D)u$, there is a simple rule which is recalled here for the reader’s convenience.

Writing $\mathcal{F}a(x,D)\mathcal{F}^{-1}(\hat{u})$ instead, the question is clearly how the support of $\mathcal{F}u$ is changed by the conjugated operator $\mathcal{F}a(x,D)\mathcal{F}^{-1}$. In terms of its distribution kernel $\mathcal{K}(\xi, \eta)$, cf. (1.15), one should expect the spectrum of $a(x,D)u$ to be contained in

$$\Xi := \text{supp} \mathcal{K} \circ \text{supp} \mathcal{F}u = \{ \xi \in \mathbb{R}^n \mid \exists \eta \in \text{supp} \hat{u}: (\xi, \eta) \in \text{supp} \mathcal{K} \}.$$  \hfill (B.1)

For $\text{supp} \mathcal{F}u \subset \mathbb{R}^n$ this was proved in [24]; but in general the closure $\Xi$ should be used instead:

**Theorem B.1.** Let $a \in S_{1,1}^\infty(\mathbb{R}^n \times \mathbb{R}^n)$ and suppose $u \in D(a(x,D))$ has the property that (2.5) holds in the topology of $\mathcal{S}'(\mathbb{R}^n)$ for some $\psi \in C_0^\infty(\mathbb{R}^n)$ equal to 1 around the origin. Then

$$\text{supp} \mathcal{F}(a(x,D)u) \subset \Xi,$$

$$\Xi = \{ \xi + \eta \mid (\xi, \eta) \in \text{supp} \mathcal{F}_{x \mapsto \xi} a, \ \eta \in \text{supp} \mathcal{F}u \}.$$ \hfill (B.2)

When $u \in \mathcal{F}^{-1} \mathcal{S}'(\mathbb{R}^n)$ the $\mathcal{S}'$-convergence holds automatically and $\Xi$ is closed for such $u$.

The reader is referred to [26] for the deduction of this from the kernel formula. Note that whilst (1.15) yields (B.3), it suffices for (B.2) to take any $v \in C_0^\infty(\mathbb{R}^n)$ with support disjoint from $\Xi$ and verify that

$$\langle \mathcal{F}a(x,D)\mathcal{F}^{-1}\hat{u}, v \rangle = \langle \mathcal{K}, v \otimes \hat{u} \rangle = 0.$$ \hfill (B.4)

Here the middle expression makes sense as $\langle (v \otimes \hat{u}) \mathcal{K}, 1 \rangle$, as noted in [26], using the remarks to [17, Def. 3.1.1]. However, the first equality sign is in general not trivial to justify: the limit in Definition 2.1 is decisive for this.

**Remark B.2.** There is a simple proof of (B.2) in the main case that $\hat{u} \in \mathcal{S}'$: If $a \in S_{1,0}^d$ and $v$ is as above, then (B.1) yields $\text{dist}(\text{supp} \mathcal{K}, \text{supp} (v \otimes \hat{u})) > 0$ since $\text{supp} \hat{u} \subset \mathbb{R}^n$. So with $\hat{u}_\varepsilon = \varphi_\varepsilon \ast \hat{u}$ for some $\varphi \in C_0^\infty(\mathbb{R}^n)$ with $\varphi(0) = 1$, $\varphi_\varepsilon = \varepsilon^{-n} \varphi(\cdot/\varepsilon)$, all sufficiently small $\varepsilon > 0$ give

$$\text{supp} \mathcal{K} \cap \text{supp} v \otimes \hat{u}_\varepsilon = \emptyset.$$ \hfill (B.5)

Therefore one has, since $\hat{u}_\varepsilon \in C_0^\infty(\mathbb{R}^n)$,

$$\langle \mathcal{F}a(x,D)\mathcal{F}^{-1}\hat{u}, v \rangle = \lim_{\varepsilon \to 0} \langle \mathcal{F}a(x,D)\mathcal{F}^{-1}\hat{u}_\varepsilon, v \rangle = \lim_{\varepsilon \to 0} \langle \mathcal{K}, v \otimes \hat{u}_\varepsilon \rangle = 0.$$ \hfill (B.6)
For general $b(x, \eta)$ in $S^0_{1,1}$ one may set $a(x, \eta) = b(x, \eta)\chi(\eta)$ for a $\chi \in C_0^\infty$ equal to 1 on an open ball containing $\text{supp} \, u$. Then $a$ is in $S^{-\infty}$ with associated kernel $\mathcal{H}_a(\xi, \eta) = \mathcal{H}_b(\xi, \eta)\chi(\eta)$ because of (1.15). Moreover, the set $\Xi$ is unchanged by this replacement, so (B.6) gives

$$\langle \mathcal{F} b(x, D) \mathcal{F}^{-1} u, v \rangle = \langle \mathcal{F} a(x, D) \mathcal{F}^{-1} u, v \rangle = \lim_{\varepsilon \to 0} \langle \mathcal{H}_a, v \otimes \hat{u}_\varepsilon \rangle = 0.$$  \hfill (B.7)

The argument in Remark B.2 clearly covers the applications of Theorem B.1 in this paper.

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