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# TRACES OF FUNCTIONS WITH A DOMINATING MIXED DERIVATIVE IN $\mathbb{R}^{3}$ 

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Abstract. We investigate traces of functions, belonging to a class of functions with dominating mixed smoothness in $\mathbb{R}^{3}$, with respect to planes in oblique position. In comparison with the classical theory for isotropic spaces a few new phenomenona occur. We shall present two different approaches. One is based on the use of the Fourier transform and restricted to $p=2$. The other one is applicable in the general case of Besov-Lizorkin-Triebel spaces and based on atomic decompositions.

Keywords: Sobolev spaces of dominating mixed smoothness, Besov and Lizorkin-Triebel classes of dominating mixed smoothness, Fourier analytic characterizations, atomic decompositions, traces on hyperplanes in oblique position

MSC 2000: 42B35, 46E35

## 1. Introduction

Sobolev spaces with dominating mixed smoothness $S_{p}^{\bar{r}} W\left(\mathbb{R}^{d}\right)$ were introduced in 1962 by S. M. Nikol'skij, see [13], [14], originally in connection with some partial differential equations. Later on there has been some interest in these types of spaces as special cases of vector-valued Sobolev spaces $\left(S_{p}^{r, \ldots, r} W\left(\mathbb{R}^{d}\right)\right.$ can be interpreted as an iterated version of the Sobolev spaces $W_{p}^{r}(\mathbb{R})$ ), see Grisvard [9], Sparr [19] and Schmeißer [17]. At the end of the eighties Triebel [20], motivated by problems in connection with eigenvalue distributions of integral operators, investigated the trace problem with respect to the diagonal $x_{1}=x_{2}$ for the Besov spaces $S_{p, 1}^{r, r} B\left(\mathbb{R}^{2}\right)$. In recent years there is an increasing interest in function spaces with a dominating mixed derivative in connection with the numerical solution of some special partial differential equations or integral equations, see e.g. Griebel, Oswald, Schiekofer [8], Yserentant [26], [27], Nitzsche [15] or Bungartz and Griebel [4].

We are interested in the description of the trace classes of $S_{p}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right)$ (and more general function spaces) with respect to a hyperplane in oblique position. For at least twenty years the situation is well understood if the trace is taken with respect to hyperplanes parallel to the coordinate axes, cf. e.g. the monographs Amanov [2], Gelman, Maz'ya [7] ( $p=2$ ) and Schmeißer, Triebel [18]. However, there is an essential difference if the hyperplane is in an oblique position. First observations in this direction have been made by Triebel [20] in the two-dimensional case, later continued by Rodriguez [16] and complemented by the first author, see [25]. To our own surprise the problem for $d=3$ turned out to be more complicated. New phenomenona occur. Whereas for $d=2$ almost all trace classes of Sobolev and Besov-Lizorkin-Triebel spaces are again Besov or Lizorkin-Triebel classes (in some limiting cases of generalized smoothness, see [25]) the situation changes with $d=3$. Here it turns out that the trace classes can be described as the sum of three different function spaces of dominating mixed smoothness. In proving such a statement we offer two different approaches. The first one is restricted to $p=2$ and uses elementary properties of the Fourier transform. In this simplified situation we are also able to establish a characterization of the trace class of $S_{2}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right)$ as an $L_{2}$-space with a weight in the Fourier image. For $p \neq 2$ one is confronted with difficult Fourier multiplier assertions. To circumvent this we apply the characterization of these classes by atoms which works also in the more general case of Besov and LizorkinTriebel spaces. However, the description of the trace classes found in this way is not very transparent. Here some further progress would be desirable.

To explain a part of the difficulties let us consider an example. We equip the hyperplane $\Gamma$ given by $x_{1}+x_{2}+x_{3}=0$ with an orthogonal basis

$$
\begin{equation*}
\mathscr{O}=\left\{\vec{\sigma}_{1}, \vec{\sigma}_{2}\right\}, \vec{\sigma}_{1}=\left(\sigma_{1,1}, \sigma_{1,2}, \sigma_{1,3}\right) \in \Gamma, \vec{\sigma}_{2}=\left(\sigma_{2,1}, \sigma_{2,2}, \sigma_{2,3}\right) \in \Gamma, \vec{\sigma}_{1} \perp \vec{\sigma}_{2} . \tag{1.1}
\end{equation*}
$$

Then we associate to this basis the corresponding "orthogonal" trace operator

$$
\begin{equation*}
\left(\operatorname{tr}_{\mathscr{O}} f\right)\left(z_{1}, z_{2}\right)=f\left(z_{1} \vec{\sigma}_{1}+z_{2} \vec{\sigma}_{2}\right), \quad z_{1}, z_{2} \in \mathbb{R} . \tag{1.2}
\end{equation*}
$$

Now we consider the following family of functions

$$
f_{\lambda}\left(x_{1}, x_{2}, x_{3}\right):=\psi\left(x_{1}\right) \psi\left(x_{2}\right) \psi\left(x_{3}\right)\left|x_{3}\right|^{\lambda}, \quad\left(x_{1}, x_{2}, x_{3}\right) \in \mathbb{R}^{3}, \quad \lambda \in \mathbb{R}
$$

where $\psi: \mathbb{R} \rightarrow \mathbb{R}$ is a smooth cut-off function supported near the origin. Such a function $f_{\lambda}$ belongs to $S_{p}^{r, r, r} W\left(\mathbb{R}^{3}\right)$ if $\lambda>r-1 / p$. But the regularity of the function
$g_{\lambda}\left(z_{1}, z_{2}\right)=\psi\left(\sigma_{1,1} z_{1}+\sigma_{2,1} z_{2}\right) \psi\left(\sigma_{1,2} z_{1}+\sigma_{2,2} z_{2}\right) \psi\left(\sigma_{1,3} z_{1}+\sigma_{2,3} z_{2}\right)\left|\sigma_{1,3} z_{1}+\sigma_{2,3} z_{2}\right|^{\lambda}$
depends on $\mathscr{O}$. The function $g_{\lambda}$ belongs to $S_{p}^{r, r} W\left(\mathbb{R}^{2}\right), \lambda>r-1 / p$, if either $\sigma_{1,3}=0$ or $\sigma_{2,3}=0$. If $\sigma_{1,3} \cdot \sigma_{2,3} \neq 0$, then $g_{\lambda}$ belongs to $S_{p}^{t, t} W\left(\mathbb{R}^{2}\right), \lambda>2 t-1 / p$. As a consequence the description of the traces of $S_{p}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right)$ to the hyperplane $x_{1}+x_{2}+x_{3}=0$ must depend on the chosen basis $\mathscr{O}$.

The paper is organized as follows. In Section 2 we start with a general discussion of the notion of the trace and continue with a detailed investigation of the trace problem for the Sobolev spaces of dominating smoothness built on $L_{2}\left(\mathbb{R}^{3}\right)$. Here we shall apply methods from Fourier analysis. In the case of $p \neq 2$, treated in Section 3, the situation becomes more complicated and we switch to the powerful but less transparent method of decompositions of functions into small building blocks like atoms. By means of those decompositions we are able to describe the trace classes for the general case of Besov and Lizorkin-Triebel classes. Our main results are contained in Theorems 2.11, 3.10, and 3.14.
2. The trace of Sobolev spaces of dominating mixed smoothness $S_{2}^{\bar{r}}\left(\mathbb{R}^{3}\right)$

### 2.1. Sobolev spaces of dominating mixed smoothness.

Let $1<p<\infty$ and $\bar{r}=\left(r_{1}, \ldots, r_{d}\right) \in \mathbb{N}_{0}^{d}\left(\mathbb{N}_{0}\right.$ denotes the natural numbers including 0 ). The Sobolev space of dominating mixed smoothness $\bar{r}=\left(r_{1}, \ldots, r_{d}\right)$ is the collection of all functions $f \in L_{p}\left(\mathbb{R}^{d}\right)$ such that

$$
D^{\alpha} f \in L_{p}\left(\mathbb{R}^{d}\right), \quad 0 \leqslant \alpha_{i} \leqslant r_{i}, \quad i=1, \ldots, d,
$$

endowed with the norm

$$
\begin{equation*}
\left\|f\left|S_{p}^{\bar{r}} W\left(\mathbb{R}^{d}\right)\left\|:=\sum_{\alpha \leqslant \bar{r}}\right\| D^{\alpha} f\right| L_{p}\left(\mathbb{R}^{d}\right)\right\| . \tag{2.1}
\end{equation*}
$$

Here $\alpha \leqslant \bar{r}$ means $\alpha_{i} \leqslant r_{i}, i=1, \ldots, d$.
The mixed derivative $\partial^{r_{1}+\ldots+r_{d}} f / \partial x_{1}^{r_{1}} \ldots \partial x_{d}^{r_{d}}$ plays a dominant part here and this fact is responsible for the name of these classes. Based on a Fourier multiplier theorem of Lizorkin one can prove a characterization of these classes using the Fourier transfom. Let $\mathscr{S}\left(\mathbb{R}^{d}\right)$ denote the class of all complex-valued rapidly decreasing infinitely differentiable functions defined on $\mathbb{R}^{d}$. By $\mathscr{S}^{\prime}\left(\mathbb{R}^{d}\right)$ we mean the collection of all tempered distributions and $\mathscr{F}$ and $\mathscr{F}^{-1}$ denote the Fourier transform and its inverse, respectively, both defined on $\mathscr{S}^{\prime}\left(\mathbb{R}^{d}\right)$. Then $f \in \mathscr{S}^{\prime}\left(\mathbb{R}^{d}\right)$ belongs to $S_{p}^{\bar{r}} W\left(\mathbb{R}^{d}\right)$ if and only if

$$
\mathscr{F}^{-1}\left(\left(1+\left|\xi_{1}\right|^{2}\right)^{r_{1} / 2} \ldots\left(1+\left|\xi_{d}\right|^{2}\right)^{r_{d} / 2} \mathscr{F} f(\xi)\right)(\cdot) \in L_{p}\left(\mathbb{R}^{d}\right) .
$$

Furthermore, the norms $\left\|f \mid S_{p}^{\bar{r}} W\left(\mathbb{R}^{d}\right)\right\|$ and

$$
\begin{equation*}
\left\|f\left|S_{p}^{\bar{r}} W\left(\mathbb{R}^{d}\right)\left\|^{*}:=\right\| \mathscr{F}^{-1}\left(\prod_{i=1}^{d}\left(1+\left|\xi_{i}\right|^{2}\right)^{r_{i} / 2} \mathscr{F} f(\xi)\right)(\cdot)\right| L_{p}\left(\mathbb{R}^{d}\right)\right\| \tag{2.2}
\end{equation*}
$$

are equivalent, cf. e.g. [18, 2.3.1]. The Fourier-analytic description can be taken to generalize these Sobolev spaces to fractional and negative order of smoothness, cf. [18, Chapt. 2]. We will take (2.2) as the definition of $S_{p}^{\bar{r}} W\left(\mathbb{R}^{d}\right)$ if $r=\left(r_{1}, \ldots, r_{d}\right)$, $r_{i} \in \mathbb{R}, i=1, \ldots, d$.

### 2.2. Some new function spaces.

As it will become clear below the description of the trace spaces will lead to some new Sobolev-type spaces. For us it will be sufficient to introduce these classes in the two-dimensional setting. For the rest of this section we concentrate on $p=2$.

Let $\mathscr{M}$ be a $2 \times 2$-matrix,

$$
\mathscr{M}=\left(\begin{array}{ll}
a & c  \tag{2.3}\\
b & d
\end{array}\right), \quad \operatorname{det} \mathscr{M} \neq 0, \quad \text { and let } \quad \vec{\eta}_{1}=\binom{a}{b}, \quad \vec{\eta}_{2}=\binom{c}{d}
$$

Let $f: \mathbb{R}^{2} \rightarrow \mathbb{R}$ be given. Then by $f \circ \mathscr{M}$ we mean the composition of $\mathscr{M}$ with $f$, i.e.

$$
(f \circ \mathscr{M})(x):=f(\mathscr{M} x)=f\left(a x_{1}+c x_{2}, b x_{1}+d x_{2}\right), \quad x=\left(x_{1}, x_{2}\right) \in \mathbb{R}^{2}
$$

Definition 2.1. Let $\mathscr{M}, \overrightarrow{\eta_{1}}, \overrightarrow{\eta_{2}}$ be as in (2.3). Let $r_{1}, r_{2} \in \mathbb{R}$. Then $S_{2}^{r_{1}, r_{2}} W(\mathscr{M}$, $\left.\mathbb{R}^{2}\right)$ denotes the collection of all tempered distributions $f \in \mathscr{S}^{\prime}\left(\mathbb{R}^{2}\right)$ such that $f \circ \mathscr{M} \in$ $S_{2}^{r_{1}, r_{2}} W\left(\mathbb{R}^{2}\right)$. We endow this class with the norm

$$
\left\|f\left|S_{2}^{r_{1}, r_{2}} W\left(\mathscr{M}, \mathbb{R}^{2}\right)\|:=\| f \circ \mathscr{M}\right| S_{2}^{r_{1}, r_{2}} W\left(\mathbb{R}^{2}\right)\right\|
$$

The following properties of these classes are immediate.

Lemma 2.2. Let $\mathscr{M}, \overrightarrow{\eta_{1}}, \overrightarrow{\eta_{2}}$ be as in (2.3). Let $r_{1}, r_{2} \in \mathbb{R}$.
(i) The classes $S_{2}^{r_{1}, r_{2}} W\left(\mathscr{M}, \mathbb{R}^{2}\right)$ are Banach spaces continuously embedded into $\mathscr{S}^{\prime}\left(\mathbb{R}^{2}\right)$.
(ii) $C_{0}^{\infty}\left(\mathbb{R}^{2}\right)$ is a dense subset of $S_{2}^{r_{1}, r_{2}} W\left(\mathscr{M}, \mathbb{R}^{2}\right)$.
(iii) A function $f \in \mathscr{S}^{\prime}\left(\mathbb{R}^{2}\right)$ belongs to $S_{2}^{r_{1}, r_{2}} W\left(\mathscr{M}, \mathbb{R}^{2}\right)$ if and only if

$$
\left(1+\left|a \xi_{1}+b \xi_{2}\right|^{2}\right)^{r_{1} / 2}\left(1+\left|c \xi_{1}+d \xi_{2}\right|^{2}\right)^{r_{2} / 2}|\mathscr{F} f(\xi)| \in L_{2}\left(\mathbb{R}^{2}\right) .
$$

Furthermore, the expression

$$
\left\|\left(1+\left|a \xi_{1}+b \xi_{2}\right|^{2}\right)^{r_{1} / 2}\left(1+\left|c \xi_{1}+d \xi_{2}\right|^{2}\right)^{r_{2} / 2} \mid \mathscr{F} f(\xi)\right\| L_{2}\left(\mathbb{R}^{2}\right) \|
$$

yields an equivalent norm in $S_{2}^{r_{1}, r_{2}} W\left(\mathscr{M}, \mathbb{R}^{2}\right)$.
For $r_{1}, r_{2} \in \mathbb{N}_{0}$ there is another interpretation. As usual, by $\partial f / \partial \vec{\eta}$ we denote the weak directional derivative of $f$ in the direction of $\vec{\eta}$.

Definition 2.3. Let $\overrightarrow{\eta_{1}}, \overrightarrow{\eta_{2}}$ be linearly independent vectors in $\mathbb{R}^{2}$. Let $r_{1}, r_{2} \in \mathbb{N}_{0}$. Then $S_{2}^{r_{1}, r_{2}} W\left(\overrightarrow{\eta_{1}}, \overrightarrow{\eta_{2}}\right)$ denotes the collection of all functions $f \in L_{2}\left(\mathbb{R}^{2}\right)$ such that

$$
\frac{\partial^{\alpha_{1}+\alpha_{2}} f}{\partial{\overrightarrow{\eta_{1}}}^{\alpha_{1}} \partial{\overrightarrow{\eta_{2}}}^{\alpha_{2}}} \in L_{2}\left(\mathbb{R}^{2}\right) \quad \text { for all } \alpha_{i} \leqslant r_{i}, i=1,2
$$

We endow this class with the norm

$$
\left\|f\left|S_{2}^{r_{1}, r_{2}} W\left(\overrightarrow{\eta_{1}}, \overrightarrow{\eta_{2}}\right)\left\|:=\sum_{\alpha_{1}=0}^{r_{1}} \sum_{\alpha_{2}=0}^{r_{2}}\right\| \frac{\partial^{\alpha_{1}+\alpha_{2}} f}{\partial{\overrightarrow{\eta_{1}}}^{\alpha_{1}} \partial{\overrightarrow{\eta_{2}}}^{\alpha_{2}}}\right| L_{2}\left(\mathbb{R}^{2}\right)\right\| .
$$

Remark 2.4. Obviously, these classes $S_{2}^{r_{1}, r_{2}} W\left(\overrightarrow{\eta_{1}}, \overrightarrow{\eta_{2}}\right)$ are Banach spaces. Let $\vec{e}_{1}, \vec{e}_{2}$ denote the elements of the canonical basis of $\mathbb{R}^{2}$. Then we have $S_{2}^{r_{1}, r_{2}} W\left(\vec{e}_{1}, \vec{e}_{2}\right)=S_{2}^{r_{1}, r_{2}} W\left(\mathbb{R}^{2}\right)$. Furthermore, $C_{0}^{\infty}\left(\mathbb{R}^{2}\right)$ is a dense set in $S_{2}^{r_{1}, r_{2}} W\left(\overrightarrow{\eta_{1}}\right.$, $\overrightarrow{\eta_{2}}$ ) for arbitrary vectors $\overrightarrow{\eta_{1}}, \overrightarrow{\eta_{2}}$.

For a smooth function $f$ and $\mathscr{M}$ and $\vec{\eta}_{1}, \vec{\eta}_{2}$ as in (2.3) it follows

$$
\frac{\partial}{\partial x_{1}}(f \circ \mathscr{M})(x)=\left\langle\nabla f(\mathscr{M} x), \vec{\eta}_{1}\right\rangle=\frac{\partial f}{\partial \vec{\eta}_{1}}(\mathscr{M} x) .
$$

By an induction argument we conclude

$$
\frac{\partial^{\alpha_{1}+\alpha_{2}}}{\partial x_{1}^{\alpha_{1}} \partial x_{2}^{\alpha_{2}}}(f \circ \mathscr{M})(x)=\frac{\partial^{\alpha_{1}+\alpha_{2}} f}{\partial{\overrightarrow{\eta_{1}}}^{\alpha_{1}} \partial{\overrightarrow{\eta_{2}}}^{\alpha_{2}}}(\mathscr{M} x) .
$$

Using the density of smooth compactly supported functions this proves the following.
Lemma 2.5. Let $\mathscr{M}, \vec{\eta}_{1}$, and $\vec{\eta}_{2}$ be as in (2.3). A function $f \in L_{2}\left(\mathbb{R}^{2}\right)$ belongs to $S_{2}^{r_{1}, r_{2}} W\left(\overrightarrow{\eta_{1}}, \overrightarrow{\eta_{2}}\right)$ if and only if the function $f \circ \mathscr{M}$ belongs to $S_{2}^{r_{1}, r_{2}} W\left(\mathbb{R}^{2}\right)$. Furthermore, the norms $\left\|f \mid S_{2}^{r_{1}, r_{2}} W\left(\overrightarrow{\eta_{1}}, \overrightarrow{\eta_{2}}\right)\right\|$ and $\left\|f \circ \mathscr{M} \mid S_{2}^{r_{1}, r_{2}} W\left(\mathbb{R}^{2}\right)\right\|$ are equivalent.

### 2.3. The trace with respect to an arbitrary orthogonal basis of the

 hyperplane.Let $A_{1}\left(\mathbb{R}^{3}\right)$ be a class of functions (distributions) defined on $\mathbb{R}^{3}$ and let $C\left(\mathbb{R}^{3}\right)$ be the collection of all continuous functions on $\mathbb{R}^{3}$. By $\widetilde{\Gamma}$ we denote a hyperplane in $\mathbb{R}^{3}$. Then we consider the mapping

$$
T:\left.f \rightarrow f\right|_{\tilde{\Gamma}}
$$

which is well-defined in the case of a continuous function $f$. The aim of this paper consists in determining a class of functions $A_{2}\left(\mathbb{R}^{2}\right) \hookrightarrow \mathscr{S}^{\prime}\left(\mathbb{R}^{2}\right)$ such that $T$, originally defined on $A_{1}\left(\mathbb{R}^{3}\right) \cap C\left(\mathbb{R}^{3}\right)$, extends to a linear, continuous and surjective mapping belonging to $\mathscr{L}\left(A_{1}\left(\mathbb{R}^{3}\right), A_{2}\left(\mathbb{R}^{2}\right)\right)$. If there exists a linear and continuous operator ext $\in \mathscr{L}\left(A_{2}\left(\mathbb{R}^{2}\right), A_{1}\left(\mathbb{R}^{3}\right)\right)$ such that $T \circ$ ext $=I$ (identity on $A_{2}\left(\mathbb{R}^{2}\right)$, we shall call $T$ a retraction and ext its corresponding coretraction.

In the monographs Amanov [2, 9.5], Gelman, Maz'ya [7, 2.3] and Schmeißer, Triebel [18, 2.4.2] the traces of function spaces with dominating mixed smoothness on hyperplanes parallel to the coordinate axes were studied. For simplicity let the hyperplane be given by $x_{3}=0$. Then the result is the following.

Proposition 2.6. Let $r_{3}>1 / 2$. Then the mapping

$$
T: f\left(x_{1}, x_{2}, x_{3}\right) \rightarrow f\left(x_{1}, x_{2}, 0\right)
$$

extends to a retraction from $S_{2}^{r_{1}, r_{2}} W\left(\mathbb{R}^{3}\right)$ onto $S_{2}^{r_{1}, r_{2}} W\left(\mathbb{R}^{2}\right)$.
Remark 2.7. A few comments are in order. First of all, $\mathscr{S}\left(\mathbb{R}^{3}\right)$ is dense in the class $S_{p}^{r_{1}, r_{2}} W\left(\mathbb{R}^{3}\right)$. So the trace operator is the unique linear extension of the mapping $T$. Secondly, there is a natural coordinate system on the hyperplane $x_{3}=0$ to measure the smoothness of the trace, namely the one induced by the unit vectors $\vec{e}_{1}=(1,0)$ and $\vec{e}_{2}=(0,1)$. Notice that the spaces $S_{p}^{r_{1}, r_{2}} W\left(\mathbb{R}^{2}\right)$ are not invariant under rotations in general.

In this paper we investigate the trace with respect to the hyperplane

$$
\Gamma:=\left\{\left(x_{1}, x_{2}, x_{3}\right) \in \mathbb{R}^{3}: x_{1}+x_{2}+x_{3}=0\right\}
$$

with $\Gamma$ as a model case for a hyperplane in an oblique position. However, taking the trace with respect to the hyperplane

$$
\Gamma_{\bar{\gamma}}:=\left\{\left(x_{1}, x_{2}, x_{3}\right): \gamma_{1} x_{1}+\gamma_{2} x_{2}+\gamma_{3} x_{3}=0\right\}, \quad \bar{\gamma}=\left(\gamma_{1}, \gamma_{2}, \gamma_{3}\right)
$$

where $\gamma_{1} \cdot \gamma_{2} \cdot \gamma_{3} \neq 0$, would give us the same result (up to the norms of the operators considered). This statement relies on the fact that the mapping

$$
f\left(x_{1}, x_{2}, x_{3}\right) \rightarrow f\left(\lambda_{1} x_{1}, \lambda_{2} x_{2}, \lambda_{3} x_{3}\right), \quad \lambda_{1} \cdot \lambda_{2} \cdot \lambda_{3} \neq 0
$$

is a bounded bijective mapping of $S_{2}^{\bar{r}} W\left(\mathbb{R}^{3}\right)$ onto itself.

The "natural" trace operators

$$
\begin{align*}
& \left(\operatorname{tr}_{1} f\right)\left(x_{2}, x_{3}\right)=f\left(-x_{2}-x_{3}, x_{2}, x_{3}\right),  \tag{2.4}\\
& \left(\operatorname{tr}_{2} f\right)\left(x_{1}, x_{3}\right)=f\left(x_{1},-x_{1}-x_{3}, x_{3}\right),  \tag{2.5}\\
& \left(\operatorname{tr}_{3} f\right)\left(x_{1}, x_{2}\right)=f\left(x_{1}, x_{2},-x_{1}-x_{2}\right), \tag{2.6}
\end{align*}
$$

and the trace operator $\operatorname{tr}_{\mathscr{O}} f$, see (1.1) and (1.2), are connected through

$$
\begin{align*}
\left(\operatorname{tr}_{\mathscr{O}} f\right)\left(z_{1}, z_{2}\right) & =f\left(z_{1} \vec{\sigma}_{1}+z_{2} \vec{\sigma}_{2}\right)  \tag{2.7}\\
& =f\left(\sigma_{1,1} z_{1}+\sigma_{2,1} z_{2}, \sigma_{1,2} z_{1}+\sigma_{2,2} z_{2}, \sigma_{1,3} z_{1}+\sigma_{2,3} z_{2}\right) \\
& =\left(\operatorname{tr}_{1} f\right)\left(\sigma_{1,2} z_{1}+\sigma_{2,2} z_{2}, \sigma_{1,3} z_{1}+\sigma_{2,3} z_{2}\right) \\
& =\left(\operatorname{tr}_{1} f\right)\left(\mathscr{R}_{1} \vec{z}\right),
\end{align*}
$$

where

$$
\mathscr{R}_{1}=\left(\begin{array}{cc}
\sigma_{1,2} & \sigma_{2,2}  \tag{2.8}\\
\sigma_{1,3} & \sigma_{2,3}
\end{array}\right) \quad \text { and } \quad \vec{z}=\binom{z_{1}}{z_{2}} .
$$

Analoguously one obtains

$$
\begin{equation*}
\left(\operatorname{tr}_{\mathscr{O}} f\right)\left(z_{1}, z_{2}\right)=\left(\operatorname{tr}_{2} f\right)\left(\mathscr{R}_{2} \vec{z}\right)=\left(\operatorname{tr}_{3} f\right)\left(\mathscr{R}_{3} \vec{z}\right), \tag{2.9}
\end{equation*}
$$

with

$$
\mathscr{R}_{2}=\left(\begin{array}{ll}
\sigma_{1,1} & \sigma_{2,1}  \tag{2.10}\\
\sigma_{1,3} & \sigma_{2,3}
\end{array}\right), \quad \mathscr{R}_{3}=\left(\begin{array}{ll}
\sigma_{1,1} & \sigma_{2,1} \\
\sigma_{1,2} & \sigma_{2,2}
\end{array}\right) .
$$

The linear independence of the vectors $\vec{\sigma}_{1}, \vec{\sigma}_{2}$, combined with $\vec{\sigma}_{1}, \vec{\sigma}_{2} \in \Gamma$, ensure that the matrices $\mathscr{R}_{1}, \mathscr{R}_{2}, \mathscr{R}_{3}$ are regular.

In what follows we shall determine the regularity of $\operatorname{tr}_{\mathscr{O}} f$ as well as of $\operatorname{tr}_{i} f$, $i=1,2,3$.

Above we considered all orthogonal bases of $\Gamma$. Probably it would be more natural to restrict to orthonormal bases. However, all function spaces under consideration here remain invariant under the change from an orthogonal to the associated orthonormal basis (up to equivalent quasi-norms). The greater generality leads to nothing new but it simplifies the calculations. For that reason we shall work with orthogonal bases.

### 2.4. The regularity of $\operatorname{tr}_{\mathscr{O}} f$.

### 2.4.1. A description of the general case.

Let $f \in C_{0}^{\infty}\left(\mathbb{R}^{3}\right)$. Now we introduce a very useful decomposition of $f$. Let $\mathscr{X}_{i}$ denote the characteristic function of the set

$$
M_{i}:=\left\{\left(\tau_{1}, \tau_{2}, \tau_{3}\right):\left|\tau_{i}\right|=\min \left(\left|\tau_{1}\right|,\left|\tau_{2}\right|,\left|\tau_{3}\right|\right)\right\}, \quad i=1,2,3
$$

Hence

$$
\left|M_{i} \cap M_{j}\right|=0, \quad i \neq j, \quad \text { and } \quad \bigcup_{i=1}^{3} M_{i}=\mathbb{R}^{3}
$$

(here $|\cdot|$ denotes the Lebesgue measure in $\mathbb{R}^{3}$ ). We put

$$
f_{i}(x):=\mathscr{F}^{-1}\left[\mathscr{X}_{i}(\xi) \mathscr{F} f(\xi)\right](x)
$$

and obtain $f=f_{1}+f_{2}+f_{3}$. We continue by defining

$$
\begin{equation*}
g_{i}\left(x_{1}, x_{2}\right)=\left(\operatorname{tr}_{i} f_{i}\right)\left(x_{1}, x_{2}\right), \quad i=1,2,3 \tag{2.11}
\end{equation*}
$$

Elementary properties of the Fourier transform yield

$$
\begin{align*}
\mathscr{F}_{2} g_{1}\left(\xi_{1}, \xi_{2}\right) & =\frac{1}{\sqrt{2 \pi}} \int_{\mathbb{R}} \mathscr{F}_{3} f_{1}\left(\tau_{1}, \xi_{1}+\tau_{1}, \xi_{2}+\tau_{1}\right) \mathrm{d} \tau_{1},  \tag{2.12}\\
\mathscr{F}_{2} g_{2}\left(\xi_{1}, \xi_{2}\right) & =\frac{1}{\sqrt{2 \pi}} \int_{\mathbb{R}} \mathscr{F}_{3} f_{2}\left(\xi_{1}+\tau_{2}, \tau_{2}, \xi_{2}+\tau_{2}\right) \mathrm{d} \tau_{2}, \\
\mathscr{F}_{2} g_{3}\left(\xi_{1}, \xi_{2}\right) & =\frac{1}{\sqrt{2 \pi}} \int_{\mathbb{R}} \mathscr{F}_{3} f_{3}\left(\xi_{1}+\tau_{3}, \xi_{2}+\tau_{3}, \tau_{3}\right) \mathrm{d} \tau_{3},
\end{align*}
$$

where $\mathscr{F}_{2} g$ denotes the Fourier transform in $\mathbb{R}^{2}$ and $\mathscr{F}_{3} f$ the Fourier transform in $\mathbb{R}^{3}$, respectively. Now we are going to check the regularity of the functions $g_{i}$. To begin with we investigate the case when $i=1$. Let $r_{1}>1 / 2$. By using Hölder's inequality and a change of variables we obtain

$$
\begin{aligned}
\int_{\mathbb{R}^{2}} & \left(1+\xi_{2}^{2}\right)^{r_{2}}\left(1+\xi_{3}^{2}\right)^{r_{3}}\left|\int_{\mathbb{R}} \mathscr{F}_{3} f_{1}\left(\tau_{1}, \xi_{2}+\tau_{1}, \xi_{3}+\tau_{1}\right) \mathrm{d} \tau_{1}\right|^{2} \mathrm{~d} \xi_{1} \mathrm{~d} \xi_{2} \\
& \leqslant c_{1} \int_{\mathbb{R}^{3}}\left(1+\xi_{2}^{2}\right)^{r_{2}}\left(1+\xi_{3}^{2}\right)^{r_{3}}\left(1+\tau_{1}^{2}\right)^{r_{1}}\left|\mathscr{F}_{3} f_{1}\left(\tau_{1}, \xi_{2}+\tau_{1}, \xi_{3}+\tau_{1}\right)\right|^{2} \mathrm{~d} \tau_{1} \mathrm{~d} \xi_{1} \mathrm{~d} \xi_{2} \\
& =c_{1} \int_{\mathbb{R}^{3}}\left[1+\left(\tau_{2}-\tau_{1}\right)^{2}\right]^{r_{2}}\left[1+\left(\tau_{3}-\tau_{1}\right)^{2}\right]^{r_{3}}\left(1+\tau_{1}^{2}\right)^{r_{1}}\left|\mathscr{F}_{3} f_{1}\left(\tau_{1}, \tau_{2}, \tau_{3}\right)\right|^{2} \mathrm{~d} \vec{\tau}
\end{aligned}
$$

with $c_{1}=\int_{\mathbb{R}}\left(1+\tau_{1}^{2}\right)^{-r_{1}} \mathrm{~d} \tau_{1}<\infty$. Finally, we observe that if $\left|\tau_{1}\right| \leqslant \min \left(\left|\tau_{2}\right|,\left|\tau_{3}\right|\right)$, then $\left|\tau_{2}-\tau_{1}\right| \leqslant 2\left|\tau_{2}\right|,\left|\tau_{3}-\tau_{1}\right| \leqslant 2\left|\tau_{3}\right|$ and

$$
\left[1+\left(\tau_{2}-\tau_{1}\right)^{2}\right]^{r_{2}}\left[1+\left(\tau_{3}-\tau_{1}\right)^{2}\right]^{r_{3}} \leqslant 4^{r_{2}+r_{3}}\left(1+\tau_{2}^{2}\right)^{r_{2}}\left(1+\tau_{3}^{2}\right)^{r_{3}}
$$

Because of supp $\mathscr{F}_{3} f_{1} \subset\left\{\left(\tau_{1}, \tau_{2}, \tau_{3}\right) \in \mathbb{R}^{3}:\left|\tau_{1}\right| \leqslant \min \left(\left|\tau_{2}\right|,\left|\tau_{3}\right|\right)\right\}$, we finally conclude

$$
\begin{equation*}
\left\|\operatorname{tr}_{1} f_{1}\left|S_{2}^{r_{2}, r_{3}} W\left(\mathbb{R}^{2}\right)\left\|\leqslant c_{2}\right\| f_{1}\right| S_{2}^{\bar{r}} W\left(\mathbb{R}^{3}\right)\right\| \leqslant c_{2}\left\|f \mid S_{2}^{\bar{r}} W\left(\mathbb{R}^{3}\right)\right\| . \tag{2.13}
\end{equation*}
$$

This proves $\operatorname{tr}_{1} f_{1} \in S_{2}^{r_{2}, r_{3}} W\left(\mathbb{R}^{2}\right)$. Similarly one obtains $\operatorname{tr}_{2} f_{2} \in S_{2}^{r_{1}, r_{3}} W\left(\mathbb{R}^{2}\right)$ (if $\left.r_{2}>1 / 2\right)$ and $\operatorname{tr}_{3} f_{3} \in S_{2}^{r_{1}, r_{2}} W\left(\mathbb{R}^{2}\right)$ (if $r_{3}>1 / 2$ ), respectively. To summarize our findings we need to recall a further notion. For three quasi-Banach spaces $A_{1}, A_{2}, A_{3} \hookrightarrow \mathscr{S}^{\prime}\left(\mathbb{R}^{2}\right)$ of tempered distributions we put

$$
A_{1}+A_{2}+A_{3}:=\left\{g \in \mathscr{S}^{\prime}\left(\mathbb{R}^{2}\right): \exists g_{i} \in A_{i}, i=1,2,3, \quad \text { s.t. } \quad g=g_{1}+g_{2}+g_{3}\right\} .
$$

We equip this space with a quasi-norm by taking

$$
\left\|g \mid A_{1}+A_{2}+A_{3}\right\|:=\inf \left\{\sum_{i=1}^{3}\left\|g_{i} \mid A_{i}\right\|: g=g_{1}+g_{2}+g_{3}, g_{i} \in A_{i}, i=1,2,3\right\}
$$

Lemma 2.8. Let $\mathscr{O}$ be an orthogonal basis of $\Gamma$ and let $\mathscr{R}_{i}, i=1,2,3$, be matrices associated with $\mathscr{O}$ by (1.1), (2.8) and (2.10).

Suppose $\min \left(r_{1}, r_{2}, r_{3}\right)>1 / 2$. Then $\operatorname{tr}_{\mathscr{O}}$ becomes a continuous mapping of $S_{2}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right)$ into

$$
\begin{equation*}
S_{2}^{r_{1}, r_{2}} W\left(\mathscr{R}_{3}^{-1}, \mathbb{R}^{2}\right)+S_{2}^{r_{1}, r_{3}} W\left(\mathscr{R}_{2}^{-1}, \mathbb{R}^{2}\right)+S_{2}^{r_{2}, r_{3}} W\left(\mathscr{R}_{1}^{-1}, \mathbb{R}^{2}\right) \tag{2.14}
\end{equation*}
$$

Proof. The boundedness of $\operatorname{tr}_{\mathscr{O}}$ follows from the identity

$$
\left(\operatorname{tr}_{\mathscr{O}} f\right)(\vec{z})=\sum_{i=1}^{3}\left(\operatorname{tr}_{\mathscr{O}} f_{i}\right)(\vec{z})=\sum_{i=1}^{3}\left(\operatorname{tr}_{i} f_{i}\right)\left(\mathscr{R}_{i} \vec{z}\right)
$$

cf. (2.7), (2.9), the definition of the spaces $S_{2}^{r_{1}, r_{2}} W\left(\mathscr{M}, \mathbb{R}^{2}\right)$ and the inequality (2.13) and its counterparts for $\operatorname{tr}_{2}$ and $\operatorname{tr}_{3}$.

The restriction $\min \left(r_{1}, r_{2}, r_{3}\right)>1 / 2$ has been convenient but is by no means necessary. Moreover, as we shall see by the next theorem the operator $\operatorname{tr}_{\mathscr{O}}$ is surjective in Lemma 2.8. The description of the trace class becomes more complicated than in Lemma 2.8 if $\min \left(r_{1}, r_{2}, r_{3}\right)<1 / 2$.

Theorem 2.9. Let $\mathscr{O}$ be an orthogonal basis of $\Gamma$ and let $\mathscr{R}_{i}, i=1,2,3$, be matrices associated with $\mathscr{O}$ by (1.1), (2.8) and (2.10).

Let $\bar{r}=\left(r_{1}, r_{2}, r_{3}\right) \in \mathbb{R}^{3}$ with $r_{i} \neq 1 / 2, i=1,2,3$, and

$$
\begin{equation*}
\min \left(r_{1}, r_{2}, r_{3}, r_{1}+r_{2}-\frac{1}{2}, r_{1}+r_{3}-\frac{1}{2}, r_{2}+r_{3}-\frac{1}{2}\right)>0 . \tag{2.15}
\end{equation*}
$$

Then

$$
\begin{equation*}
\operatorname{tr}_{\mathscr{O}} \in \mathscr{L}\left(S_{2}^{\bar{r}} W\left(\mathbb{R}^{3}\right), S^{1}\left(\mathbb{R}^{2}\right)+S^{2}\left(\mathbb{R}^{2}\right)+S^{3}\left(\mathbb{R}^{2}\right)\right) \tag{2.16}
\end{equation*}
$$

where

$$
S^{1}\left(\mathbb{R}^{2}\right):= \begin{cases}S_{2}^{r_{2}, r_{3}} W\left(\mathscr{R}_{1}^{-1}, \mathbb{R}^{2}\right), & \text { if } r_{1}>\frac{1}{2}, \\ S_{2}^{r_{2}, r_{3}+r_{1}-\frac{1}{2}} W\left(\mathscr{R}_{1}^{-1}, \mathbb{R}^{2}\right) \cap S_{2}^{r_{2}+r_{1}-\frac{1}{2}, r_{3}} W\left(\mathscr{R}_{1}^{-1}, \mathbb{R}^{2}\right), & \text { if } r_{1}<\frac{1}{2}\end{cases}
$$

and similarly for $S^{2}$ and $S^{3}$.
Conversely, to each function $g \in S^{1}\left(\mathbb{R}^{2}\right)+S^{2}\left(\mathbb{R}^{2}\right)+S^{3}\left(\mathbb{R}^{2}\right)$ there exists a function $f \in S_{2}^{\bar{r}} W\left(\mathbb{R}^{3}\right)$ such that $\operatorname{tr}_{\mathscr{O}} f=g$.

Proof. Step 1. Preparations. For $\alpha, \beta, t \in \mathbb{R}$ we define

$$
I(\alpha, \beta, t):=\int_{-\infty}^{\infty}\left(1+(t+\tau)^{2}\right)^{-\alpha}\left(1+\tau^{2}\right)^{-\beta} \mathrm{d} \tau
$$

If $\alpha+\beta>\frac{1}{2}, \beta<\frac{1}{2}$, elementary calculations yield

$$
I(\alpha, \beta, t) \leqslant c \begin{cases}\left(1+t^{2}\right)^{-\beta} & \text { if } \alpha>\frac{1}{2}  \tag{2.17}\\ \left(1+t^{2}\right)^{-\beta}(1+\log (1+|t|)) & \text { if } \alpha=\frac{1}{2} \\ \left(1+t^{2}\right)^{-(\alpha+\beta)+1 / 2} & \text { if } \alpha<\frac{1}{2}\end{cases}
$$

for some $c$ independent of $t$.
Step 2. The boundedness of $\operatorname{tr}_{\mathscr{O}}$ in the case of $\min \left(r_{1}, r_{2}, r_{3}\right)>\frac{1}{2}$ has been proven before.

Now we suppose $0<r_{1}<\frac{1}{2}$. We proceed as at the beginning of this subsection and obtain

$$
\begin{aligned}
& \int_{\mathbb{R}^{2}}\left(1+\xi_{2}^{2}\right)^{r_{2}}\left(1+\xi_{3}^{2}\right)^{r_{3}+r_{1}-\frac{1}{2}}\left|\int_{\mathbb{R}} \mathscr{F}_{3} f_{1}\left(\tau_{1}, \xi_{2}+\tau_{1}, \xi_{3}+\tau_{1}\right) \mathrm{d} \tau_{1}\right|^{2} \mathrm{~d} \xi_{1} \mathrm{~d} \xi_{2} \\
& \leqslant \int_{\mathbb{R}^{3}}\left(1+\xi_{2}^{2}\right)^{r_{2}}\left(1+\xi_{3}^{2}\right)^{r_{3}+r_{1}-\frac{1}{2}} I\left(\alpha, r_{1}, \xi_{3}\right)\left(1+\tau_{1}^{2}\right)^{r_{1}}\left(1+\left(\tau_{1}+\xi_{3}\right)^{2}\right)^{\alpha} \\
& \quad \times\left|\mathscr{F}_{3} f_{1}\left(\tau_{1}, \xi_{2}+\tau_{1}, \xi_{3}+\tau_{1}\right)\right|^{2} \mathrm{~d} \tau_{1} \mathrm{~d} \xi_{1} \mathrm{~d} \xi_{2} \\
& \leqslant c_{1} \int_{\mathbb{R}^{3}}\left[1+\left(\tau_{2}-\tau_{1}\right)^{2}\right]^{r_{2}}\left[1+\left(\tau_{3}-\tau_{1}\right)^{2}\right]^{r_{3}-\alpha}\left(1+\tau_{1}^{2}\right)^{r_{1}}\left(1+\tau_{3}^{2}\right)^{\alpha}\left|\mathscr{F}_{3} f_{1}\left(\tau_{1}, \tau_{2}, \tau_{3}\right)\right|^{2} \mathrm{~d} \vec{\tau}
\end{aligned}
$$

where we have used (2.17) with some $\alpha$ satisfying $\frac{1}{2}-r_{1}<\alpha<\frac{1}{2}$. Choosing $\alpha$ sufficiently close to $\frac{1}{2}-r_{1}$ our restriction $r_{1}+r_{3}>\frac{1}{2}$, see (2.15), guarantees $r_{3}-\alpha \geqslant 0$. Furthermore, taking into account the information on the support of $\mathscr{F} f_{1}$ we arrive at

$$
\left\|\operatorname{tr}_{1} f_{1}\left|S_{2}^{r_{2}, r_{3}+r_{1}-\frac{1}{2}} W\left(\mathbb{R}^{2}\right)\left\|\leqslant c_{2}\right\| f_{1}\right| S_{2}^{\bar{r}} W\left(\mathbb{R}^{2}\right)\right\| \leqslant c_{2}\left\|f \mid S_{2}^{\bar{r}} W\left(\mathbb{R}^{2}\right)\right\|
$$

with some $c$ independent of $f$. Interchanging the roles of $\xi_{1}$ and $\xi_{2}$ also

$$
\left\|\operatorname{tr}_{1} f_{1}\left|S_{2}^{r_{2}+r_{1}-\frac{1}{2}, r_{3}} W\left(\mathbb{R}^{2}\right)\left\|\leqslant c_{3}\right\| f\right| S_{2}^{\bar{r}} W\left(\mathbb{R}^{2}\right)\right\|
$$

follows. Moreover, by symmetry we obtain the needed estimates of $\operatorname{tr}_{i} f_{1}, i=2,3$, as well. This completes the proof of the boundedness.

Step 3. Construction of an extension operator.
Substep 3.1. Construction of a linear extension operator for $S_{2}^{r_{1}, r_{2}} W\left(\mathbb{R}^{2}\right)$. Let $\varphi \in C_{0}^{\infty}(\mathbb{R})$ be a function such that $\int \varphi(t) \mathrm{d} t=\sqrt{2 \pi}$. Then, for $g \in C_{0}^{\infty}\left(\mathbb{R}^{2}\right)$ and $x \in \mathbb{R}^{3}$, we define

$$
\begin{aligned}
& f_{1}(x)=\operatorname{ext}_{1}^{*} g(x):=\mathscr{F}_{3}^{-1}\left[\varphi\left(\xi_{1}\right) \mathscr{F}_{2} g\left(\xi_{2}-\xi_{1}, \xi_{3}-\xi_{1}\right)\right](x), \\
& f_{2}(x)=\operatorname{ext}_{2}^{*} g(x):=\mathscr{F}_{3}^{-1}\left[\varphi\left(\xi_{2}\right) \mathscr{F}_{2} g\left(\xi_{1}-\xi_{2}, \xi_{3}-\xi_{2}\right)\right](x), \\
& f_{3}(x)=\operatorname{ext}_{3}^{*} g(x):=\mathscr{F}_{3}^{-1}\left[\varphi\left(\xi_{3}\right) \mathscr{F}_{2} g\left(\xi_{1}-\xi_{3}, \xi_{2}-\xi_{3}\right)\right](x) .
\end{aligned}
$$

Hence, e.g. for $f_{3}$, we conclude

$$
\frac{1}{\sqrt{2 \pi}} \int_{\mathbb{R}} \mathscr{F}_{3} f_{3}\left(\xi_{1}+\tau_{3}, \xi_{2}+\tau_{3}, \tau_{3}\right) \mathrm{d} \tau_{3}=\mathscr{F}_{2} g\left(\xi_{1}, \xi_{2}\right)
$$

and from this identity we derive

$$
g\left(x_{1}, x_{2}\right)=\left(\operatorname{tr}_{3} f_{3}\right)\left(x_{1}, x_{2}\right)=f_{3}\left(x_{1}, x_{2},-x_{1}-x_{2}\right), \quad\left(x_{1}, x_{2}\right) \in \mathbb{R}^{2}
$$

Similarly

$$
g=\operatorname{tr}_{1} f_{1} \quad \text { and } \quad g=\operatorname{tr}_{2} f_{2}
$$

The regularity of $\operatorname{ext}_{3}^{*} g$ is easily checked in view of

$$
\begin{aligned}
\int_{\mathbb{R}^{3}} & \left(1+\xi_{1}^{2}\right)^{r_{1}}\left(1+\xi_{2}^{2}\right)^{r_{2}}\left(1+\xi_{3}^{2}\right)^{r_{3}}\left|\mathscr{F}_{3} \operatorname{ext}^{*} g(\xi)\right|^{2} \mathrm{~d} \vec{\xi} \\
& =\int_{\mathbb{R}^{3}}\left(1+\left|\xi_{1}+\tau_{3}\right|^{2}\right)^{r_{1}}\left(1+\left|\xi_{2}+\tau_{3}\right|^{2}\right)^{r_{2}}\left(1+\tau_{3}^{2}\right)^{r_{3}}\left|\varphi\left(\tau_{3}\right) \mathscr{F}_{2} g\left(\xi_{1}, \xi_{2}\right)\right|^{2} \mathrm{~d} \xi_{1} \mathrm{~d} \xi_{2} \mathrm{~d} \tau_{3} \\
& \leqslant c_{1} \int_{\mathbb{R}}\left(1+\left|\tau_{3}\right|^{2}\right)^{r_{1}+r_{2}+r_{3}}\left|\varphi\left(\tau_{3}\right)\right|^{2} \mathrm{~d} \tau_{3}\left\|g \mid S_{2}^{r_{1}, r_{2}} W\left(\mathbb{R}^{2}\right)\right\|^{2} \\
& \leqslant c_{2}\left\|g \mid S_{2}^{r_{1}, r_{2}} W\left(\mathbb{R}^{2}\right)\right\|^{2},
\end{aligned}
$$

where we also used the fact that $\varphi$ has compact support. This proves $\operatorname{ext}_{3}^{*} \in$ $\mathscr{L}\left(S_{2}^{r_{1}, r_{2}} W\left(\mathbb{R}^{2}\right), S_{2}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right)\right)$ for any $r_{3} \in \mathbb{R}$. Similarly, ext ${ }_{1}^{*} \in \mathscr{L}\left(S_{2}^{r_{2}, r_{3}} W\left(\mathbb{R}^{2}\right)\right.$, $\left.S_{2}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right)\right)$ for any $r_{1}$ and ext ${ }_{2}^{*} \in \mathscr{L}\left(S_{2}^{r_{1}, r_{3}} W\left(\mathbb{R}^{2}\right), S_{2}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right)\right)$ for any $r_{2}$, respectively.

Substep 3.2. Construction of an extension operator in the case of $\min \left(r_{1}\right.$, $\left.r_{2}, r_{3}\right)>1 / 2$. We shall use the abbreviations $A_{1}=S_{2}^{r_{2}, r_{3}} W\left(\mathscr{R}_{1}^{-1}, \mathbb{R}^{2}\right), A_{2}=$ $S_{2}^{r_{1}, r_{3}} W\left(\mathscr{R}_{2}^{-1}, \mathbb{R}^{2}\right)$ and $A_{3}=S_{2}^{r_{1}, r_{2}} W\left(\mathscr{R}_{3}^{-1}, \mathbb{R}^{2}\right)$. Let $g \in A_{1}+A_{2}+A_{3}$. Further, let $g=g_{1}+g_{2}+g_{3}$, where

$$
g_{i} \in A_{i}, \quad i=1,2,3, \quad \text { and } \quad\left\|g\left|A_{1}+A_{2}+A_{3}\left\|\leqslant 2 \sum_{i=1}^{3}\right\| g_{i}\right| A_{i}\right\| .
$$

By definition $g_{1}\left(\mathscr{R}_{1}^{-1} \cdot\right) \in S_{2}^{r_{2}, r_{3}} W\left(\mathbb{R}^{2}\right)$ and consequently, by Step 3.1, $f_{1}:=\operatorname{ext}_{1}^{*}$ $g_{1}\left(\mathscr{R}_{1}^{-1} \cdot\right) \in S_{2}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right)$. Similarly, $f_{2}:=\operatorname{ext}_{2}^{*} g_{2}\left(\mathscr{R}_{2}^{-1} \cdot\right), f_{3}:=\operatorname{ext}_{3}^{*} g_{3}\left(\mathscr{R}_{3}^{-1} \cdot\right) \in$ $S_{2}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right)$. We put $f:=f_{1}+f_{2}+f_{3}$. Because of

$$
\begin{aligned}
\operatorname{tr}_{\mathscr{O}} f & =\sum_{i=1}^{3} \operatorname{tr}_{\mathscr{O}} f_{i}=\sum_{i=1}^{3}\left(\operatorname{tr}_{i} f_{i}\right)\left(\mathscr{R}_{i} \cdot\right) \\
& =\sum_{i=1}^{3}\left(\operatorname{tr}_{i} \operatorname{ext}_{i}^{*} g_{i}\left(\mathscr{R}_{i}^{-1} \cdot\right)\right)\left(\mathscr{R}_{i} \cdot\right)=\sum_{i=1}^{3} g_{i}=g,
\end{aligned}
$$

see Substep 3.1, this proves the existence of a bounded extension of $g$ if $\min \left(r_{1}\right.$, $\left.r_{2}, r_{3}\right)>\frac{1}{2}$.

Substep 3.3. Let $0<r_{1}<\frac{1}{2}$. We shall use the abbreviations $A_{1}=$ $S_{2}^{r_{2}, r_{3}+r_{1}-\frac{1}{2}} W\left(\mathbb{R}^{2}\right), A_{2}=S_{2}^{r_{2}+r_{1}-\frac{1}{2}, r_{3}} W\left(\mathbb{R}^{2}\right) . \quad$ By the arguments from the previous substep (and by symmetry) it will be sufficient to construct a function $f_{1} \in S_{2}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right)$ such that $\operatorname{tr}_{1} f_{1}=g_{1}\left(\mathscr{R}_{1}^{-1}.\right) \in A_{1} \cap A_{2}$. To shorten the notation we write $h_{1}$ instead of $g_{1}\left(\mathscr{R}_{1}^{-1} \cdot\right)$. To begin with we define two subsets of $\mathbb{R}^{3}$

$$
\begin{aligned}
\Omega_{1}:= & \left\{\left(\xi_{1}, \xi_{2}, \xi_{3}\right):\left|\xi_{2}-\xi_{1}\right| \leqslant\left|\xi_{3}-\xi_{1}\right|, \frac{\left|\xi_{2}-\xi_{1}\right|}{4} \leqslant\left|\xi_{1}\right| \leqslant \frac{\left|\xi_{2}-\xi_{1}\right|}{2}\right. \\
& \text { if } \left.\left|\xi_{2}-\xi_{1}\right| \geqslant 1,\left|\xi_{1}\right| \leqslant 1 \text { if }\left|\xi_{2}-\xi_{1}\right|<1\right\}, \\
\Omega_{2}:= & \left\{\left(\xi_{1}, \xi_{2}, \xi_{3}\right):\left|\xi_{3}-\xi_{1}\right|<\left|\xi_{2}-\xi_{1}\right|, \frac{\left|\xi_{3}-\xi_{1}\right|}{4} \leqslant\left|\xi_{1}\right| \leqslant \frac{\left|\xi_{3}-\xi_{1}\right|}{2}\right. \\
& \text { if } \left.\left|\xi_{3}-\xi_{1}\right| \geqslant 1,\left|\xi_{1}\right| \leqslant 1 \text { if }\left|\xi_{3}-\xi_{1}\right|<1\right\} .
\end{aligned}
$$

Obviously, these sets are disjoint. Let $\mathscr{X}_{i}$ denote the characteristic function of $\Omega_{i}$, $i=1,2$. Then we define

$$
\begin{aligned}
f_{1}(x):= & \int \mathrm{e}^{\mathrm{i} x \xi} \mathscr{F}_{2} h_{1}\left(\xi_{2}-\xi_{1}, \xi_{3}-\xi_{1}\right) \\
& \times\left(\mathscr{X}_{1}(\xi) \frac{\left(1+\left(\xi_{2}-\xi_{1}\right)^{2}\right)^{r_{2}+r_{1}-1 / 2}\left(1+\left(\xi_{3}-\xi_{1}\right)^{2}\right)^{r_{3}}}{\left(1+\xi_{1}^{2}\right)^{r_{1}}\left(1+\xi_{2}^{2}\right)^{r_{2}}\left(1+\xi_{3}^{2}\right)^{r_{3}}} H^{1}\left(\xi_{2}-\xi_{1}, \xi_{3}-\xi_{1}\right)\right. \\
+ & \left.\mathscr{X}_{2}(\xi) \frac{\left(1+\left(\xi_{2}-\xi_{1}\right)^{2}\right)^{r_{2}}\left(1+\left(\xi_{3}-\xi_{1}\right)^{2}\right)^{r_{3}+r_{1}-1 / 2}}{\left(1+\xi_{1}^{2}\right)^{r_{1}}\left(1+\xi_{2}^{2}\right)^{r_{2}}\left(1+\xi_{3}^{2}\right)^{r_{3}}} H^{2}\left(\xi_{2}-\xi_{1}, \xi_{3}-\xi_{1}\right)\right) \mathrm{d} \vec{\xi},
\end{aligned}
$$

where the functions $H^{1}, H^{2}$ will be chosen later. First we prove $\operatorname{tr}_{1} f_{1}=h_{1}$. It is sufficient to assume $h_{1} \in C_{0}^{\infty}\left(\mathbb{R}^{2}\right)$. Setting $\tau_{2}=\xi_{2}-\xi_{1}$ and $\tau_{3}=\xi_{3}-\xi_{1}$ we find

$$
\begin{aligned}
& f_{1}\left(-x_{2}-x_{3}, x_{2}, x_{3}\right) \\
& =\int_{\left|\tau_{2}\right| \leqslant\left|\tau_{3}\right|} \mathrm{e}^{\mathrm{i}\left(x_{2} \tau_{2}+x_{3} \tau_{3}\right)} \mathscr{F}_{2} h_{1}\left(\tau_{2}, \tau_{3}\right)\left(1+\tau_{2}^{2}\right)^{r_{2}+r_{1}-1 / 2}\left(1+\tau_{3}^{2}\right)^{r_{3}} H^{1}\left(\tau_{2}, \tau_{3}\right) \\
& \quad \times \int_{I\left(\tau_{2}\right)} \frac{1}{\left(1+\xi_{1}^{2}\right)^{r_{1}}\left(1+\left(\tau_{2}+\xi_{1}\right)^{2}\right)^{r_{2}}\left(1+\left(\tau_{3}+\xi_{1}\right)^{2}\right)^{r_{3}}} \mathrm{~d} \xi_{1} \mathrm{~d} \tau_{2} \mathrm{~d} \tau_{3} \\
& \quad+\int_{\left|\tau_{3}\right|<\left|\tau_{2}\right|} \mathrm{e}^{\mathrm{i}\left(x_{2} \tau_{2}+x_{3} \tau_{3}\right)} \mathscr{F}_{2} h_{1}\left(\tau_{2}, \tau_{3}\right)\left(1+\tau_{2}^{2}\right)^{r_{2}}\left(1+\tau_{3}^{2}\right)^{r_{3}+r_{1}-1 / 2} H^{2}\left(\tau_{2}, \tau_{3}\right) \\
& \quad \times \int_{I\left(\tau_{3}\right)} \frac{1}{\left(1+\xi_{1}^{2}\right)^{r_{1}}\left(1+\left(\tau_{2}+\xi_{1}\right)^{2}\right)^{r_{2}}\left(1+\left(\tau_{3}+\xi_{1}\right)^{2}\right)^{r_{3}}} \mathrm{~d} \xi_{1} \mathrm{~d} \tau_{2} \mathrm{~d} \tau_{3}
\end{aligned}
$$

with $I\left(\tau_{2}\right)$ and $I\left(\tau_{3}\right)$ being appropriate subsets of $\mathbb{R}$. The functions $H^{1}$ and $H^{2}$ are determined through the identities

$$
\begin{aligned}
& H^{1}\left(\tau_{2}, \tau_{3}\right)=\frac{1}{2 \pi}\left(\int_{I\left(\tau_{2}\right)} \frac{\left(1+\tau_{2}^{2}\right)^{r_{2}+r_{1}-1 / 2}\left(1+\tau_{3}^{2}\right)^{r_{3}}}{\left(1+\xi_{1}^{2}\right)^{r_{1}}\left(1+\left(\tau_{2}+\xi_{1}\right)^{2}\right)^{r_{2}}\left(1+\left(\tau_{3}+\xi_{1}\right)^{2}\right)^{r_{3}}} \mathrm{~d} \xi_{1}\right)^{-1} \\
& H^{2}\left(\tau_{2}, \tau_{3}\right)=\frac{1}{2 \pi}\left(\int_{I\left(\tau_{3}\right)} \frac{\left(1+\tau_{2}^{2}\right)^{r_{2}}\left(1+\tau_{3}^{2}\right)^{r_{3}+r_{1}-1 / 2}}{\left(1+\xi_{1}^{2}\right)^{r_{1}}\left(1+\left(\tau_{2}+\xi_{1}\right)^{2}\right)^{r_{2}}\left(1+\left(\tau_{3}+\xi_{1}\right)^{2}\right)^{r_{3}}} \mathrm{~d} \xi_{1}\right)^{-1}
\end{aligned}
$$

As a consequence we obtain

$$
f_{1}\left(-x_{2}-x_{3}, x_{2}, x_{3}\right)=\frac{1}{2 \pi} \int_{\mathbb{R}^{2}} \mathrm{e}^{\mathrm{i}\left(x_{2} \tau_{2}+x_{3} \tau_{3}\right.} \mathscr{F}_{2} h_{1}\left(\tau_{2}, \tau_{3}\right) \mathrm{d} \tau_{2} \mathrm{~d} \tau_{3}=h_{1}\left(x_{2}, x_{3}\right)
$$

as claimed. From the definition of the sets $\Omega_{i}$ we derive the existence of two positive constants $c_{1}$ and $c_{2}$ such that for all $\tau_{2}, \tau_{3}$

$$
c_{1} \leqslant H^{1}\left(\tau_{2}, \tau_{3}\right) \leqslant c_{2}
$$

as well as

$$
c_{1} \leqslant H^{2}\left(\tau_{2}, \tau_{3}\right) \leqslant c_{2} .
$$

This will be used to prove that $f_{1}$ is sufficiently regular. Indeed, we have

$$
\begin{aligned}
& \int_{\mathbb{R}^{3}}\left(1+\xi_{1}^{2}\right)^{r_{1}}\left(1+\xi_{2}^{2}\right)^{r_{2}}\left(1+\xi_{3}^{2}\right)^{r_{3}}\left|\mathscr{F}_{3} f_{1}(\xi)\right|^{2} \mathrm{~d} \vec{\xi} \\
& =\int_{\mathbb{R}^{3}}\left(1+\xi_{1}^{2}\right)^{-r_{1}}\left(1+\xi_{2}^{2}\right)^{-r_{2}}\left(1+\xi_{3}^{2}\right)^{-r_{3}}\left|\mathscr{F}_{2} h_{1}\left(\xi_{2}-\xi_{1}, \xi_{3}-\xi_{1}\right)\right|^{2} \\
& \quad \times\left(\mathscr{X}_{1}(\xi)\left|H^{1}\left(\xi_{2}-\xi_{1}, \xi_{3}-\xi_{1}\right)\right|^{2}\left(1+\left(\xi_{2}-\xi_{1}\right)^{2}\right)^{2 r_{2}+2 r_{1}-1}\left(1+\left(\xi_{3}-\xi_{1}\right)^{2}\right)^{2 r_{3}}\right. \\
& \left.\quad+\mathscr{X}_{2}(\xi)\left|H^{2}\left(\xi_{2}-\xi_{1}, \xi_{3}-\xi_{1}\right)\right|^{2}\left(1+\left(\xi_{2}-\xi_{1}\right)^{2}\right)^{2 r_{2}}\left(1+\left(\xi_{3}-\xi_{1}\right)^{2}\right)^{2 r_{3}+2 r_{1}-1}\right) \mathrm{d} \vec{\xi} \\
& =: J_{1}+J_{2} .
\end{aligned}
$$

A change of coordinates, the boundedness of $H^{1}$ and the definition of $\Omega_{1}$ yield

$$
\begin{aligned}
& J_{1} \leqslant c_{2}^{2} \\
& \int_{\left|\tau_{2}\right| \leqslant\left|\tau_{3}\right|}\left|\mathscr{F}_{2} h_{1}\left(\tau_{2}, \tau_{3}\right)\right|^{2}\left(1+\tau_{2}^{2}\right)^{2 r_{2}+2 r_{1}-1}\left(1+\tau_{3}^{2}\right)^{2 r_{3}} \\
& \times \int_{I\left(\tau_{2}\right)}\left(1+\xi_{1}^{2}\right)^{-r_{1}}\left(1+\left(\tau_{2}+\xi_{1}\right)^{2}\right)^{-r_{2}}\left(1+\left(\tau_{3}+\xi_{1}\right)^{2}\right)^{-r_{3}} \mathrm{~d} \xi_{1} \mathrm{~d} \tau_{2} \mathrm{~d} \tau_{3} \\
& \leqslant c_{3} \int_{\mathbb{R}^{2}}\left|\mathscr{F}_{2} h_{1}\left(\tau_{2}, \tau_{3}\right)\right|^{2}\left(1+\tau_{2}^{2}\right)^{r_{2}+r_{1}-1 / 2}\left(1+\tau_{3}^{2}\right)^{r_{3}} \mathrm{~d} \tau_{2} \mathrm{~d} \tau_{3} .
\end{aligned}
$$

The estimate of $J_{2}$ is similar. Hence

$$
\left\|f_{1} \mid S_{2}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right)\right\| \leqslant c_{4}\left(\left\|h_{1}\left|S_{2}^{r_{2}, r_{3}+r_{1}-\frac{1}{2}} W\left(\mathbb{R}^{2}\right)\|+\| h_{1}\right| S_{2}^{r_{2}+r_{1}-\frac{1}{2}, r_{3}} W\left(\mathbb{R}^{2}\right)\right\|\right)
$$

with some constant $c_{4}$ independent of $h_{1}$. This proves the boundedness of the extension.

Remark 2.10. Let us mention that we have not shown the existence of a linear and continuous extension operator. The step in which $g$ is split into the three functions $g_{1}, g_{2}$ and $g_{3}$ need not be linear. This problem will be investigated in the next subsection.

### 2.4.2. A description of the trace classes on the Fourier side.

For simplicity we concentrate on the situation $\min \left(r_{1}, r_{2}, r_{3}\right)>1 / 2$. The sum $S_{2}^{r_{1}, r_{2}} W\left(\mathscr{R}_{3}^{-1}, \mathbb{R}^{2}\right)+S_{2}^{r_{1}, r_{3}} W\left(\mathscr{R}_{2}^{-1}, \mathbb{R}^{2}\right)+S_{2}^{r_{2}, r_{3}} W\left(\mathscr{R}_{1}^{-1}, \mathbb{R}^{2}\right)$ is not direct. It is obvious that

$$
C_{0}^{\infty}\left(\mathbb{R}^{3}\right) \subset\left(S_{2}^{r_{1}, r_{2}} W\left(\mathscr{R}_{3}^{-1}, \mathbb{R}^{2}\right) \cap S_{2}^{r_{1}, r_{3}} W\left(\mathscr{R}_{2}^{-1}, \mathbb{R}^{2}\right) \cap S_{2}^{r_{2}, r_{3}} W\left(\mathscr{R}_{1}^{-1}, \mathbb{R}^{2}\right)\right) .
$$

At this moment it is not clear whether the connection between $g$ and its optimal decomposition $g_{1}+g_{2}+g_{3}$ can be realized in a linear way. But that can be seen easily by the Fourier-analytic description of the trace space.

Let $\mathscr{O}$ be an orthogonal basis of $\Gamma=\left\{\left(x_{1}, x_{2}, x_{3}\right) \in \mathbb{R}^{3}: x_{1}+x_{2}+x_{3}=0\right\}$ and let $\mathscr{R}_{1}, \mathscr{R}_{2}, \mathscr{R}_{3}$ be the matrices associated with $\mathscr{O}$. First, we notice that $g_{3} \in$ $S_{2}^{r_{1}, r_{2}} W\left(\mathscr{R}_{3}^{-1}, \mathbb{R}^{2}\right)$ if, and only if,

$$
\begin{equation*}
\underbrace{\left[1+\left(\sigma_{2,2} \xi_{1}-\sigma_{1,2} \xi_{2}\right)^{2}\right]^{r_{1} / 2}\left[1+\left(\sigma_{2,1} \xi_{1}-\sigma_{1,1} \xi_{2}\right)^{2}\right]^{r_{2} / 2}}_{m_{3}\left(\xi_{1}, \xi_{2}\right)} \mathscr{F} g_{3}\left(\xi_{1}, \xi_{2}\right) \in L_{2}\left(\mathbb{R}^{2}\right) \tag{2.18}
\end{equation*}
$$

cf. Lemma 2.2(iii). Similarly, $g_{1} \in S_{2}^{r_{2}, r_{3}} W\left(\mathscr{R}_{1}^{-1}, \mathbb{R}^{2}\right)$ if, and only if,

$$
\begin{equation*}
\underbrace{\left[1+\left(\sigma_{2,3} \xi_{1}-\sigma_{1,3} \xi_{2}\right)^{2}\right]^{r_{2} / 2}\left[1+\left(\sigma_{2,2} \xi_{1}-\sigma_{1,2} \xi_{2}\right)^{2}\right]^{r_{3} / 2}}_{m_{1}\left(\xi_{1}, \xi_{2}\right)} \mathscr{F} g_{1}\left(\xi_{1}, \xi_{2}\right) \in L_{2}\left(\mathbb{R}^{2}\right) \tag{2.19}
\end{equation*}
$$

and $g_{2} \in S_{2}^{r_{1}, r_{3}} W\left(\mathscr{R}_{2}^{-1}, \mathbb{R}^{2}\right)$ if, and only if,

$$
\begin{equation*}
\underbrace{\left[1+\left(\sigma_{2,3} \xi_{1}-\sigma_{1,3} \xi_{2}\right)^{2}\right]^{r_{1} / 2}\left[1+\left(\sigma_{2,1} \xi_{1}-\sigma_{1,1} \xi_{2}\right)^{2}\right]^{r_{3} / 2}}_{m_{2}\left(\xi_{1}, \xi_{2}\right)} \mathscr{F} g_{2}\left(\xi_{1}, \xi_{2}\right) \in L_{2}\left(\mathbb{R}^{2}\right) \tag{2.20}
\end{equation*}
$$

In view of these characterizations we define

$$
\begin{equation*}
m\left(\xi_{1}, \xi_{2}\right):=\min \left(m_{1}\left(\xi_{1}, \xi_{2}\right), m_{2}\left(\xi_{1}, \xi_{2}\right), m_{3}\left(\xi_{1}, \xi_{2}\right)\right) \tag{2.21}
\end{equation*}
$$

and

$$
\begin{equation*}
L_{2}\left(\mathbb{R}^{2}, m\right):=\left\{g \in L_{2}\left(\mathbb{R}^{2}\right): m \mathscr{F} g \in L_{2}\left(\mathbb{R}^{2}\right)\right\} \tag{2.22}
\end{equation*}
$$

equipped with the natural norm

$$
\left\|g\left|L_{2}\left(\mathbb{R}^{2}, m\right)\|:=\| m \mathscr{F} g\right| L_{2}\left(\mathbb{R}^{2}\right)\right\|
$$

Now we arrive at the main result of this section.

Theorem 2.11. Let $\mathscr{O}$ be an orthogonal basis of $\Gamma$ and let $\mathscr{R}_{i}, i=1,2,3$, be matrices associated with $\mathscr{O}$ by (1.1), (2.8) and (2.10). Suppose (2.15) and $r_{i} \neq \frac{1}{2}$, $i=1,2,3$. Then there exists a continuous function $m$ such that $\operatorname{tr}_{\mathscr{O}}$ becomes a retraction of $S_{2}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right)$ onto $L_{2}\left(\mathbb{R}^{2}, m\right)$. There is a bounded linear extension operator ext $\in \mathscr{L}\left(L_{2}\left(\mathbb{R}^{2}, m\right), S_{2}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right)\right.$ ) such that $\operatorname{tr}_{\mathscr{O}} \circ$ ext $=I$ (identity on $L_{2}\left(\mathbb{R}^{2}, m\right)$ ).

Proof. We concentrate on the case $\min \left(r_{1}, r_{2}, r_{3}\right)>\frac{1}{2}$. Then the function $m$ is given by (2.21). The modifications which have to be made for the general situation are obvious. We omit the details.

Step 1. Boundedness. Again we shall use the abbreviations $A_{1}=S_{2}^{r_{2}, r_{3}} W\left(\mathscr{R}_{1}^{-1}\right.$, $\left.\mathbb{R}^{2}\right), A_{2}=S_{2}^{r_{1}, r_{3}} W\left(\mathscr{R}_{2}^{-1}, \mathbb{R}^{2}\right)$ and $A_{3}=S_{2}^{r_{1}, r_{2}} W\left(\mathscr{R}_{3}^{-1}, \mathbb{R}^{2}\right)$. Let $g \in A_{1}+A_{2}+A_{3}$ and let $g=g_{1}+g_{2}+g_{3}$ be an optimal decomposition of $g$ with $g_{i} \in A_{i}$. Then

$$
m(\xi)|\mathscr{F} g(\xi)| \leqslant \sum_{i=1}^{3} m_{i}(\xi)\left|\mathscr{F} g_{i}(\xi)\right|, \quad \xi \in \mathbb{R}^{2}
$$

But this implies

$$
\left\|g\left|L_{2}\left(\mathbb{R}^{2}, m\right)\left\|\leqslant \sum_{i=1}^{3}\right\| m_{i} \mathscr{F} g_{i}\right| L_{2}\left(\mathbb{R}^{2}\right)\right\| \leqslant c \sum_{i=1}^{3}\left\|g_{i} \mid A_{i}\right\|,
$$

with some $c$ independent of $g$.

Vice versa, if $g \in L_{2}\left(\mathbb{R}^{2}, m\right)$, then we define

$$
\begin{equation*}
\Omega_{i}:=\left\{\left(\xi_{1}, \xi_{2}\right): m_{i}\left(\xi_{1}, \xi_{2}\right)=m\left(\xi_{1}, \xi_{2}\right)\right\} \tag{2.23}
\end{equation*}
$$

$\mathscr{X}_{i}$ denotes its characteristic function, and

$$
\begin{equation*}
g_{i}(x):=\mathscr{F}^{-1}\left[\mathscr{X}_{i}(\xi) \mathscr{F} g(\xi)\right](x), \quad i=1,2,3 . \tag{2.24}
\end{equation*}
$$

Thanks to

$$
\left|\Omega_{i} \cap \Omega_{j}\right|=0, \quad i \neq j, \quad \text { and } \quad \bigcup_{i=1}^{3} \Omega_{i}=\mathbb{R}^{2}
$$

$\left(|\cdot|\right.$ being the Lebesgue measure in $\left.\mathbb{R}^{2}\right)$ this implies $g=g_{1}+g_{2}+g_{3}$ and

$$
\left\|m_{i} \mathscr{F} g_{i}\left|L_{2}\left(\mathbb{R}^{2}\right)\|\leqslant\| m \mathscr{F} g\right| L_{2}\left(\mathbb{R}^{2}\right)\right\|, \quad i=1,2,3
$$

Summarizing we have proved the coincidence of $S_{2}^{r_{1}, r_{2}} W\left(\mathscr{R}_{3}^{-1}, \mathbb{R}^{2}\right)+S_{2}^{r_{1}, r_{3}} W\left(\mathscr{R}_{2}^{-1}\right.$, $\left.\mathbb{R}^{2}\right)+S_{2}^{r_{2}, r_{3}} W\left(\mathscr{R}_{1}^{-1}, \mathbb{R}^{2}\right)$ and $L_{2}\left(\mathbb{R}^{2}, m\right)$ in the sense of equivalent norms. Hence $\operatorname{tr}_{\mathscr{O}} \in \mathscr{L}\left(S_{2}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right), L_{2}\left(\mathbb{R}^{2}, m\right)\right)$.

Step 2. The linear extension. Since the mappings $g \rightarrow g_{i}, i=1,2,3$, cf. (2.24), are linear and continuous, the extension operator constructed in the proof of Theorem 2.9 is linear and bounded as well.

### 2.4.3. The trace space for a dominating direction.

This subsection contains an additional observation of minor importance. So we concentrate on $\min \left(r_{1}, r_{2}, r_{3}\right)>\frac{1}{2}$.

A simplified description of the trace spaces can be given if one of the parameters $r_{1}, r_{2}, r_{3}$ is dominating the sum of the others.

Lemma 2.12. Let $\mathscr{O}$ be an orthogonal basis of $\Gamma$ and let $\mathscr{R}_{i}, i=1,2,3$, be matrices associated with $\mathscr{O}$ by (1.1), (2.8) and (2.10). Then the embeddings

$$
\begin{aligned}
& S_{2}^{r_{1}, r_{3}} W\left(\mathscr{R}_{2}^{-1}, \mathbb{R}^{2}\right) \hookrightarrow S_{2}^{r_{1}, r_{2}} W\left(\mathscr{R}_{3}^{-1}, \mathbb{R}^{2}\right) \quad \text { and } \\
& S_{2}^{r_{2}, r_{3}} W\left(\mathscr{R}_{1}^{-1}, \mathbb{R}^{2}\right) \hookrightarrow S_{2}^{r_{1}, r_{2}} W\left(\mathscr{R}_{3}^{-1}, \mathbb{R}^{2}\right)
\end{aligned}
$$

exist if, and only if, $r_{3} \geqslant r_{1}+r_{2}$.
Proof. Again we work in the Fourier image. Let $m_{1}, m_{2}$ and $m_{3}$ be the functions defined in (2.18)-(2.20). Then the first embedding is equivalent to the boundedness of $m_{3} / m_{2}$ and the second is equivalent to the boundedness of $m_{3} / m_{1}$, respectively.

Let us turn to the boundedness of the first quotient. By a change of coordinates

$$
y_{1}:=\sigma_{2,2} \xi_{1}-\sigma_{1,2} \xi_{2} \quad \text { and } \quad y_{2}:=\sigma_{2,3} \xi_{1}-\sigma_{1,3} \xi_{2}
$$

and taking care of $\vec{\sigma}_{1}, \vec{\sigma}_{2} \in \Gamma$ the boundedness of $m_{3} / m_{2}$ is equivalent to

$$
\sup _{y_{1}, y_{2} \in \mathbb{R}} \frac{\left(1+y_{1}^{2}\right)^{r_{1}}\left[1+\left(y_{1}+y_{2}\right)^{2}\right]^{r_{2}}}{\left(1+y_{1}^{2}\right)^{r_{3}}\left(1+y_{2}^{2}\right)^{r_{2}}}<\infty .
$$

With $y_{2}=0$ the necessity of $r_{3} \geqslant r_{1}+r_{2}$ follows. Sufficiency can be derived from

$$
1+\left(y_{1}+y_{2}\right)^{2} \leqslant 2\left(1+y_{1}^{2}\right)\left(1+y_{2}^{2}\right)
$$

Theorem 2.13. Let $\mathscr{O}$ be an orthogonal basis of $\Gamma=\left\{\left(x_{1}, x_{2}, x_{3}\right) \in \mathbb{R}^{3}: x_{1}+\right.$ $\left.x_{2}+x_{3}=0\right\}$ and let $\mathscr{R}_{3}$ be the matrix associated with $\mathscr{O}$ by (1.1) and (2.10). Let $\min \left(r_{1}, r_{2}, r_{3}\right)>1 / 2$ and suppose $r_{3} \geqslant r_{1}+r_{2}$. Then $\operatorname{tr}_{\mathscr{O}}$ becomes a retraction of $S_{2}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right)$ onto $\left.S_{2}^{r_{1}, r_{2}} W\left(\mathscr{R}_{3}^{-1}, \mathbb{R}^{2}\right)\right)$ and

$$
\left.S_{2}^{r_{1}, r_{2}} W\left(\mathscr{R}_{3}^{-1}, \mathbb{R}^{2}\right)\right)=L_{2}\left(\mathbb{R}^{2}, m_{3}\right) \quad \text { (equivalent norms). }
$$

Proof. From Lemma 2.12 we derive

$$
S_{2}^{r_{1}, r_{2}} W\left(\mathscr{R}_{3}, \mathbb{R}^{2}\right)+S_{2}^{r_{2}, r_{3}} W\left(\mathscr{R}_{1}, \mathbb{R}^{2}\right)+S_{2}^{r_{1}, r_{3}} W\left(\mathscr{R}_{2}, \mathbb{R}^{2}\right)=S_{2}^{r_{1}, r_{2}} W\left(\mathscr{R}_{3}, \mathbb{R}^{2}\right)
$$

with equivalent norms. Now the statement follows from Theorems 2.9 and 2.11. The last identity has been derived in (2.18).

Also $\operatorname{tr}_{1}, \operatorname{tr}_{2}$ and $\operatorname{tr}_{3}$ have additional properties if one of the smoothness parameters dominates the sum of the others.

Theorem 2.14. Let $\mathscr{O}$ be an orthogonal basis of $\Gamma=\left\{\left(x_{1}, x_{2}, x_{3}\right) \in \mathbb{R}^{3}: x_{1}+\right.$ $\left.x_{2}+x_{3}=0\right\}$ and let $\mathscr{R}_{3}$ be the matrix associated with $\mathscr{O}$ by (1.1) and (2.10). Let $\min \left(r_{1}, r_{2}, r_{3}\right)>1 / 2$ and suppose $r_{3} \geqslant r_{1}+r_{2}$.

Then $\operatorname{tr}_{3}$ becomes a retraction of $S_{2}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right)$ onto $S_{2}^{r_{1}, r_{2}} W\left(\mathbb{R}^{2}\right)$, i.e. there exists a linear extension operator ext* $\in \mathscr{L}\left(S_{2}^{r_{1}, r_{2}} W\left(\mathbb{R}^{2}\right), S_{2}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right)\right)$ s.t. $\operatorname{tr}_{3} \circ \mathrm{ext}^{*}=I$.

Proof. Step 1. Boundedness of $\operatorname{tr}_{3}$. To show that, we use again (2.12). Furthermore, taking $h\left(\xi_{1}, \xi_{2}, \tau_{3}\right):=\mathscr{F} f\left(\xi_{1}+\tau_{3}, \xi_{2}+\tau_{3}, \tau_{3}\right)$, it will be enough to show the existence of some positive constant $c$ such that

$$
\begin{align*}
& \int_{\mathbb{R}^{2}}\left(1+y_{1}^{2}\right)^{r_{1}}\left(1+y_{2}^{2}\right)^{r_{2}}\left|\int_{\mathbb{R}} h\left(y_{1}, y_{2}, y_{3}\right) \mathrm{d} y_{3}\right|^{2} \mathrm{~d} y_{1} \mathrm{~d} y_{2}  \tag{2.25}\\
& \leqslant c \int_{\mathbb{R}^{3}}\left[1+\left(y_{1}+y_{3}\right)^{2}\right]^{r_{1}}\left[1+\left(y_{2}+y_{3}\right)^{2}\right]^{r_{2}}\left[1+y_{3}^{2}\right]^{r_{3}}\left|h\left(y_{1}, y_{2}, y_{3}\right)\right|^{2} \mathrm{~d} y_{1} \mathrm{~d} y_{2} \mathrm{~d} y_{3}
\end{align*}
$$

Let us denote

$$
\Theta_{1}\left(y_{1}, y_{2}\right):=\left(1+y_{1}^{2}\right)^{r_{1}}\left(1+y_{2}^{2}\right)^{r_{2}}
$$

and

$$
\Theta_{2}\left(y_{1}, y_{2}, y_{3}\right):=\left[1+\left(y_{1}+y_{3}\right)^{2}\right]^{r_{1}}\left[1+\left(y_{2}+y_{3}\right)^{2}\right]^{r_{2}}\left(1+y_{3}^{2}\right)^{r_{3}}
$$

respectively. Then Hölder's inequality leads to

$$
\begin{aligned}
\left(1+y_{1}^{2}\right)^{r_{1}} & \left(1+y_{2}^{2}\right)^{r_{2}}\left(\int_{\mathbb{R}}\left|h\left(y_{1}, y_{2}, y_{3}\right)\right| \mathrm{d} y_{3}\right)^{2} \\
& =\left(\int_{\mathbb{R}} \frac{\sqrt{\Theta_{1}\left(y_{1}, y_{2}\right)}}{\sqrt{\Theta_{2}\left(y_{1}, y_{2}, y_{3}\right)}} \sqrt{\Theta_{2}\left(y_{1}, y_{2}, y_{3}\right)}\left|h\left(y_{1}, y_{2}, y_{3}\right)\right| \mathrm{d} y_{3}\right)^{2} \\
& \leqslant \underbrace{\left(\sup _{y_{1}, y_{2} \in \mathbb{R}} \int_{\mathbb{R}} \frac{\Theta_{1}\left(y_{1}, y_{2}\right)}{\Theta_{2}\left(y_{1}, y_{2}, y_{3}\right)} \mathrm{d} y_{3}\right)} \int_{\mathbb{R}} \Theta_{2}\left(y_{1}, y_{2}, y_{3}\right)\left|h\left(y_{1}, y_{2}, y_{3}\right)\right|^{2} \mathrm{~d} y_{3} . \\
& :=\Theta\left(r_{1}, r_{2}, r_{3}\right)
\end{aligned}
$$

If $\Theta\left(r_{1}, r_{2}, r_{3}\right)<\infty$, then it is enough to integrate this inequality with respect to $y_{1}, y_{2} \in \mathbb{R}$ to obtain (2.25). To prove finiteness of $\Theta\left(r_{1}, r_{2}, r_{3}\right)$ under the given restrictions is elementary.

Step 2. Surjectivity of $\operatorname{tr}_{3}$. Here we make use of the operator ext ${ }_{3}^{*}$, defined in the proof of Theorem 2.9, Substep 3.1.

Remark 2.15. By symmetry we have similar statements with respect to $\operatorname{tr}_{1}$ as well as to $\operatorname{tr}_{2}$, e.g. if $\min \left(r_{1}, r_{2}, r_{3}\right)>\frac{1}{2}$ and $r_{2} \geqslant r_{1}+r_{3}$ then $\operatorname{tr}_{2}$ becomes a retraction of $S_{2}^{r_{1}, r_{2}, r_{3}} W\left(\mathbb{R}^{3}\right)$ onto $S_{2}^{r_{1}, r_{3}} W\left(\mathbb{R}^{2}\right)$.

### 2.4.4. An example.

We consider the orthogonal basis $\vec{\sigma}_{1}:=(1,-1,0)$, and $\vec{\sigma}_{2}:=(1,1,-2)$ of $\Gamma$. Then the functions $m_{i}, i=1,2,3$, defined in (2.18)-(2.20), are given by

$$
\begin{aligned}
& m_{1}^{2}\left(\xi_{1}, \xi_{2}\right)=\left[1+\left(2 \xi_{1}\right)^{2}\right]^{r_{2}}\left[1+\left(\xi_{1}+\xi_{2}\right)^{2}\right]^{r_{3}} \\
& m_{2}^{2}\left(\xi_{1}, \xi_{2}\right)=\left[1+\left(2 \xi_{1}\right)^{2}\right]^{r_{1}}\left[1+\left(\xi_{1}-\xi_{2}\right)^{2}\right]^{r_{3}} \\
& m_{3}^{2}\left(\xi_{1}, \xi_{2}\right)=\left[1+\left(\xi_{1}+\xi_{2}\right)^{2}\right]^{r_{1}}\left[1+\left(\xi_{1}-\xi_{2}\right)^{2}\right]^{r_{2}}
\end{aligned}
$$

Let $r_{1}=r_{2}=r_{3}=1$ and define

$$
\begin{gathered}
w\left(\xi_{1}, \xi_{2}\right):=\min \left(1+5 \xi_{1}^{2}+\xi_{2}^{2}+2 \xi_{1} \xi_{2}+4 \xi_{1}^{4}+4 \xi_{1}^{2} \xi_{2}^{2}+8 \xi_{1}^{3} \xi_{2}\right. \\
1+5 \xi_{1}^{2}+\xi_{2}^{2}-2 \xi_{1} \xi_{2}+4 \xi_{1}^{4}+4 \xi_{1}^{2} \xi_{2}^{2}-8 \xi_{1}^{3} \xi_{2} \\
\left.1+2 \xi_{1}^{2}+2 \xi_{2}^{2}+\xi_{1}^{4}-2 \xi_{1}^{2} \xi_{2}^{2}+\xi_{2}^{4}\right)
\end{gathered}
$$

cf. (2.21). Hence, the trace space of the Sobolev space $S_{2}^{1,1,1} W\left(\mathbb{R}^{3}\right)$ with respect to this orthogonal basis is the collection of all functions $g \in L_{2}\left(\mathbb{R}^{2}\right)$ such that

$$
\int_{\mathbb{R}^{2}} w\left(\xi_{1}, \xi_{2}\right)\left|\mathscr{F} g\left(\xi_{1}, \xi_{2}\right)\right|^{2} \mathrm{~d} \xi<\infty
$$

Furthermore, the trace space of the Sobolev space $S_{2}^{1,1,2} W\left(\mathbb{R}^{3}\right)$ with respect to this orthogonal basis is the collection of all functions $g \in L_{2}\left(\mathbb{R}^{2}\right)$ such that

$$
\int_{\mathbb{R}^{2}}\left(1+2 \xi_{1}^{2}+2 \xi_{2}^{2}+\xi_{1}^{4}-2 \xi_{1}^{2} \xi_{2}^{2}+\xi_{2}^{4}\right)\left|\mathscr{F} g\left(\xi_{1}, \xi_{2}\right)\right|^{2} \mathrm{~d} \xi<\infty
$$

## 3. Besov and Triebel-Lizorkin spaces

Now we turn to the general case of Besov and Triebel-Lizorkin spaces. To begin with we recall the Fourier-analytic definition as well as the characterization by atoms of these classes. Since we shall need the spaces for $d=3$ and for $d=2$ we shall work for a while with the general $d$-dimensional case.

### 3.1. Notation.

As usual, $\mathbb{R}^{d}$ denotes the $d$-dimensional real Euclidean space, $\mathbb{N}$ the collection of all natural numbers and $\mathbb{N}_{0}=\mathbb{N} \cup\{0\}$. The letter $\mathbb{Z}$ stands for the set of all integers and $\mathbb{C}$ denotes the complex numbers.

If $x, y \in \mathbb{R}^{d}$, we write $x>y$ if, and only if, $x_{i}>y_{i}$ for every $i=1, \ldots, d$. Similarly, we define the relations $x \geqslant y, x<y, x \leqslant y$. Finally, in slight abuse of notation, we write $x>\lambda$ for $x \in \mathbb{R}^{d}, \lambda \in \mathbb{R}$ if $x_{i}>\lambda, i=1, \ldots, d$. For a real number $x$ we denote by $x_{+}:=\max (x, 0)$ the nonnegative part.

Let $S\left(\mathbb{R}^{d}\right)$ be the Schwartz space of all complex-valued rapidly decreasing infinitely differentiable functions on $\mathbb{R}^{d}$.

### 3.2. The Fourier-analytic approach.

Let $\varphi \in \mathscr{S}(\mathbb{R})$ with

$$
\begin{equation*}
\varphi(t)=1 \quad \text { if }|t| \leqslant 1 \quad \text { and } \quad \varphi(t)=0 \quad \text { if }|t| \geqslant \frac{3}{2} \tag{3.1}
\end{equation*}
$$

We put $\varphi_{0}=\varphi, \varphi_{1}(t)=\varphi(t / 2)-\varphi(t)$ and

$$
\varphi_{j}(t):=\varphi_{1}\left(2^{-j+1} t\right), \quad t \in \mathbb{R}, \quad j \in \mathbb{N}
$$

Hence we have $\sum_{j=0}^{\infty} \varphi_{j}(t)=1$ for all $t \in \mathbb{R}$. For $\bar{k}=\left(k_{1}, \ldots, k_{d}\right) \in \mathbb{N}_{0}^{d}$ and $x=$ $\left(x_{1}, \ldots, x_{d}\right) \in \mathbb{R}^{d}$ we define $\varphi_{\bar{k}}(x):=\varphi_{k_{1}}\left(x_{1}\right) \cdot \ldots \cdot \varphi_{k_{d}}\left(x_{d}\right)$. Then, since

$$
\begin{equation*}
\sum_{\bar{k} \in \mathbb{N}_{0}^{d}} \varphi_{\bar{k}}(x)=1 \quad \text { for every } x \in \mathbb{R}^{d}, \tag{3.2}
\end{equation*}
$$

the system $\left\{\varphi_{\bar{k}}\right\}_{\bar{k} \in \mathbb{N}_{0}^{d}}$ forms a smooth dyadic resolution of unity. This will be used to define the classes of functions we are interested in.

Definition 3.1. Let $\bar{r}=\left(r_{1}, \ldots, r_{d}\right) \in \mathbb{R}^{d}$, and $0<q \leqslant \infty$.
(i) Let $0<p \leqslant \infty$. Then the Besov space of dominating mixed smoothness $S_{p, q}^{\bar{r}} B\left(\mathbb{R}^{d}\right)$ is the collection of all $f \in \mathscr{S}^{\prime}\left(\mathbb{R}^{d}\right)$ such that

$$
\begin{align*}
\left\|f \mid S_{p, q}^{\bar{r}} B\left(\mathbb{R}^{d}\right)\right\|_{\varphi} & =\left(\sum_{\bar{k} \in \mathbb{N}_{0}^{d}} 2^{q \overline{\bar{k}} \cdot \bar{r}}\left\|\mathscr{F}^{-1}\left[\varphi_{\bar{k}} \mathscr{F} f\right] \mid L_{p}\left(\mathbb{R}^{d}\right)\right\|^{q}\right)^{1 / q}  \tag{3.3}\\
& =\left\|2^{\bar{k} \cdot \bar{r}} \mathscr{F}^{-1}\left[\varphi_{\bar{k}} \mathscr{F} f\right] \mid \ell_{q}\left(L_{p}\right)\right\|
\end{align*}
$$

is finite.
(ii) Let $0<p<\infty$. Then the Triebel-Lizorkin space of dominating mixed smoothness $S_{p, q}^{\bar{r}} F\left(\mathbb{R}^{d}\right)$ is the collection of all $f \in \mathscr{S}^{\prime}\left(\mathbb{R}^{d}\right)$ such that

$$
\begin{align*}
\left\|f \mid S_{p, q}^{\bar{r}} F\left(\mathbb{R}^{d}\right)\right\|_{\varphi} & =\left\|\left(\sum_{\bar{k} \in \mathbb{N}_{0}^{d}}\left|2^{\bar{k} \cdot \bar{r}} \mathscr{F}^{-1}\left[\varphi_{\bar{k}} \mathscr{F} f\right](\cdot)\right|^{q}\right)^{1 / q} \mid L_{p}\left(\mathbb{R}^{d}\right)\right\|  \tag{3.4}\\
& =\| 2^{\bar{k} \cdot \bar{r} \mathscr{F}^{-1}\left[\varphi_{\bar{k}} \mathscr{F} f\right] \mid L_{p}\left(\ell_{q}\right) \|}
\end{align*}
$$

is finite.
Remark 3.2. 1. Sometimes, we shall write $S_{p, q}^{\bar{r}} A\left(\mathbb{R}^{d}\right)$ instead of $S_{p, q}^{\bar{r}} B\left(\mathbb{R}^{d}\right)$ or $S_{p, q}^{\bar{r}} F\left(\mathbb{R}^{d}\right)$.
2. Different functions $\varphi$ (with properties described above) lead to equivalent quasinorms on $S_{p, q}^{\bar{r}} A\left(\mathbb{R}^{d}\right)$. We shall write $\left\|f \mid S_{p, q}^{\bar{r}} A\left(\mathbb{R}^{d}\right)\right\|$ meaning one of these quasinorms (which one is in general of no importance in our context). For details see [18, Section 2.2.3].
3. For a systematic investigation of these classes we refer to the monographs [2] and [18]. More recent developments may be found in [3], [11] and [23, 24, 25].
4. For $1<p<\infty$ we have the coincidence of $S_{p, 2}^{\bar{r}} F\left(\mathbb{R}^{d}\right)$ and the Sobolev space $S_{p}^{\bar{r}} W\left(\mathbb{R}^{d}\right)$ in the sense of equivalent norms, cf. [12] and [18, 2.3.1].

### 3.3. Atomic decomposition.

In the mid-eighties Frazier and Jawerth [5] have been the first who studied atomic decompositions of Besov spaces. One of the applications has been a description of the solution of the trace problem with respect to hyperplanes in the isotropic situation. Here we follow the same philosophy. We shall make use of the characterization of Besov and Lizorkin-Triebel spaces by means of atoms for studying the properties of $\operatorname{tr}_{\mathscr{O}}$.

Atomic decomposition techniques allow a certain discretization. Function spaces are replaced by sequence spaces. This method has been studied in various situations by now, cf. [5], [6], [1], [21] for isotropic spaces of Besov and Lizorkin-Triebel type and [10] for some generalizations in various directions. Besov and Lizorkin-Triebel spaces of dominating mixed smoothness have been characterized in such a way in [24].

### 3.3.1. Sequence spaces.

For $\bar{\nu} \in \mathbb{N}_{0}^{d}, \bar{m} \in \mathbb{Z}^{d}$ we denote by $Q_{\bar{\nu} \bar{m}}$ the cube with the centre at the point $2^{-\bar{\nu}} \bar{m}=\left(2^{-\nu_{1}} m_{1}, \ldots, 2^{-\nu_{d}} m_{d}\right)$, sides parallel to the coordinate axes and of lengths $2^{-\nu_{1}}, \ldots, 2^{-\nu_{d}}$. We denote by $\chi_{\bar{\nu} \bar{m}}=\chi_{Q_{\bar{\nu} \bar{m}}}$ the characteristic function of $Q_{\bar{\nu} \bar{m}}$ and by $c Q_{\bar{\nu} \bar{m}}$ we mean a cube concentric with $Q_{\bar{\nu} \bar{m}}$ with sides $c$ times larger.

Definition 3.3. If $0<p, q \leqslant \infty, \bar{r} \in \mathbb{R}^{d}$ and

$$
\begin{equation*}
\lambda=\left\{\lambda_{\bar{\nu} \bar{m}} \in \mathbb{C}: \bar{\nu} \in \mathbb{N}_{0}^{d}, \bar{m} \in \mathbb{Z}^{d}\right\}, \tag{3.5}
\end{equation*}
$$

then we define

$$
\begin{equation*}
s_{p q}^{\bar{r}} b:=\left\{\lambda:\left\|\lambda \mid s_{p q}^{\bar{r}} b\right\|=\left(\sum_{\bar{\nu} \in \mathbb{N}_{0}^{d}} 2^{\overline{\bar{v}} \cdot\left(\bar{r}-\frac{1}{p}\right) q}\left(\sum_{\bar{m} \in \mathbb{Z}^{d}}\left|\lambda_{\bar{\nu} \bar{m}}\right|^{p}\right)^{q / p}\right)^{1 / q}<\infty\right\} \tag{3.6}
\end{equation*}
$$

and

$$
\begin{equation*}
s_{p q}^{\bar{r}} f:=\left\{\lambda:\left\|\lambda\left|s_{p q}^{\bar{r}} f\|=\|\left(\sum_{\bar{\nu} \in \mathbb{N}_{0}^{d}} \sum_{\bar{m} \in \mathbb{Z}^{d}}\left|2^{\bar{\nu} \cdot \bar{r}} \lambda_{\bar{\nu} \bar{m}} \chi_{\bar{\nu} \bar{m}}(\cdot)\right|^{q}\right)^{1 / q}\right| L_{p}\left(\mathbb{R}^{d}\right)\right\|<\infty\right\} \tag{3.7}
\end{equation*}
$$

with the usual modification for $p$ and/or $q$ equal to $\infty$.
Remark 3.4. We shall use the same convention as in the case of the distribution spaces: from time to time we shall write $\left\|\lambda \mid s_{p q}^{\bar{r}} a\right\|$ instead of $\left\|\lambda \mid s_{p q}^{\bar{r}} b\right\|$ or $\left\|\lambda \mid s_{p q}^{\bar{r}} f\right\|$, respectively.

### 3.3.2. Atomic decompositions.

We will be very brief and refer for details to [23] and [24]. Here we concentrate on the "regular" case, i.e.

$$
\bar{r}> \begin{cases}\sigma_{p}=\max (1 / p-1,0) & \text { in the } B \text {-case }  \tag{3.8}\\ \sigma_{p q}=\max (1 / \min (p, q)-1,0) & \text { in the } F \text {-case }\end{cases}
$$

The phrase "regular" indicates that only those distribution spaces are considered which consist of regular distributions. Then, compared with the general case, no moment conditions have to be satisfied by the elementary building blocks called atoms. As usual, $[x]$ denotes the integer part of the real number $x$. If $Q$ is a cube and $\delta$ is a positive real number then $\delta Q$ denotes the cube with the same center as $Q$, sides parallel to those of $Q$ and sidelength multiplied by $\delta$.

Definition 3.5. Let $\bar{K}=\left(K_{1}, \ldots, K_{d}\right) \in \mathbb{N}_{0}^{d}$ and $\delta>1$. A $\bar{K}$-times differentiable complex-valued function $a(x)$ is called $\bar{K}$-atom related to $Q_{\bar{\nu} \bar{m}}$ if

$$
\begin{equation*}
\operatorname{supp} a \subset \delta Q_{\bar{\nu} \bar{m}} \tag{3.9}
\end{equation*}
$$

and

$$
\begin{equation*}
\sup _{x \in \mathbb{R}^{d}}\left|D^{\alpha} a(x)\right| \leqslant 2^{\alpha \cdot \bar{\nu}} \quad \text { for } 0 \leqslant \alpha \leqslant \bar{K} \tag{3.10}
\end{equation*}
$$

Theorem 3.6. Let $0<p, q \leqslant \infty\left(p<\infty\right.$ in the $F$-case), and $\bar{r} \in \mathbb{R}^{d}$ with (3.8). Fix $\bar{K} \in \mathbb{N}_{0}^{d}$ with

$$
\begin{equation*}
K_{i} \geqslant\left(1+\left[r_{i}\right]\right)_{+} \quad i=1, \ldots, d, \tag{3.11}
\end{equation*}
$$

and $\delta$ sufficiently large.
Then $f \in \mathscr{S}^{\prime}\left(\mathbb{R}^{d}\right)$ belongs to $S_{p, q}^{\bar{r}} A\left(\mathbb{R}^{d}\right)$ if, and only if, it can be represented as

$$
\begin{equation*}
f=\sum_{\bar{\nu} \in \mathbb{N}_{0}^{d}} \sum_{\bar{m} \in \mathbb{Z}^{d}} \lambda_{\bar{\nu} \bar{m}} a_{\bar{\nu} \bar{m}}(x) \quad\left(\text { convergence in } S^{\prime}\left(\mathbb{R}^{d}\right)\right) \tag{3.12}
\end{equation*}
$$

where $\left\{a_{\bar{\nu} \bar{m}}(x)\right\}_{\bar{\nu} \in \mathbb{N}_{0}^{d}, \bar{m} \in \mathbb{Z}^{d}}$ are $\bar{K}$-atoms related to $Q_{\bar{\nu} \bar{m}}$ and $\lambda \in s_{p q}^{\bar{r}} a$. Furthermore,

$$
\inf \left\|\lambda \mid s_{p q}^{\bar{r}} a\right\|,
$$

where the infimum is taken over all admissible representations (3.12), yields an equivalent quasi-norm in $S_{p, q}^{\bar{r}} A\left(\mathbb{R}^{d}\right)$.

Remark 3.7. To explain our philosophy, let the function $a$ be a $\bar{K}=\left(K_{1}, K_{2}\right.$, $\left.K_{3}\right)$-atom related to $Q_{\bar{\nu} \bar{m}}$, where $\nu=\left(\nu_{1}, \nu_{2}, \nu_{3}\right)$ and $m=\left(m_{1}, m_{2}, m_{3}\right)$. Then

$$
\left(\operatorname{tr}_{3} a\right)\left(x_{1}, x_{2}\right)=a\left(x_{1}, x_{2},-\left(x_{1}+x_{2}\right)\right)
$$

becomes a $\left(K_{1}, K_{2}\right)$-atom with respect to $Q_{\left(\nu_{1}, \nu_{2}\right),\left(m_{1}, m_{2}\right)}$ if $K_{3} \geqslant K_{1}+K_{2}$ and $\nu_{3} \leqslant \min \left(\nu_{1}, \nu_{2}\right)$. Similarly $\operatorname{tr}_{2} a\left(\operatorname{tr}_{1} a\right)$ becomes a $\left(K_{1}, K_{3}\right)$-atom $\left(\left(K_{2}, K_{3}\right)\right.$-atom $)$ with respect to $Q_{\left(\nu_{1}, \nu_{3}\right),\left(m_{1}, m_{3}\right)}\left(Q_{\left(\nu_{2}, \nu_{3}\right),\left(m_{2}, m_{3}\right)}\right)$ if $K_{2} \geqslant K_{1}+K_{3}\left(K_{1} \geqslant K_{2}+K_{3}\right)$ and $\nu_{2} \leqslant \min \left(\nu_{1}, \nu_{3}\right)\left(\nu_{1} \leqslant \min \left(\nu_{2}, \nu_{3}\right)\right)$. This simple observation will motivate an appropriate decomposition of the atomic decomposition of a function which turns out to be a basic step in our proof of the boundedness of $\operatorname{tr}_{\mathscr{O}}$.

### 3.4. Traces of Besov spaces of dominating mixed smoothness.

For a better comparison we recall the properties of the mapping $f\left(x_{1}, x_{2}, x_{3}\right) \mapsto$ $f\left(x_{1}, x_{2}, 0\right)$ in this general context, cf. e.g. Amanov [2, 9.5] and Schmeißer, Triebel [18, 2.4.2] (further references are given in [18, Remark 2.4.2]).

Proposition 3.8. Let $0<q \leqslant \infty$.
(i) Let $0<p \leqslant \infty$ and $r_{3}>1 / p$. Then the mapping

$$
T: f\left(x_{1}, x_{2}, x_{3}\right) \rightarrow f\left(x_{1}, x_{2}, 0\right)
$$

extends to a retraction from $S_{p, q}^{r_{1}, r_{2}, r_{3}} B\left(\mathbb{R}^{3}\right)$ onto $S_{p, q}^{r_{1}, r_{2}} B\left(\mathbb{R}^{2}\right)$.
(ii) Let $0<p<\infty$ and $r_{3}>1 / p$. Then the mapping $T$ extends to a retraction from $S_{p, q}^{r_{1}, r_{2}, r_{3}} F\left(\mathbb{R}^{3}\right)$ onto $S_{p, q}^{r_{1}, r_{2}} F\left(\mathbb{R}^{2}\right)$.

As mentioned in Introduction, to reflect the underlying geometry of our problem, we have to define some new spaces with dominating mixed smoothness, cf. Subsection 2.2 for $p=2$.

Definition 3.9. Let $0<q \leqslant \infty, 0<p \leqslant \infty$ in the B-case and $0<p<\infty$ in the $F$-case. Let $\mathscr{R}$ be a $(2,2)$-matrix with $\operatorname{det} \mathscr{R} \neq 0$. Then we put

$$
\begin{aligned}
S_{p, q}^{\bar{r}} A\left(\mathscr{R}, \mathbb{R}^{2}\right) & :=\left\{f \in S^{\prime}\left(\mathbb{R}^{2}\right): f \circ \mathscr{R} \in S_{p, q}^{\bar{r}} A\left(\mathbb{R}^{2}\right)\right\}, \\
\left\|f \mid S_{p, q}^{\bar{r}} A\left(\mathscr{R}, \mathbb{R}^{2}\right)\right\| & :=\left\|f \circ \mathscr{R} \mid S_{p, q}^{\bar{r}} A\left(\mathbb{R}^{2}\right)\right\| .
\end{aligned}
$$

Recall that for $p=q=2$ we have coincidence of $S_{2,2}^{\bar{r}} B\left(\mathscr{R}, \mathbb{R}^{2}\right)$ with $S_{2}^{\bar{r}} W\left(\mathscr{R}, \mathbb{R}^{2}\right)$ in the sense of equivalent norms, cf. [12] or [18, Thm. 2.3.1]. By means of these classes we are able to describe the trace classes for Besov as well as for Lizorkin-Triebel classes.

The counterpart of Theorem 2.9 for Besov spaces is as follows.

Theorem 3.10. Let $\mathscr{O}$ be an orthogonal basis of $\Gamma$ and let $\mathscr{R}_{i}, i=1,2,3$, be matrices associated with $\mathscr{O}$ by (1.1), (2.8) and (2.10).

Let $0<p, q \leqslant \infty$ and $\bar{r}=\left(r_{1}, r_{2}, r_{3}\right) \in \mathbb{R}^{3}$ with $r_{i} \neq 1 / p, i=1,2,3$, and

$$
\begin{equation*}
\min \left(r_{1}, r_{2}, r_{3}, r_{1}+r_{2}-\frac{1}{p}, r_{1}+r_{3}-\frac{1}{p}, r_{2}+r_{3}-\frac{1}{p}\right)>\sigma_{p} \tag{3.13}
\end{equation*}
$$

Then

$$
\begin{equation*}
\operatorname{tr}_{\mathscr{O}} \in \mathscr{L}\left(S_{p, q}^{\bar{r}} B\left(\mathbb{R}^{3}\right), S^{1}\left(\mathbb{R}^{2}\right)+S^{2}\left(\mathbb{R}^{2}\right)+S^{3}\left(\mathbb{R}^{2}\right)\right) \tag{3.14}
\end{equation*}
$$

where

$$
S^{1}\left(\mathbb{R}^{2}\right):= \begin{cases}S_{p, q}^{r_{2}, r_{3}} B\left(\mathscr{R}_{1}^{-1}, \mathbb{R}^{2}\right) & \text { if } r_{1}>1 / p \\ S_{p, q}^{r_{2}, r_{3}+r_{1}-\frac{1}{p}} B\left(\mathscr{R}_{1}^{-1}, \mathbb{R}^{2}\right) \cap S_{p, q}^{r_{2}+r_{1}-\frac{1}{p}, r_{3}} B\left(\mathscr{R}_{1}^{-1}, \mathbb{R}^{2}\right) & \text { if } r_{1}<1 / p\end{cases}
$$

and similarly for $S^{2}$ and $S^{3}$.
Conversely, to each function $g \in S^{1}\left(\mathbb{R}^{2}\right)+S^{2}\left(\mathbb{R}^{2}\right)+S^{3}\left(\mathbb{R}^{2}\right)$ there exists a function $f \in S_{p, q}^{\bar{r}} B\left(\mathbb{R}^{3}\right)$ such that $\operatorname{tr}_{\mathscr{O}} f=g$.

Proof. The restrictions in (3.13) are guaranteeing that we may apply Theorem 3.6 for $S_{p, q}^{\bar{r}} B\left(\mathbb{R}^{3}\right)$ as well as for all spaces appearing in the definition of the target spaces but taken with the identity matrix instead of $\mathscr{R}_{i}^{-1}, i \in\{1,2,3\}$.

Step 1. According to Theorem 3.6, each $f \in S_{p, q}^{\bar{r}} B\left(\mathbb{R}^{3}\right)$ may be decomposed into

$$
\begin{equation*}
f=\sum_{\bar{\nu} \in \mathbb{N}_{0}^{d}} \sum_{\bar{m} \in \mathbb{Z}^{d}} \lambda_{\bar{\nu} \bar{m}} a_{\bar{\nu} \bar{m}}(x), \tag{3.15}
\end{equation*}
$$

with

$$
\begin{equation*}
\left\|\lambda\left|s_{p, q}^{\bar{r}} b\|\leqslant c\| f\right| S_{p, q}^{\bar{r}} B\left(\mathbb{R}^{3}\right)\right\| \tag{3.16}
\end{equation*}
$$

with some constant $c$ independent of $f$. We require some additional regularity of the atoms, cf. Definition 3.5:

$$
\begin{equation*}
K_{i} \geqslant \max \left(\left[r_{1}\right]+\left[r_{2}\right]+2,\left[r_{1}\right]+\left[r_{3}\right]+2,\left[r_{2}\right]+\left[r_{3}\right]+2\right), \quad i=1,2,3 \tag{3.17}
\end{equation*}
$$

In view of Remark 3.7 we decompose $f$ into three parts $f_{i}, i=1,2,3$, where

$$
\begin{align*}
f_{1}(x) & :=\sum_{\nu_{1}=0}^{\infty} \sum_{\nu_{2}=\nu_{1}}^{\infty} \sum_{\nu_{3}=\nu_{1}}^{\infty} \sum_{\bar{m} \in \mathbb{Z}^{3}} \lambda_{\bar{\nu} \bar{m}} a_{\bar{\nu} \bar{m}}(x),  \tag{3.18}\\
f_{2}(x) & :=\sum_{\nu_{2}=0}^{\infty} \sum_{\nu_{1}=\nu_{2}+1}^{\infty} \sum_{\nu_{3}=\nu_{2}}^{\infty} \sum_{\bar{m} \in \mathbb{Z}^{3}} \lambda_{\bar{\nu} \bar{m}} a_{\bar{\nu} \bar{m}}(x),  \tag{3.19}\\
f_{3}(x) & :=\sum_{\nu_{3}=0}^{\infty} \sum_{\nu_{1}=\nu_{3}+1}^{\infty} \sum_{\nu_{2}=\nu_{3}+1}^{\infty} \sum_{\bar{m} \in \mathbb{Z}^{3}} \lambda_{\bar{\nu} \bar{m}} a_{\bar{\nu} \bar{m}}(x) . \tag{3.20}
\end{align*}
$$

This allows us to decompose $\operatorname{tr}_{\mathscr{O}} f$ into (see (2.7))

$$
\begin{equation*}
\left(\operatorname{tr}_{\mathscr{O}} f\right)\left(z_{1}, z_{2}\right)=\sum_{i=1}^{3}\left(\operatorname{tr}_{i} f_{i}\right)\left(\mathscr{R}_{i} \vec{z}\right) \tag{3.21}
\end{equation*}
$$

So, to establish (3.14) it is enough to prove the existence of a constant $c$ independent of $f$ such that

$$
\begin{equation*}
\left\|\operatorname{tr}_{1} f_{1}\left|S_{p, q}^{r_{2}, r_{3}} B\left(\mathbb{R}^{2}\right)\|\leqslant c\| f\right| S_{p, q}^{\bar{r}} B\left(\mathbb{R}^{3}\right)\right\| \tag{3.22}
\end{equation*}
$$

if $r_{1}>1 / p$ and

$$
\begin{align*}
& \left\|\operatorname{tr}_{1} f_{1}\left|S_{p, q}^{r_{2}, r_{3}+r_{1}-\frac{1}{p}} B\left(\mathbb{R}^{2}\right)\|\leqslant c\| f\right| S_{p, q}^{\bar{r}} B\left(\mathbb{R}^{3}\right)\right\|,  \tag{3.23}\\
& \left\|\operatorname{tr}_{1} f_{1}\left|S_{p, q}^{r_{2}+r_{1}-\frac{1}{p}, r_{3}} B\left(\mathbb{R}^{2}\right)\|\leqslant c\| f\right| S_{p, q}^{\bar{r}} B\left(\mathbb{R}^{3}\right)\right\| \tag{3.24}
\end{align*}
$$

if $r_{1}<1 / p$ and the corresponding analogues for $\operatorname{tr}_{i} f_{i}, i=2,3$.
Step 2. Proof of (3.22)-(3.24). We proceed similar to [25]. For brevity we put

$$
\begin{aligned}
& \Upsilon_{1}:=\left\{\bar{\nu} \in \mathbb{N}_{0}^{3}: \nu_{1} \leqslant \min \left(\nu_{2}, \nu_{3}\right)\right\}, \\
& \Upsilon_{2}:=\left\{\bar{\nu} \in \mathbb{N}_{0}^{3}: \nu_{2} \leqslant \min \left(\nu_{1}, \nu_{3}\right)\right\}, \\
& \Upsilon_{3}:=\left\{\bar{\nu} \in \mathbb{N}_{0}^{3}: \nu_{3} \leqslant \min \left(\nu_{1}, \nu_{2}\right)\right\} .
\end{aligned}
$$

Then

$$
\begin{equation*}
\operatorname{tr}_{1} f_{1}\left(x_{2}, x_{3}\right)=\sum_{\bar{\nu} \in \Upsilon_{1}} \sum_{\bar{m} \in B_{\bar{\nu}}} \lambda_{\bar{\nu} \bar{m}} a_{\bar{\nu} \bar{m}}\left(-x_{2}-x_{3}, x_{2}, x_{3}\right) \tag{3.25}
\end{equation*}
$$

where

$$
\begin{equation*}
B_{\bar{\nu}}:=\left\{\bar{m} \in \mathbb{Z}^{3}: \operatorname{supp} a_{\bar{\nu} \bar{m}} \cap \Gamma \neq \emptyset\right\} . \tag{3.26}
\end{equation*}
$$

Due to (3.9), for given $\bar{\nu} \in \Upsilon_{1}$ and $m_{2}, m_{3} \in \mathbb{Z}$, there are at most $N$ integers $m_{1} \in \mathbb{Z}$, such that $\bar{m}=\left(m_{1}, m_{2}, m_{3}\right) \in B_{\bar{\nu}}$. The number $N$ does not depend on $\bar{\nu}$ and $m_{2}, m_{3}$. To simplify notation we shall work only with one number $m_{1}$, denoted by $m_{1}\left(\bar{\nu}, m_{2}, m_{3}\right)$ or simply by $m_{1}$ if the values of $\bar{\nu}, m_{2}$ and $m_{3}$ are clear from context. Rewriting (3.25) this gives

$$
\begin{align*}
& \operatorname{tr}_{1} f_{1}\left(x_{2}, x_{3}\right)  \tag{3.27}\\
& =\sum_{\nu_{2}=0}^{\infty} \sum_{\nu_{3}=0}^{\infty} \sum_{\left(m_{2}, m_{3}\right) \in \mathbb{Z}^{2}} \sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)} \lambda_{\bar{\nu}\left(m_{1}, m_{2}, m_{3}\right)} a_{\bar{\nu}\left(m_{1}, m_{2}, m_{3}\right)}\left(-x_{2}-x_{3}, x_{2}, x_{3}\right) \\
& =\sum_{\nu_{2}=0}^{\infty} \sum_{\nu_{3}=0}^{\infty} \sum_{\left(m_{2}, m_{3}\right) \in \mathbb{Z}^{2}} \gamma_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)} b_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}\left(x_{2}, x_{3}\right),
\end{align*}
$$

where

$$
\begin{equation*}
\gamma_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}=\sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)}\left|\lambda_{\bar{\nu}\left(m_{1}, m_{2}, m_{3}\right)}\right|, \tag{3.28}
\end{equation*}
$$

$b_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}\left(x_{2}, x_{3}\right)=0$ if $\gamma_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}=0$, and

$$
\begin{aligned}
& b_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}\left(x_{2}, x_{3}\right) \\
& \quad=\frac{1}{\gamma_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}} \sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)} \\
& \lambda_{\bar{\nu}\left(m_{1}, m_{2}, m_{3}\right)} a_{\bar{\nu}\left(m_{1}, m_{2}, m_{3}\right)}\left(-x_{2}-x_{3}, x_{2}, x_{3}\right)
\end{aligned}
$$

if $\gamma_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}>0$. We recall that in this sum $m_{1}$ is an abbreviation for $m_{1}\left(\bar{\nu}, m_{2}, m_{3}\right)$.

Step 3. We claim that

1. $b_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}$ are atoms in the sense of Definition 3.5 related to $\left(\nu_{2}, \nu_{3}\right)$, $\left(m_{2}\right.$, $m_{3}$ ) up to a general constant.
2. $\left\|\gamma\left|s_{p, q}^{r_{2}, r_{3}} b\|\leqslant c\| \lambda\right| s_{p, q}^{\bar{r}} b\right\|$ if $r_{1}>1 / p$,
3. $\left\|\gamma\left|s_{p, q}^{r_{2}, r_{3}+r_{1}-1 / p} b\|\leqslant c\| \lambda\right| s_{p, q}^{\bar{r}} b\right\|$ and $\left\|\gamma\left|s_{p, q}^{r_{2}+r_{1}-1 / p, r_{3}} b\|\leqslant c\| \lambda\right| s_{p, q}^{\bar{r}} b\right\|$ if $r_{1}<1 / p$.

Substep 3.1. The proof of the first assertion is elementary, see Remark 3.7. Two comments are in order. The first one concerns regularity. If the components of $\bar{K}$ are large enough then $b$ is sufficiently smooth to satisfy (3.10) for some $\widetilde{K}$ such that we can apply Theorem 3.6 with respect to the target space, cf. (3.17). The second comment concerns the estimate (3.10). As claimed this estimate is satisfied by the functions $b_{\left(\nu_{2}, \nu_{3}\right),\left(m_{2}, m_{3}\right)}$ up to a general constant $c_{\alpha}$ depending on $\alpha$. Since we need to control a finite number of derivatives only we conclude that $C b_{\left(\nu_{2}, \nu_{3}\right),\left(m_{2}, m_{3}\right)}$ are atoms with $C^{-1}:=\max _{\alpha} c_{\alpha}$. This is enough for our purpose.

Substep 3.2. Let $r_{1}>1 / p$. Let $r_{1}-1 / p=\varepsilon_{1}+\varepsilon_{2}, \varepsilon_{i}>0, i=1,2$. Obviuosly, $\varepsilon_{1}>0$ guarantees

$$
\sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)}\left|\lambda_{\bar{\nu} \bar{m}}\right| \leqslant c_{1}\left(\sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)}\left|2^{\nu_{1} \varepsilon_{1}} \lambda_{\bar{\nu} \bar{m}}\right|^{p}\right)^{1 / p}
$$

Next we use $\varepsilon_{2}>0$ and obtain

$$
\begin{aligned}
\left\|\gamma \mid s_{p, q}^{r_{2}, r_{3}} b\right\|^{q} & =\sum_{\nu_{2}=0}^{\infty} \sum_{\nu_{3}=0}^{\infty} 2^{\left[\nu_{2}\left(r_{2}-\frac{1}{p}\right)+\nu_{3}\left(r_{3}-\frac{1}{p}\right)\right] q}\left(\sum_{\left(m_{2}, m_{3}\right) \in \mathbb{Z}^{2}}\left|\gamma_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}\right|^{p}\right)^{q / p} \\
& \leqslant c_{2} \sum_{\nu_{2}=0}^{\infty} \sum_{\nu_{3}=0}^{\infty} 2^{\left[\nu_{2}\left(r_{2}-\frac{1}{p}\right)+\nu_{3}\left(r_{3}-\frac{1}{p}\right)\right] q}\left(\sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)} 2^{\nu_{1} \varepsilon_{1} p} \sum_{\left(m_{2}, m_{3}\right) \in \mathbb{Z}^{2}}\left|\lambda_{\bar{\nu} \bar{m}}\right|^{p}\right)^{q / p} \\
& \leqslant c_{3} \sum_{\bar{\nu} \in \Upsilon_{1}} 2^{\bar{\nu} \cdot\left(\bar{r}-\frac{1}{p}\right) q}\left(\sum_{\left(m_{2}, m_{3}\right) \in \mathbb{Z}^{2}}\left|\lambda_{\bar{\nu} \bar{m}}\right|^{p}\right)^{q / p} \leqslant c_{3}\left\|\lambda \mid s_{p, q}^{\bar{r}} b\right\|^{q} .
\end{aligned}
$$

Substep 3.3. Let $r_{1}<1 / p$. To begin with let $p \geqslant 1$. The triangle inequality yields

$$
\begin{equation*}
\left(\sum_{\left(m_{2}, m_{3}\right) \in \mathbb{Z}^{2}}\left(\sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)}\left|\lambda_{\bar{\nu} \bar{m}}\right|\right)^{p}\right)^{1 / p} \leqslant \sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)}\left(\sum_{\left(m_{2}, m_{3}\right) \in \mathbb{Z}^{2}}\left|\lambda_{\bar{\nu} \bar{m}}\right|^{p}\right)^{1 / p} \tag{3.29}
\end{equation*}
$$

If now $q \leqslant 1$, we get

$$
\begin{aligned}
\left\|\gamma \left\lvert\, s_{p, q}^{r_{2}, r_{3}+r_{1}-\frac{1}{p}} b\right.\right\|^{q} & \leqslant \sum_{\nu_{2}=0}^{\infty} \sum_{\nu_{3}=0}^{\infty} 2^{\left[\nu_{2}\left(r_{2}-\frac{1}{p}\right)+\nu_{3}\left(r_{3}+r_{1}-\frac{1}{p}-\frac{1}{p}\right)\right] q} \\
& \leqslant \sum_{\bar{\nu} \in \Upsilon_{1}} 2^{\bar{\nu} \cdot\left(\bar{r}-\frac{1}{p}\right) q}\left(\sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)}\left(\sum_{\left.\left(m_{2}, m_{3}\right) \in \mathbb{m}_{3}\right) \in \mathbb{Z}^{2}}\left|\lambda_{\bar{\nu} \bar{m}}\right|^{p}\right)^{1 / p}\right]^{q} \\
& \left.\left.\lambda_{\bar{\nu} \bar{m}}\right|^{p}\right)^{q / p} \leqslant\left\|\lambda \mid s_{p, q}^{\bar{r}} b\right\|^{q}
\end{aligned}
$$

For $q>1$, we denote

$$
\varrho_{\bar{\nu}}:=\left(\sum_{\left(m_{2}, m_{3}\right) \in \mathbb{Z}^{2}}\left|\lambda_{\bar{\nu} \bar{m}}\right|^{p}\right)^{1 / p}
$$

and apply Hölder's inequality to obtain

$$
\begin{aligned}
&\left\|\gamma \left\lvert\, s_{p, q}^{r_{2}, r_{3}+r_{1}-\frac{1}{p}} b\right.\right\|^{q} \\
& \leqslant \sum_{\nu_{2}=0}^{\infty} \sum_{\nu_{3}=0}^{\infty} 2^{\left[\nu_{2}\left(r_{2}-\frac{1}{p}\right)+\nu_{3}\left(r_{3}+r_{1}-\frac{1}{p}-\frac{1}{p}\right)\right] q}\left[\sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)} 2^{\left(\nu_{3}-\nu_{1}\right)\left(r_{1}-\frac{1}{p}\right)} 2^{\left(\nu_{1}-\nu_{3}\right)\left(r_{1}-\frac{1}{p}\right)} \varrho_{\bar{\nu}}\right]^{q} \\
& \leqslant\left[\sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)} 2^{\left(\nu_{3}-\nu_{1}\right)\left(r_{1}-\frac{1}{p}\right) q^{\prime}}\right]^{q / q^{\prime}} \sum_{\bar{\nu} \in \Upsilon_{1}} 2^{\bar{\nu} \cdot\left(\bar{r}-\frac{1}{p}\right) q} \varrho_{\nu}^{q} \\
& \leqslant c\left\|\lambda \mid s_{p, q}^{\bar{r}} b\right\|^{q} .
\end{aligned}
$$

This proves our claims if $p \geqslant 1$. Now let $p<1$. We substitute (3.29) by

$$
\begin{aligned}
(3.30)\left(\sum_{\left(m_{2}, m_{3}\right) \in \mathbb{Z}^{2}}\left(\sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)}\left|\lambda_{\bar{\nu} \bar{m}}\right|\right)^{p}\right)^{1 / p} & \leqslant\left(\sum_{\left(m_{2}, m_{3}\right) \in \mathbb{Z}^{2}} \sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)}\left|\lambda_{\bar{\nu} \bar{m}}\right|^{p}\right)^{1 / p} \\
& =\left(\sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)} \varrho_{\bar{\nu}}^{p}\right)^{1 / p}
\end{aligned}
$$

If $q \leqslant p$ the monotonicity of the $\ell_{r}$-quasinorms yields

$$
\begin{align*}
\left\|\gamma \left\lvert\, s_{p, q}^{r_{2}, r_{3}+r_{1}-\frac{1}{p}} b\right.\right\|^{q} & \leqslant \sum_{\nu_{2}=0}^{\infty} \sum_{\nu_{3}=0}^{\infty} 2^{\left[\nu_{2}\left(r_{2}-\frac{1}{p}\right)+\nu_{3}\left(r_{3}+r_{1}-\frac{1}{p}-\frac{1}{p}\right)\right] q}\left[\sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)} \varrho_{\bar{\nu}}^{p}\right]^{q / p}  \tag{3.31}\\
& \leqslant \sum_{\bar{\nu} \in \Upsilon_{1}} 2^{\bar{\nu} \cdot\left(\bar{r}-\frac{1}{p}\right) q} \varrho_{\nu}^{q} \leqslant\left\|\lambda \mid s_{p, q}^{\bar{r}} b\right\|^{q} .
\end{align*}
$$

And for $q>p$, we combine (3.30) with Hölder's inequality

$$
\sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)} 2^{\left(\nu_{3}-\nu_{1}\right)\left(r_{1}-\frac{1}{p}\right) p} 2^{\left(\nu_{1}-\nu_{3}\right)\left(r_{1}-\frac{1}{p}\right) p} \varrho_{\bar{\nu}}^{p} \leqslant c\left(\sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)} 2^{\left(\nu_{1}-\nu_{3}\right)\left(r_{1}-\frac{1}{p}\right) q} \varrho_{\bar{\nu}}^{q}\right)^{p / q}
$$

to derive (3.31) again. Moreover, the second estimate in Claim 3 follows by interchanging the roles of $r_{2}$ and $r_{3}$. This completes the estimates claimed for $\gamma$.

Step 4. We shall prove the estimate for $\operatorname{tr}_{1} f_{1}$. In case when $r_{1}>1 / p$ we argue, by using Claim 2 and Theorem 3.6, first for $d=2$ and later for $d=3$, as follows

$$
\left\|\operatorname{tr}_{1} f_{1}\left|S_{p, q}^{r_{2}, r_{3}} B\left(\mathbb{R}^{2}\right)\left\|\leqslant c_{1}\right\| \gamma\right| s_{p, q}^{r_{2}, r_{3}} b\right\| \leqslant c_{2}\left\|\lambda\left|s_{p, q}^{\bar{r}} b\left\|\leqslant c_{3}\right\| f\right| S_{p, q}^{\bar{r}} B\left(\mathbb{R}^{3}\right)\right\| .
$$

Mutatis mutandis the case $r_{1}<1 / p$ can be treated. The estimates of $\operatorname{tr}_{i} f_{i}, i=2,3$ follow by symmetry.

Step 5. Now we construct the (non-)linear extension operator. We start with a function $g \in S^{1}\left(\mathbb{R}^{2}\right)+S^{2}\left(\mathbb{R}^{2}\right)+S^{3}\left(\mathbb{R}^{2}\right)$. Then there are $g_{i} \in S^{i}\left(\mathbb{R}^{2}\right), i=1,2,3$, such that $g_{i} \in S^{i}\left(\mathbb{R}^{2}\right), g=g_{1}+g_{2}+g_{3}$ and

$$
\left\|g_{i}\left|S^{i}\left(\mathbb{R}^{2}\right)\|\leqslant 2\| g\right| S^{1}\left(\mathbb{R}^{2}\right)+S^{2}\left(\mathbb{R}^{2}\right)+S^{3}\left(\mathbb{R}^{2}\right)\right\|
$$

We shall extend each $g_{i}$ separately. That is, we construct three functions $f_{1}, f_{2}, f_{3} \in$ $S_{p, q}^{\bar{r}} B\left(\mathbb{R}^{3}\right)$ such that $\operatorname{tr}_{\mathscr{O}} f_{i}=g_{i}, i=1,2,3$. The desirable extension will then be given by $f=f_{1}+f_{2}+f_{3}$.

Substep 5.1. We restrict ourselves to $i=1$, the other cases follow by symmetry. To begin with we treat the case $r_{1}>1 / p$. We put $h_{1}:=g_{1} \circ \mathscr{R}_{1}^{-1}$. Then $h_{1} \in$ $S_{p, q}^{r_{2}, r_{3}} B\left(\mathbb{R}^{2}\right)$ and, according to (2.7), we get

$$
g_{1}\left(z_{1}, z_{2}\right)=\left(\operatorname{tr}_{\mathscr{O}} f_{1}\right)\left(z_{1}, z_{2}\right)=\left(\operatorname{tr}_{1} f_{1}\right)\left(\mathscr{R}_{1} \vec{z}\right)
$$

for all $\vec{z}=\left(z_{1}, z_{2}\right) \in \mathbb{R}^{2}$ if, and only if,

$$
g_{1}\left(\mathscr{R}_{1}^{-1} \vec{z}\right)=h_{1}(\vec{z})=\left(\operatorname{tr}_{1} f_{1}\right)\left(z_{1}, z_{2}\right)=f_{1}\left(-z_{1}-z_{2}, z_{1}, z_{2}\right), \quad \vec{z}=\left(z_{1}, z_{2}\right) \in \mathbb{R}^{2} .
$$

Hence, our original task, namely to find $f_{1}$ such that $\operatorname{tr}_{\mathscr{O}} f_{1}=g_{1}$, where $g_{1} \in$ $S_{p, q}^{r_{2}, r_{3}} B\left(\mathscr{R}_{1}^{-1}, \mathbb{R}^{2}\right)$ is given, can be replaced by searching for $f_{1}$ such that $\operatorname{tr}_{1} f_{1}=h_{1}$, where $h_{1} \in S_{p, q}^{r_{2}, r_{3}} B\left(\mathbb{R}^{2}\right)$. Again we make use of atomic decompositions. According to Theorem 3.6 we can decompose

$$
h_{1}\left(x_{2}, x_{3}\right)=\sum_{\left(\nu_{2}, \nu_{3}\right) \in \mathbb{N}_{0}^{2}} \sum_{\left(m_{2}, m_{3}\right) \in \mathbb{Z}^{2}} \gamma_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)} b_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}\left(x_{2}, x_{3}\right),
$$

where

$$
c_{1}\left\|\gamma\left|s_{p, q}^{r_{2}, r_{3}} b\|\leqslant\| h_{1}\right| S_{p, q}^{r_{2}, r_{3}} B\left(\mathbb{R}^{2}\right)\right\| \leqslant c_{2}\left\|\gamma \mid s_{p, q}^{r_{2}, r_{3}} b\right\|
$$

for certain positive constants $c_{1}$ and $c_{2}$ independent of $h_{1}$. Now we choose an integer $m_{1}$ such that $\left|2^{-\nu_{1}} m_{1}+2^{-\nu_{2}} m_{2}+2^{-\nu_{3}} m_{3}\right| \leqslant 2^{-\nu_{1}}$ and define

$$
a_{\bar{\nu} \bar{m}}\left(x_{1}, x_{2}, x_{3}\right):=\psi\left(2^{\nu_{1}} x_{1}-m_{1}\right) b_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}\left(x_{2}, x_{3}\right),
$$

where

$$
\psi \in \mathscr{S}(\mathbb{R}), \quad \operatorname{supp} \psi \subset[-2(1+\delta), 2(1+\delta)], \quad \psi(t)=1 \text { if } t \in[-(1+\delta),(1+\delta)]
$$

and $\delta$ is the number from (3.9). For $\nu_{1} \leqslant \min \left(\nu_{2}, \nu_{3}\right)$ some easy calculations yield

$$
a_{\bar{\nu} \bar{m}}\left(-x_{2}-x_{3}, x_{2}, x_{3}\right)=b_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}\left(x_{2}, x_{3}\right), \quad\left(x_{2}, x_{3}\right) \in \mathbb{R}^{2} .
$$

If the first component of $\bar{m}$ differs from this specific $m_{1}$ then we define $a_{\bar{\nu} \bar{m}} \equiv 0$. Further, we put

$$
\lambda_{\left(\nu_{1}, \nu_{2}, \nu_{3}\right)\left(m_{1}, m_{2}, m_{3}\right)}:= \begin{cases}\gamma_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)} & \text { if } \nu_{1}=0  \tag{3.32}\\ 0 & \text { otherwise }\end{cases}
$$

and

$$
f_{1}:=\operatorname{ext} h_{1}=\sum_{\bar{\nu} \in \Upsilon_{1}} \sum_{\bar{m} \in \mathbb{Z}^{3}} \lambda_{\bar{\nu} \bar{m}} a_{\bar{\nu} \bar{m}} .
$$

Then

$$
\begin{equation*}
\left\|\operatorname{ext} h_{1}\left|S_{p, q}^{\bar{r}} B\left(\mathbb{R}^{3}\right)\left\|\leqslant C_{1}\right\| \lambda\right| s_{p, q}^{\bar{r}} b\right\|=C_{1}\left\|\gamma\left|s_{p, q}^{r_{2}, r_{3}} b\left\|\leqslant C_{2}\right\| h_{1}\right| S_{p, q}^{r_{2}, r_{3}} B\left(\mathbb{R}^{2}\right)\right\| \tag{3.33}
\end{equation*}
$$

This shows that $f_{1}$ represents an appropriate extension of $h_{1}$ if $r_{1}>1 / p$.
Substep 5.2. Let $r_{1}<1 / p$. First of all notice that this time $h_{1} \in S_{p, q}^{r_{2}, r_{3}+r_{1}-1 / p}$ $B\left(\mathbb{R}^{2}\right) \cap S_{p, q}^{r_{2}+r_{1}-1 / p, r_{3}} B\left(\mathbb{R}^{2}\right)$. We have to modify the definition of $\lambda$, cf. (3.32). This time we use

$$
\lambda_{\left(\nu_{1}, \nu_{2}, \nu_{3}\right)\left(m_{1}, m_{2}, m_{3}\right)}:= \begin{cases}\gamma_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)} & \text { if } \nu_{1}=\min \left(\nu_{2}, \nu_{3}\right)  \tag{3.34}\\ 0 & \text { otherwise },\end{cases}
$$

for the specific value of $m_{1}$ as chosen in Substep 5.1. In all other cases we put $\lambda_{\bar{\nu} \bar{m}}=0$. Then

$$
\begin{aligned}
&\left\|\operatorname{ext} h_{1}\left|S_{p, q}^{\bar{r}} B\left(\mathbb{R}^{3}\right)\left\|^{q} \leqslant C_{1}\right\| \lambda\right| s_{p, q}^{\bar{r}} b\right\|^{q} \\
&= C_{1} \sum_{\nu_{2}=0}^{\infty} \sum_{\nu_{3}=0}^{\infty} 2^{\left[\min \left(\nu_{2}, \nu_{3}\right)\left(r_{1}-1 / p\right)+\nu_{2}\left(r_{2}-\frac{1}{p}\right)+\nu_{3}\left(r_{3}-\frac{1}{p}\right)\right] q}\left(\sum_{\left(m_{2}, m_{3}\right) \in \mathbb{Z}^{2}}\left|\gamma_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}\right|^{p}\right)^{q / p} \\
&= C_{1}\left(\sum_{\nu_{2}=0}^{\infty} \sum_{\nu_{3}=\nu_{2}}^{\infty} 2^{\left[\nu_{2}\left(r_{1}-1 / p\right)+\nu_{2}\left(r_{2}-\frac{1}{p}\right)+\nu_{3}\left(r_{3}-\frac{1}{p}\right)\right] q}\left(\sum_{\left(m_{2}, m_{3}\right) \in \mathbb{Z}^{2}}\left|\gamma_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}\right|^{p}\right)^{q / p}\right. \\
&\left.+\sum_{\nu_{3}=0}^{\infty} \sum_{\nu_{2}=\nu_{3}+1}^{\infty} 2^{\left[\nu_{3}\left(r_{1}-1 / p\right)+\nu_{2}\left(r_{2}-\frac{1}{p}\right)+\nu_{3}\left(r_{3}-\frac{1}{p}\right)\right] q}\left(\sum_{\left(m_{2}, m_{3}\right) \in \mathbb{Z}^{2}}\left|\gamma_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}\right|^{p}\right)^{q / p}\right) \\
& \leqslant C_{1}\left(\left\|\gamma\left|s_{p, q}^{r_{2}+r_{1}-\frac{1}{p}, r_{3}} b\|+\| \gamma\right| s_{p, q}^{r_{2}, r_{3}+r_{1}-\frac{1}{p}} b\right\|\right)^{q} \\
& \leqslant C_{2}\left(\left\|h_{1}\left|S_{p, q}^{r_{2}+r_{1}-\frac{1}{p}, r_{3}} B\left(\mathbb{R}^{2}\right)\|+\| h_{1}\right| S_{p, q}^{r_{2}, r_{3}+r_{1}-\frac{1}{p}} B\left(\mathbb{R}^{2}\right)\right\|\right)^{q} .
\end{aligned}
$$

Hence, also in this situation we have an appropriate extension of $g_{1}$. The modifications for an extension of $g_{2}$ and $g_{3}$ are obvious.

Remark 3.11. The reader may notice that the only possible failure of linearity of the extension operator comes from the (generally non-linear) decomposition of $g$ into $g=g_{1}+g_{2}+g_{3}$.

It remains to consider the limiting cases where at least one of the $r_{i}$ equals $1 / p$. We concentrate on the more simple situation where $0<p, q \leqslant 1$.

Proposition 3.12. Let $\mathscr{O}$ be an orthogonal basis of $\Gamma$ and let $\mathscr{R}_{i}, i=1,2,3$, be matrices associated with $\mathscr{O}$ by (1.1), (2.8) and (2.10).

Let $0<p, q \leqslant 1$. Then the statement of Theorem 3.10 remains true without the assumption $r_{i} \neq 1 / p, i=1,2,3$.

Proof. The proof of Theorem 3.10 extends to the present situation since in Substep 3.2 one can work with $\varepsilon_{1}=\varepsilon_{2}=0$.

Remark 3.13. Proposition 3.12 does not extend to values of $p$ larger than 1 . In analogy to the two-dimensional situation, cf. [25] for details, more complicated spaces occur. We omit the details.

### 3.5. Traces of Lizorkin-Triebel spaces.

Now we turn to the Lizorkin-Triebel classes. To prove an analog of Theorem 3.10 for these spaces we can proceed in the same way as in case of the Besov spaces. We shall describe only the modifications needed.

Theorem 3.14. Let $\mathscr{O}$ be an orthogonal basis of $\Gamma$ and let $\mathscr{R}_{i}, i=1,2,3$, be matrices associated with $\mathscr{O}$ by (1.1), (2.8) and (2.10). Let $0<p<\infty$ and $0<q \leqslant \infty$. Let $\bar{r}=\left(r_{1}, r_{2}, r_{3}\right) \in \mathbb{R}^{3}$ with

$$
\begin{equation*}
\min \left(r_{1}, r_{2}, r_{3}\right)>\max \left(\frac{1}{p}, \sigma_{p q}\right) \tag{3.35}
\end{equation*}
$$

Then

$$
\begin{align*}
\operatorname{tr}_{\mathscr{O}} \in \mathscr{L}\left(S_{p, q}^{\bar{r}} F\left(\mathbb{R}^{3}\right), S_{p, q}^{r_{2}, r_{3}} F\left(\mathscr{R}_{1}^{-1}, \mathbb{R}^{2}\right)\right. & +S_{p, q}^{r_{1}, r_{3}} F\left(\mathscr{R}_{2}^{-1}, \mathbb{R}^{2}\right)  \tag{3.36}\\
& \left.+S_{p, q}^{r_{1}, r_{2}} F\left(\mathscr{R}_{3}^{-1}, \mathbb{R}^{2}\right)\right) .
\end{align*}
$$

Conversely, to each function $g \in S_{p, q}^{r_{2}, r_{3}} F\left(\mathscr{R}_{1}^{-1}, \mathbb{R}^{2}\right)+S_{p, q}^{r_{1}, r_{3}} F\left(\mathscr{R}_{2}^{-1}, \mathbb{R}^{2}\right)+S_{p, q}^{r_{1}, r_{2}} \times$ $F\left(\mathscr{R}_{3}^{-1}, \mathbb{R}^{2}\right)$ there exists a function $f \in S_{p, q}^{\bar{r}} F\left(\mathbb{R}^{3}\right)$ such that $\operatorname{tr}_{\mathscr{O}} f=g$.

Proof. We shall use the same notation as in the proof of Theorem 3.10.
Step 1. Boundedness. In Step 1 of the proof of Theorem 3.10 we simply change the letter $B$ to $F$. In Step 2 nothing changes and we concentrate on Step 3 now. We have to prove that

$$
\begin{equation*}
\left\|\gamma\left|s_{p, q}^{r_{2}, r_{3}} f\|\leqslant c\| \lambda\right| s_{p, q}^{r_{1}, r_{2}, r_{3}} f\right\| \tag{3.37}
\end{equation*}
$$

with some $c$ independent of $\lambda$.
Instead we shall prove a pointwise inequality. So, first we fix a point $\left(x_{2}, x_{3}\right) \in \mathbb{R}^{2}$. Then there is only one element $\left(m_{2}, m_{3}\right) \in \mathbb{Z}^{2}$ such that $\chi_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}\left(x_{2}, x_{3}\right)=1$. We denote $\gamma_{\left(\nu_{2}, \nu_{3}\right)}=\gamma_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}$. Similarly, for each $\bar{\nu}=\left(\nu_{1}, \nu_{2}, \nu_{3}\right)$, there is a unique $\bar{m}(\bar{\nu})=\left(m_{1}, m_{2}, m_{3}\right)$ such that $\chi_{\left(\nu_{1}, \nu_{2}, \nu_{3}\right)\left(m_{1}, m_{2}, m_{3}\right)}\left(x_{1}, x_{2}, x_{3}\right)=1$ and $\bar{m} \in B_{\bar{\nu}}$. We denote $\lambda_{\bar{\nu}}=\lambda_{\bar{\nu} \bar{m}}$.

Substep 2.1. Let $r_{1}>1 / p$ and $0<q \leqslant 1$. Then

$$
\left|\gamma_{\left(\nu_{2}, \nu_{3}\right)}\right|^{q}=\left(\sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)}\left|\lambda_{\left(\nu_{1}, \nu_{2}, \nu_{3}\right)}\right|\right)^{q} \leqslant \sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)}\left|\lambda_{\left(\nu_{1}, \nu_{2}, \nu_{3}\right)}\right|^{q},
$$

and

$$
\left(\sum_{\nu_{2}, \nu_{3}=0}^{\infty} 2^{\left[\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left|\gamma_{\left(\nu_{2}, \nu_{3}\right)}\right|^{q}\right)^{p / q} \leqslant\left(\sum_{\nu_{1}=0}^{\infty} \sum_{\nu_{2}, \nu_{3}=\nu_{1}}^{\infty} 2^{\left[\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left|\lambda_{\bar{\nu}}\right|^{q}\right)^{p / q} .
$$

To continue we distinguish two cases. Let $0<p \leqslant q$. Then

$$
\begin{aligned}
\left(\sum_{\nu_{1}=0}^{\infty} \sum_{\nu_{2}, \nu_{3}=\nu_{1}}^{\infty} 2^{\left[\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left|\lambda_{\bar{\nu}}\right|^{q}\right)^{p / q} & \leqslant \sum_{\nu_{1}=0}^{\infty}\left(\sum_{\nu_{2}, \nu_{3}=\nu_{1}}^{\infty} 2^{\left[\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left|\lambda_{\bar{\nu}}\right|^{q}\right)^{p / q} \\
& \leqslant \sum_{\nu_{1}=0}^{\infty} 2^{-\nu_{1} r_{1} p}\left(\sum_{\nu_{2}, \nu_{3}=\nu_{1}}^{\infty} 2^{\left[\nu_{1} r_{1}+\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left|\lambda_{\bar{\nu}}\right|^{q}\right)^{p / q}
\end{aligned}
$$

Now let $0<q<p<\infty$. With $0<\varepsilon<r_{1} p-1$ and applying Hölder's inequality we find

$$
\begin{aligned}
&\left(\sum_{\nu_{1}=0}^{\infty} \sum_{\nu_{2}, \nu_{3}=\nu_{1}}^{\infty} 2^{\left[\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left|\lambda_{\bar{\nu}}\right|^{q}\right)^{p / q} \\
& \leqslant c \sum_{\nu_{1}=0}^{\infty} 2^{-\nu_{1}\left(r_{1} p-\varepsilon\right)}\left(\sum_{\nu_{2}, \nu_{3}=\nu_{1}}^{\infty} 2^{\left[\nu_{1} r_{1}+\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left|\lambda_{\bar{\nu}}\right|^{q}\right)^{p / q} \\
& \leqslant c \sum_{\nu_{1}=0}^{\infty} 2^{-\nu_{1}}\left(\sum_{\nu_{2}, \nu_{3}=\nu_{1}}^{\infty} 2^{\left[\nu_{1} r_{1}+\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left|\lambda_{\bar{\nu}}\right|^{q}\right)^{p / q}
\end{aligned}
$$

Substep 2.2. If $q>1$ we use triangle inequality

$$
\begin{gathered}
\left(\sum_{\nu_{2}, \nu_{3}=0}^{\infty} 2^{\left[\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left(\sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)}\left|\lambda_{\bar{\nu}}\right|\right)^{q}\right)^{1 / q} \leqslant \sum_{\nu_{1}=0}^{\infty}\left(\sum_{\nu_{2}, \nu_{3}=\nu_{1}}^{\infty} 2^{\left[\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left|\lambda_{\bar{\nu}}\right|^{q}\right)^{1 / q} \\
\leqslant \sum_{\nu_{1}=0}^{\infty} 2^{-\nu_{1} r_{1}}\left(\sum_{\nu_{2}, \nu_{3}=\nu_{1}}^{\infty} 2^{\left[\nu_{1} r_{1}+\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left|\lambda_{\bar{\nu}}\right|^{q}\right)^{1 / q}
\end{gathered}
$$

If $0<p \leqslant 1$

$$
\begin{aligned}
& \left(\sum_{\nu_{2}, \nu_{3}=0}^{\infty} 2^{\left[\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left(\sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)}\left|\lambda_{\bar{\nu}}\right|\right)^{q}\right)^{p / q} \\
& \quad \leqslant \sum_{\nu_{1}=0}^{\infty} 2^{-\nu_{1} r_{1} p}\left(\sum_{\nu_{2}, \nu_{3}=\nu_{1}}^{\infty} 2^{\left[\nu_{1} r_{1}+\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left|\lambda_{\bar{\nu}}\right|^{q}\right)^{p / q}
\end{aligned}
$$

follows. If $p>1$ we apply again Hölder's inequality and find

$$
\begin{aligned}
\left(\sum_{\nu_{2}, \nu_{3}=0}^{\infty} 2^{\left[\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\right. & \left.\left(\sum_{\nu_{1}=0}^{\min \left(\nu_{2}, \nu_{3}\right)}\left|\lambda_{\bar{\nu}}\right|\right)^{q}\right)^{p / q} \\
& \leqslant c \sum_{\nu_{1}=0}^{\infty} 2^{-\nu_{1}\left(r_{1} p-\varepsilon\right)}\left(\sum_{\nu_{2}, \nu_{3}=\nu_{1}}^{\infty} 2^{\left[\nu_{1} r_{1}+\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left|\lambda_{\bar{\nu}}\right|^{q}\right)^{p / q} \\
& \leqslant c \sum_{\nu_{1}=0}^{\infty} 2^{-\nu_{1}}\left(\sum_{\nu_{2}, \nu_{3}=\nu_{1}}^{\infty} 2^{\left[\nu_{1} r_{1}+\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left|\lambda_{\bar{\nu}}\right|^{q}\right)^{p / q}
\end{aligned}
$$

Substep 2.3. Summarizing in all situations we have found

$$
\begin{align*}
& \left(\sum_{\nu_{2}, \nu_{3}=0}^{\infty} 2^{\left[\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left|\gamma_{\left(\nu_{2}, \nu_{3}\right)}\right|^{q}\right)^{p / q}  \tag{3.38}\\
& \leqslant c \sum_{\nu_{1}=0}^{\infty} 2^{-\nu_{1}}\left(\sum_{\nu_{2}, \nu_{3}=\nu_{1}}^{\infty} 2^{\left[\nu_{1} r_{1}+\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left|\lambda_{\bar{\nu}}\right|^{q}\right)^{p / q}
\end{align*}
$$

where $c$ does not depend on $\lambda$. We have to show that this inequality implies (3.37). For fixed $\left(x_{2}, x_{3}\right)$ we choose a sequence of intervals $I_{\nu_{1}}$ such that

$$
I_{\nu_{1}} \cap I_{\nu_{1}^{\prime}}=\emptyset, \quad \nu_{1} \neq \nu_{1}^{\prime}, \quad\left|I_{\nu_{1}}\right| \geqslant c 2^{-\nu_{1}}
$$

for some $c>0$ and

$$
\left\{\left(x_{1}, x_{2}, x_{3}\right): x_{1} \in I_{\nu_{1}}\right\} \subset Q_{\bar{\nu} \bar{m}}, \quad \bar{\nu} \in \Upsilon_{1}, \quad \bar{m} \in B_{\bar{\nu}}
$$

Then (3.38) implies

$$
\begin{aligned}
& \left(\sum_{\nu_{2}, \nu_{3}=0}^{\infty} 2^{\left[\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left|\gamma_{\left(\nu_{2}, \nu_{3}\right)} \chi_{\left(\nu_{2}, \nu_{3}\right)\left(m_{2}, m_{3}\right)}\left(x_{2}, x_{3}\right)\right|^{q}\right)^{p / q} \\
& \quad \leqslant c \sum_{\nu_{1}=0}^{\infty} \int_{I_{\nu_{1}}}\left(\sum_{\nu_{2}, \nu_{3}=\nu_{1}}^{\infty} 2^{\left[\nu_{1} r_{1}+\nu_{2} r_{2}+\nu_{3} r_{3}\right] q}\left|\lambda_{\bar{\nu}} \chi_{\bar{\nu}}\left(x_{1}, x_{2}, x_{3}\right)\right|^{q}\right)^{p / q} \mathrm{~d} x_{1} .
\end{aligned}
$$

Integration with respect to $x_{2}$ and $x_{3}$ completes the proof of the boundedness of $\operatorname{tr}_{1} f_{1}$. The rest is the same as in the $B$-case.

Step 2. The extension. Here the same construction as in the $B$-case can be applied, cf. Substep 5.1 of the proof of Theorem 3.10.

The above proof can be used also if that some of the $r_{i}$ coincide with $1 / p$, at least under additional restrictions on $p$ and $q$.

Proposition 3.15. Let $\mathscr{O}$ be an orthogonal basis of $\Gamma$ and let $\mathscr{R}_{i}, i=1,2,3$, be matrices associated with $\mathscr{O}$ by (1.1), (2.8) and (2.10).

Let $0<p \leqslant \min (1, q)$. Then the statement of Theorem 3.14 remains true under the weaker restriction

$$
\min \left(r_{1}, r_{2}, r_{3}\right) \geqslant \frac{1}{p} \quad \text { and } \quad \min \left(r_{1}, r_{2}, r_{3}\right)>\sigma_{p, q}
$$

Remark 3.16. A final remark. In the general situation of the Besov-LizorkinTriebel spaces we have proved a full counterpart of Theorem 2.9. In fact, it is not only a counterpart. Based on the identities $S_{2}^{\bar{r}} W\left(\mathbb{R}^{3}\right)=S_{2,2}^{\bar{r}} F\left(\mathbb{R}^{3}\right)=S_{2,2}^{\bar{r}} B\left(\mathbb{R}^{3}\right)$ (in the sense of equivalent norms) we have given a new proof of Theorem 2.9. Because of $S_{p}^{\bar{r}} W\left(\mathbb{R}^{3}\right)=S_{p, 2}^{\bar{r}} F\left(\mathbb{R}^{3}\right), 1<p<\infty$ (also in the sense of equivalent norms), Theorem 3.14 contains the extension to Sobolev spaces of dominating mixed smoothness with $p$ different from 2. However, we do not have counterparts of Theorems 2.11 and 2.13 , respectively. Here a good description of the spaces $S_{p, q}^{r_{1}, r_{2}} A\left(\mathscr{R}, \mathbb{R}^{2}\right)$ in terms of atoms would be desirable, see Lemma 2.2 (iii) for the Fourier-analytic counterpart.

## References

[1] D. R. Adams and L. I. Hedberg: Function spaces and potential theory. Springer, Berlin, 1996.
[2] T. I. Amanov: Spaces of differentiable functions with dominating mixed derivatives. Nauka Kaz. SSR, Alma-Ata, 1976. (In Russian.)
[3] D. B. Bazarkhanov: Characterizations of the Nikol'skii-Besov and Lizorkin-Triebel function spaces of mixed smoothness. Trudy Mat. Inst. Steklov 243 (2003), 53-65 ; 2Dipc. Steklov Inst. 243 (2003), 53-65.
[4] H.-J. Bungartz and M. Griebel: Sparse grids. Acta Numerica (2004), 1-123.
[5] M. Frazier and B. Jawerth: Decomposition of Besov spaces. Indiana Univ. Math. J. 34 (1985), 777-799.
[6] M. Frazier and B. Jawerth: A discrete transform and decomposition of distribution spaces. J. Funct. Anal. 93 (1990), 34-170.
[7] I. W. Gelman and W. G. Maz'ya: Abschätzungen für Differentialoperatoren im Halbraum. Akademie-Verlag, Berlin, 1981.
[8] M. Griebel, P. Oswald and T. Schiekofer: Sparse grids for boundary integral equations. Numer. Math 83 (1999), 279-312.
[9] P. Grisvard: Commutativité de deux fonctuers d'interpolation et applications. J. Math. Pures Appl. 45 (1966), 143-290.
[10] L.I. Hedberg and Y. Netrusov: An axiomatic approach to function spaces, spectral synthesis, and Luzin approximation. Memoirs of the AMS (to appear).
[11] R. Hochmuth: A $\varphi$-transform result for spaces with dominating mixed smoothness properties. Result. Math. 33 (1998), 106-119.
[12] P. I. Lizorkin and S. M. Nikol'skij: Classification of differentiable functions on the basis of spaces with dominating mixed smoothness. Trudy Mat. Inst. Steklov 77 (1965), 143-167.
[13] S. M. Nikol'skij: On boundary properties of differentiable functions of several variables. Dokl. Akad. Nauk SSSR 146 (1962), 542-545.
[14] S. M. Nikol'skij: On stable boundary values of differentiable functions of several variables. Mat. Sbornik 61 (1963), 224-252.
[15] P.-A. Nitsche: Sparse approximation of singularity functions. Constr. Approx. 21 (2005), 63-81.
zbl
[16] M. C. Rodríguez Fernández: U̇ber die Spur von Funktionen mit dominierenden gemischten Glattheitseigenschaften auf der Diagonale, Ph.D-thesis, Jena. 1997.
[17] H.-J. Schmeisser: Vector-valued Sobolev and Besov spaces. Teubner-Texte zur Math. 96 (1987), 4-44.
[18] H.-J. Schmeisser and H. Triebel: Topics in Fourier analysis and function spaces. Wiley, Chichester, 1987.
[19] G. Sparr: Interpolation of several Banach spaces. Ann. Mat. Pura Appl. 99 (1974), 247-316.
[20] H. Triebel: A diagonal embedding theorem for function spaces with dominating mixed smoothness properties. Banach Center Publications, Warsaw 22 (1989), 475-486.
[21] H. Triebel: Fractals and Spectra. Birkhäuser, Basel, 1997.
[22] H. Triebel: The Structure of Functions. Birkhäuser, Basel, 2001.
[23] J. Vybiral: Characterisations of function spaces with dominating mixed smoothness properties. Jenaer Schriften zur Mathematik und Informatik, Math/Inf/15/03, 2003.
[24] J. Vybíral: Function spaces with dominating mixed smoothness. Diss. Math. 436 (2006), 1-73.
[25] J. Vybíral: A diagonal embedding theorem for function spaces with dominating mixed smoothness. Funct. et Approx. 33 (2005), 101-120.
[27] H. Yserentant: On the regularity of the electronic Schrödinger equation in Hilbert spaces of mixed derivatives. Numer. Math. 98 (2004), 731-759.
[28] H. Yserentant: Sparse grid spaces for the numerical solution of the electronic Schrödinger equation. Numer. Math. 101 (2005), 381-389.

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