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## DIMENSION-FREE $I^{p}$ ESTIMATES FOR VEC TORS OF RIESZ TRANSFORMS ASSOchated wIIH ORTHOcONAL EXPANSIONS

# DIMENSION-FREE $L^{p}$ ESTIMATES FOR VECTORS OF RIESZ TRANSFORMS ASSOCIATED WITH ORTHOGONAL EXPANSIONS 

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An explicit Bellman function is used to prove a bilinear embedding theorem for operators associated with general multidimensional orthogonal expansions on product spaces. This is then applied to obtain $L^{p}$ boundedness, $1<p<\infty$, of appropriate vectorial Riesz transforms, in particular in the case of Jacobi polynomials. Our estimates for the $L^{p}$ norms of these Riesz transforms are both dimension-free and linear in $\max (p, p /(p-1))$. The approach we present allows us to avoid the use of both differential forms and general spectral multipliers.

## 1. Introduction

The classical Riesz transforms on $\mathbb{R}^{d}$ are the operators

$$
R_{i} f(x)=\partial_{x_{i}}\left(-\Delta_{\mathbb{R}^{d}}\right)^{-1 / 2} f(x), \quad i=1, \ldots, d .
$$

E. M. Stein [1983] proved that the vector of Riesz transforms

$$
\boldsymbol{R} f=\left(R_{1} f, \ldots, R_{d} f\right)
$$

has $L^{p}$ bounds which are independent of the dimension. More precisely

$$
\begin{equation*}
\|\boldsymbol{R} f\|_{L^{p}\left(\mathbb{R}^{d}\right)} \leqslant C_{p}\|f\|_{L^{p}\left(\mathbb{R}^{d}\right)}, \quad 1<p<\infty \tag{1-1}
\end{equation*}
$$

where $C_{p}$ is independent of the dimension $d$. Note that (1-1) is formally the same as the a priori bound

$$
\||\nabla f|\|_{L^{p}\left(\mathbb{R}^{d}\right)} \leqslant C_{p}\left\|(-\Delta)^{1 / 2} f\right\|_{L^{p}\left(\mathbb{R}^{d}\right)} .
$$

Later it was realized that, for $1<p<2$, one may take $C_{p} \leqslant C(p-1)^{-1}$ in (1-1); see [Bañuelos 1986; Duoandikoetxea and Rubio de Francia 1985]. It is worth mentioning that the best constant in (1-1) remains unknown when $d \geq 2$; the best results to date are given in [Bañuelos and Wang 1995] (see also [Dragičević and Volberg 2006] for an analytic proof) and [Iwaniec and Martin 1996].

The main goal of this paper is to generalize (1-1) to product settings different from $\mathbb{R}^{d}=\mathbb{R} \times \cdots \times \mathbb{R}$ with the product Lebesgue measure. Our starting point is the observation that the classical Riesz transform can be written as $R_{i}=\delta_{i}\left(\sum_{i=1}^{d} L_{i}\right)^{-1 / 2}$, where $\delta_{i}=\partial_{x_{i}}$, and $L_{i}=\delta_{i}^{*} \delta_{i}$. The generalized Riesz transforms we pursue are of the same form,

$$
\begin{equation*}
R_{i}=\delta_{i} L^{-1 / 2}, \quad i=1, \ldots, d, \tag{1-2}
\end{equation*}
$$

[^0]with $\delta_{i}$ being an operator on $L^{2}\left(X_{i}, \mu_{i}\right)$,
$$
L_{i}=\delta_{i}^{*} \delta_{i}+a_{i} \quad \text { and } \quad L=\sum_{i=1}^{d} L_{i}
$$

Here $a_{i}$ is a nonnegative constant. The adjoint $\delta_{i}^{*}$ is taken with respect to the inner product on $L^{2}\left(X_{i}, \mu_{i}\right)$, where $\mu_{i}$ is a nonnegative Borel measure on $X_{i}$ such that $d \mu_{i}\left(x_{i}\right)=w_{i}\left(x_{i}\right) d x_{i}$ for some positive and smooth function $w_{i}$ on $X_{i}$. To be precise, if 0 is an $L^{2}$ eigenvalue of $L$, then the definition of $R_{i}$ needs to be slightly modified; this is properly explained in the next section. Throughout the paper we assume that each $X_{i}, i=1, \ldots, d$, is an open interval in $\mathbb{R}$, an open half-line in $\mathbb{R}$, or the real line; we also set $X=X_{1} \times \cdots \times X_{d}$ and $\mu=\mu_{1} \otimes \cdots \otimes \mu_{d}$. We consider $\delta_{i}$ given by

$$
\delta_{i} f(x)=p_{i}\left(x_{i}\right) \partial_{x_{i}}+q_{i}\left(x_{i}\right), \quad x_{i} \in X_{i},
$$

for some real-valued functions $p_{i} \in C^{\infty}\left(X_{i}\right)$ and $q_{i} \in C^{\infty}\left(X_{i}\right)$. We remark that a significant difference between the classical Riesz transforms and the general Riesz transforms (1-2) lies in the fact that the operators $\delta_{i}$ and $\delta_{i}^{*}$ do not need to commute.

There are two assumptions which are critical to our results. Firstly, a computation, see [Nowak and Stempak 2006, p. 683], shows that the commutator [ $\delta_{i}, \delta_{i}^{*}$ ] is a function which we call $v_{i}$. We assume that $v_{i}$ is nonnegative; see (A1). Secondly, it is not hard to see that $L=\sum_{i=1}^{d} L_{i}$ may be written as $L=\widetilde{L}+r$, where $\widetilde{L}$ is a purely differential operator (without a zero-order potential term) and $r$ is the potential term. We impose that $\sum_{i=1}^{d} q_{i}^{2}$ is controlled pointwise from above by a constant times $r$, namely $\sum_{i=1}^{d} q_{i}^{2} \leq K \cdot r$ for some $K \geqslant 0$; see (A2). In several cases we will consider, we can take $K=1$ or $K=0$. In particular if $q_{1}=\cdots=q_{d}=0$ then the bound (A2) holds with $K=0$. When 0 is not an $L^{2}$ eigenvalue of $L$, our main result can be summarized as follows.

Main result (informal). Set $p^{*}=\max (p, p /(p-1))$. Then the vectorial Riesz transform $\boldsymbol{R} f=$ ( $R_{1} f, \ldots, R_{d} f$ ) with $R_{i}$ given by (1-2) satisfies the bounds

$$
\|\boldsymbol{R} f\|_{L^{p}(X, \mu)} \leqslant 24(1+\sqrt{K})\left(p^{*}-1\right)\|f\|_{L^{p}(X, \mu)}, \quad 1<p<\infty .
$$

In other words, introducing $\delta f=\left(\delta_{1} f, \ldots, \delta_{d} f\right)$, we have

$$
\||\delta f|\|_{L^{p}(X, \mu)} \leqslant 24(1+\sqrt{K})\left(p^{*}-1\right)\left\|L^{1 / 2} f\right\|_{L^{p}(X, \mu)}, \quad 1<p<\infty .
$$

The rigorous statement of our main result is contained in Theorem 1. In order to prove it we need some extra technical assumptions. For the sake of clarity of the presentation we decided to concentrate on the case of orthogonal expansions, when each of the operators $L_{i}=\delta_{i}^{*} \delta_{i}+a_{i}$ has a decomposition in terms of an orthonormal basis. Our precise setting is described in detail in Section 2. We follow the approach of [Nowak and Stempak 2006]; in fact the present paper may be thought of as an $L^{p}$ counterpart for a large part of the $L^{2}$ results from that paper. Adding the technical assumptions (T1), (T2), and (T3) to the crucial assumptions (A1) and (A2), we state our main result, Theorem 1, in Section 3. In all the cases we will consider, the projection $\Pi$ appearing in Theorem 1 is the identity operator or has its $L^{p}$ norm bounded by 2 for all $1 \leq p \leq \infty$. Moreover, we have $\Pi=I$ if and only if 0 is not an $L^{2}$ eigenvalue of $L$.

From Theorem 1 we obtain several new dimension-free bounds on $L^{p}, 1<p<\infty$, for vectors of Riesz transforms connected with classical multidimensional orthogonal expansions. For more details we refer to the examples in Section 5. For instance, in Section 5.3 we obtain the dimension-free boundedness for the vector of Riesz transforms in the case of Jacobi polynomial expansions. This answers a question left open in [Nowak and Sjögren 2008]. Moreover, the approach we present gives a unified way to treat dimension-free estimates for vectors of Riesz transforms. In most of the previous cases, separate papers were written for each of the classical orthogonal expansions. More unified approaches were recently presented by Forzani, Sasso, and Scotto in [Forzani et al. 2015] and by the author in [Wróbel 2014]. However, these papers treat only dimension-free estimates for scalar Riesz transforms and not for the vector of Riesz transforms.

Let us remark that Theorem 1 formally cannot be applied to some cases where the crucial assumptions on $v_{i}$ and $r$ continue to hold. This is true when $L$ has a purely continuous spectrum, for instance for the classical Riesz transforms on $\mathbb{R}^{d}$ (when $v_{i}=0$ and $r=0$ ). However, it is not difficult to modify the proof of Theorem 1 so that it remains valid for the classical Riesz transforms. We believe that a similar procedure can be applied to other cases outside the scope of Theorem 1, as long as the crucial assumptions (A1) and (A2) are satisfied.

We deduce Theorem 1 from a bilinear embedding theorem (see Theorem 3) together with a bilinear formula (see Proposition 2). The main tool that is used to prove Theorem 3 is the Bellman function technique. This method was introduced to harmonic analysis by Nazarov, Treil, and Volberg [Nazarov et al. 1999]. Before that paper, Bellman functions appeared implicitly in [Burkholder 1984; 1988; 1991]. The proof of Theorem 3 is presented in Section 4 and is based on subtle properties of a particular Bellman function. This approach was devised by Dragičević and Volberg [2006; 2011; 2012]. Carbonaro and Dragičević [2013; 2016a; 2016b; 2017] developed the method further. The approach from [Carbonaro and Dragičević 2013] was recently adapted by Mauceri and Spinelli [2015] to the case of the Laguerre operator. Our paper generalizes simultaneously [Dragičević and Volberg 2012] (as we admit a nonnegative potential $r$ ) and [Dragičević and Volberg 2006; Mauceri and Spinelli 2015] (as we consider general $p_{i}$ in $\left.\delta_{i}=p_{i} \partial_{x_{i}}+q_{i}\right)$.

In some applications of the Bellman function method, the authors needed to prove dimension-free bounds on $L^{p}$ for certain spectral multipliers related to the considered operators; see [Dragičević and Volberg 2006; 2012] for such a situation. In other papers mentioned in the previous paragraph they needed to consider operators acting on differential forms; see [Carbonaro and Dragičević 2013; Mauceri and Spinelli 2015]. One of the merits of our approach is that we avoid using both general spectral multipliers and differential forms. This is achieved by means of the bilinear formula from Proposition 2. This formula relates the Riesz transform $R_{i}$ with an integral where only $\delta_{i}$ and two kinds of semigroups (one for $L$ and one for $L+v_{i}$ ) are present; see (3-1).

For the sake of simplicity we use a Bellman function with real entries in Section 4. Thus our main results, Theorems 1 and 3, apply to real-valued functions. Of course they can be easily extended to complex-valued functions with the constants being twice as large. One may improve the estimates further by using a Bellman function with complex arguments, as was done in [Dragičević and Volberg 2006; 2011; 2012].

Notation. We finish this section by introducing the general notation used in the paper. By $\mathbb{N}$ we denote the set of nonnegative integers. For $N \in \mathbb{N}$ and $Y$ being an open subset of $\mathbb{R}^{N}$, the symbol $C^{n}(Y), n \in \mathbb{N}$, denotes the space of real-valued functions which have continuous partial derivatives in $Y$ up to the order $n$. In particular $C^{0}(Y)=C(Y)$ denotes the space of continuous functions on $Y$ equipped with the supremum norm. By $C^{\infty}(Y)$ we mean the space of infinitely differentiable functions on $Y$. Whenever we say that $v$ is a measure on $Y$ we mean that $v$ is a Borel measure on $Y$. The symbols $\nabla f$ and Hess $f$ stand for the gradient and the Hessian of a function $f: \mathbb{R}^{N} \rightarrow \mathbb{R}$. For $a, b \in \mathbb{R}^{N}$, we denote by $\langle a, b\rangle$ the inner product on $\mathbb{R}^{N}$ and set $|a|^{2}=\langle a, a\rangle$. The actual $N$ should be clear from the context (in fact we always have $N \in\{1, d, d+1\})$. For $p \in(1, \infty)$ we set

$$
p^{*}=\max \left(p, \frac{p}{p-1}\right) .
$$

## 2. Preliminaries

All the functions we consider are real-valued. Our notations will closely follow that of [Nowak and Stempak 2006].

For $i=1, \ldots, d$, let $X_{i}$ be the real line $\mathbb{R}$, an open half-line in $\mathbb{R}$ or an open interval in $\mathbb{R}$ of the form

$$
X_{i}=\left(\sigma_{i}, \Sigma_{i}\right), \quad \text { where }-\infty \leqslant \sigma_{i}<\Sigma_{i} \leqslant \infty .
$$

Consider the measure spaces $\left(X_{i}, \mathcal{B}_{i}, \mu_{i}\right)$, where $\mathcal{B}_{i}$ denotes the $\sigma$-algebra of Borel subsets of $X_{i}$ and $\mu_{i}$ is a Borel measure on $X_{i}$. We impose that $d \mu_{i}\left(x_{i}\right)=w_{i}\left(x_{i}\right) d x_{i}$, where $w_{i}$ is a positive $C^{\infty}$ function on $X_{i}$. Note that in [Nowak and Stempak 2006] the authors assumed that $X_{1}=\cdots=X_{d}$; this is, however, not needed in our paper. Throughout the article we let

$$
X=X_{1} \times \cdots \times X_{d}, \quad \mu=\mu_{1} \otimes \cdots \otimes \mu_{d}
$$

and abbreviate

$$
L^{p}:=L^{p}(X, \mu), \quad\|\cdot\|_{p}=\|\cdot\|_{L^{p}}, \quad \text { and } \quad\|\cdot\|_{p \rightarrow p}=\|\cdot\|_{L^{p} \rightarrow L^{p}}
$$

This notation is also used for vector-valued functions. Namely, if $\boldsymbol{g}=\left(g_{1}, \ldots, g_{N}\right): X \rightarrow \mathbb{R}^{N}$ for some $N \in \mathbb{N}$, then

$$
\|\boldsymbol{g}\|_{p}=\left(\int_{X}|\boldsymbol{g}(x)|^{p} d \mu(x)\right)^{1 / p}, \quad \text { with }|\boldsymbol{g}(x)|=\left(\sum_{i=1}^{N}\left|g_{i}(x)\right|^{2}\right)^{1 / 2} .
$$

We shall also write $\langle f, g\rangle_{L^{2}}$ for $\langle f, g\rangle_{L^{2}(X, \mu)}$.
Let $\delta_{i}, i=1, \ldots, d$, be the operators acting on $C_{c}^{\infty}\left(X_{i}\right)$ functions via

$$
\delta_{i}=p_{i} \partial_{x_{i}}+q_{i} .
$$

Here $p_{i}$ and $q_{i}$ are real-valued functions on $X_{i}$, with $p_{i}, q_{i} \in C^{\infty}\left(X_{i}\right)$. We assume that $p_{i}\left(x_{i}\right) \neq 0$ for $x_{i} \in X_{i}$. We shall also denote by $p$ and $q$ the exponents of $L^{p}$ and $L^{q}$ spaces. This will not lead to any confusion as the functions $p_{i}$ and $q_{i}$ will always appear with the index $i=1, \ldots, d$.

Let $\delta_{i}^{*}$ be the formal adjoint of $\delta_{i}$ with respect to the inner product on $L^{2}\left(X_{i}, \mu_{i}\right)$; i.e.,

$$
\delta_{i}^{*} f=-\frac{1}{w_{i}} \partial_{x_{i}}\left(p_{i} w_{i} f\right)+q_{i} f, \quad f \in C_{c}^{\infty}\left(X_{i}\right)
$$

A simple calculation, see [Nowak and Stempak 2006, p. 683], shows that the commutator

$$
\begin{equation*}
\left[\delta_{i}, \delta_{i}^{*}\right]=\delta_{i} \delta_{i}^{*}-\delta_{i}^{*} \delta_{i}=p_{i}\left(2 q_{i}^{\prime}-\left(p_{i} \frac{w_{i}^{\prime}}{w_{i}}\right)^{\prime}-p_{i}^{\prime \prime}\right)=: v_{i} \tag{2-1}
\end{equation*}
$$

is a locally integrable function ( 0 -order operator). Most of the assumptions made in this section are of a technical nature. The first of the two assumptions that are crucial to our results is

$$
\begin{equation*}
\text { the functions } v_{i}, i=1, \ldots, d \text {, are nonnegative. } \tag{A1}
\end{equation*}
$$

The property (A1) has been (explicitly or implicitly) instrumental for establishing the main results in [Harboure et al. 2004; Mauceri and Spinelli 2015; Nowak and Sjögren 2008; Stempak and Wróbel 2013]. It is also explicitly stated by Forzani, Sasso, and Scotto as Assumption H1(c) in [Forzani et al. 2015].

For a scalar $a_{i} \geq 0$ we let $L_{i}$ and $L$ be given on $C_{c}^{\infty}(X)$ by

$$
L_{i}:=\delta_{i}^{*} \delta_{i}+a_{i}, \quad L=\sum_{i=1}^{d} L_{i}
$$

Here each $L_{i}$ can be considered to act either on $C_{c}^{\infty}\left(X_{i}\right)$ or on $C_{c}^{\infty}(X)$; thus the definition of $L$ makes sense. Note that both $L_{i}$ and $L$ are symmetric on $C_{c}^{\infty}(X)$ with respect to the inner product on $L^{2}$. We assume that for each $i=1, \ldots, d$, there is an orthonormal basis $\left\{\varphi_{k_{i}}^{i}\right\}_{k_{i} \in \mathbb{N}}$ which consists of $L^{2}\left(X_{i}, \mu_{i}\right)$ eigenvectors of $L_{i}$ that correspond to nonnegative eigenvalues $\left\{\lambda_{k_{i}}^{i}\right\}_{k_{i} \in \mathbb{N}}$; i.e.,

$$
L_{i} \varphi_{k_{i}}^{i}=\lambda_{k_{i}}^{i} \varphi_{k_{i}}^{i}
$$

Then, it must be that $\lambda_{k_{i}} \geq a_{i}$ for $k_{i} \in \mathbb{N}$ and $i=1, \ldots, d$. We require that the sequence $\left\{\lambda_{k_{i}}^{i}\right\}_{k_{i} \in \mathbb{N}}$ is strictly increasing and that $\lim _{k_{i} \rightarrow \infty} \lambda_{k_{i}}^{i}=\infty$. Note that our assumptions on $p_{i}, q_{i}$, and $w_{i}$ imply that $L_{i}$ is hypoelliptic. Therefore we have $\varphi_{k_{i}}^{i} \in C^{\infty}\left(X_{i}\right)$. Setting, for $k=\left(k_{1}, \ldots, k_{d}\right) \in \mathbb{N}^{d}$,

$$
\begin{equation*}
\varphi_{k}=\varphi_{k_{1}}^{1} \otimes \cdots \otimes \varphi_{k_{d}}^{d} \tag{2-2}
\end{equation*}
$$

we obtain an orthonormal basis of eigenvectors on $L^{2}$ for the operator $L=L_{1}+\cdots+L_{d}$. The eigenvalue corresponding to $\varphi_{k}$ is

$$
\lambda_{k}:=\lambda_{k_{1}}^{1}+\cdots+\lambda_{k_{d}}^{d}
$$

so that $L \varphi_{k}=\lambda_{k} \varphi_{k}$. We consider the self-adjoint extension of $L$ (still denoted by the same symbol) given by

$$
L f=\sum_{k \in \mathbb{N}^{d}} \lambda_{k}\left\langle f, \varphi_{k}\right\rangle_{L^{2}} \varphi_{k}
$$

on the domain

$$
\operatorname{Dom}(L)=\left\{f \in L^{2}: \sum_{k \in \mathbb{N}^{d}}\left|\lambda_{k}\right|^{2}\left|\left\langle f, \varphi_{k}\right\rangle_{L^{2}}\right|^{2}<\infty\right\}
$$

We assume that the eigenfunctions $\varphi_{k_{i}}^{i}, i=1, \ldots, d$, are such that

$$
\begin{equation*}
\left\langle\delta_{i} \varphi_{k_{i}}^{i}, \delta_{i} \varphi_{m_{i}}^{i}\right\rangle_{L^{2}\left(X_{i}, \mu_{i}\right)}=\left\langle\delta_{i}^{*} \delta_{i} \varphi_{k_{i}}^{i}, \varphi_{m_{i}}^{i}\right\rangle_{L^{2}\left(X_{i}, \mu_{i}\right)} \tag{T1}
\end{equation*}
$$

for $i=1, \ldots, d$, and $k_{i}, m_{i} \in \mathbb{N}$; see [Nowak and Stempak 2006, (2.8)]. The condition (T1) implies that the functions

$$
\begin{equation*}
\delta_{i} \varphi_{k}=\varphi_{k_{1}}^{1} \otimes \cdots \otimes \delta_{i} \varphi_{k_{i}}^{i} \otimes \cdots \otimes \varphi_{k_{d}}^{d} \tag{2-3}
\end{equation*}
$$

are pairwise orthogonal on $L^{2}$ and

$$
\left\langle\delta_{i} \varphi_{k}, \delta_{i} \varphi_{k}\right\rangle_{L^{2}}=\lambda_{k_{i}}^{i}-a_{i} ;
$$

see [Nowak and Stempak 2006, Lemmas 5 and 6]. Moreover, since $\varphi_{k} \in C^{\infty}(X)$, we also see that $\delta_{i} \varphi_{k} \in C^{\infty}(X)$.

We remark that our assumptions differ slightly from those in [Nowak and Stempak 2006]. Namely, we assume that the coefficients $p_{i}, q_{i}$, and the weight $w_{i}$ are $C^{\infty}$ functions, whereas they considered $p_{i}, q_{i}, w_{i}$ that possessed only a finite order of smoothness. The smoothness of these functions is in fact needed to easily conclude that $L_{i}$ is hypoelliptic and that $\varphi_{k} \in C^{\infty}(X)$, which is an issue that was overlooked ${ }^{\dagger}$ in [Nowak and Stempak 2006].

We also impose a boundary condition on the functions $\varphi_{k_{i}}^{i}$ and $\delta_{i} \varphi_{k_{i}}^{i}$. Namely, we require that for each $i=1, \ldots, d$, if $z_{i} \in\left\{\sigma_{i}, \Sigma_{i}\right\}$, then

$$
\begin{array}{r}
\lim _{x_{i} \rightarrow z_{i}}\left[\left(1+\left|\varphi_{k_{i}}^{i}\right|^{s_{1}}+\left|\delta_{i} \varphi_{k_{i}}^{i}\right|^{s_{2}}\right)\left(p_{i}^{2} w_{i} \partial_{x_{i}} \varphi_{k_{i}}^{i}\right)\right]\left(x_{i}\right)=0, \\
\lim _{x_{i} \rightarrow z_{i}}\left[\left(1+\left|\varphi_{k_{i}}^{i}\right|^{s_{1}}+\mid \delta_{i} \varphi_{k_{i}}^{i}{ }^{S_{2}}\right)\left(p_{i}^{2} w_{i} \partial_{x_{i}} \delta_{i} \varphi_{k_{i}}^{i}\right)\right]\left(x_{i}\right)=0 \tag{T2}
\end{array}
$$

for all $k_{i} \in \mathbb{N}$ and $s_{1}, s_{2}>0$. Condition (T2) is close to Assumption H1(a) from [Forzani et al. 2015]. Observe that the term $\left|\varphi_{k_{i}}^{i}\right|^{s_{1}}+\left|\delta_{i} \varphi_{k_{i}}^{i}\right|^{s_{2}}$ in (T2) is significant only when the functions $\varphi_{k_{i}}^{i}$ and $\delta_{i} \varphi_{k_{i}}^{i}$ are unbounded on $X_{i}$.

Let

$$
A=a_{1}+\cdots+a_{d}, \quad \Lambda_{0}=\lambda_{0}^{1}+\cdots+\lambda_{0}^{d} .
$$

Then $\Lambda_{0}$ is the smallest eigenvalue of $L$. We set

$$
\mathbb{N}_{\Lambda}^{d}= \begin{cases}\mathbb{N}^{d}, & \Lambda_{0}>0, \\ \mathbb{N}^{d} \backslash\{(0, \ldots, 0)\}, & \Lambda_{0}=0\end{cases}
$$

and define

$$
\Pi f=\sum_{k \in \mathbb{N}_{\Lambda}^{d}}\left\langle f, \varphi_{k}\right\rangle_{L^{2}} \varphi_{k}
$$

Then in the case $\Lambda_{0}>0$ we have $\Pi=I$, while in the case $\Lambda_{0}=0$ the operator $\Pi$ is the projection onto the orthogonal complement of the vector $\varphi_{(0, \ldots, 0)}$. The Riesz transforms studied in this paper are formally

[^1]of the form
$$
R_{i}:=\delta_{i} L^{-1 / 2} \Pi \text {, }
$$
while the rigorous definition of $R_{i}$ is
$$
R_{i} f=\sum_{k \in \mathbb{N}_{\Lambda}^{d}} \lambda_{k}^{-1 / 2}\left\langle f, \varphi_{k}\right\rangle_{L^{2}} \delta_{i} \varphi_{k}
$$

In many of the considered cases, $\Pi \equiv I$ so that $R_{i}=\delta_{i} L^{-1 / 2}$.
It was proved in [Nowak and Stempak 2006, Proposition 1] that the vector of Riesz transforms

$$
\boldsymbol{R} f=\left(R_{1} f, \ldots, R_{d} f\right)
$$

satisfies

$$
\|\boldsymbol{R} f\|_{2 \rightarrow 2} \leqslant\|f\|_{2} .
$$

The main goal of this paper is to prove similar estimates for $p$ in place of 2 . We aim for these estimates to be dimension-free and linear in $p^{*}$. More precisely, we shall prove that for $1<p<\infty$ it holds that

$$
\|\boldsymbol{R} f\|_{p \rightarrow p} \leqslant C\left(p^{*}-1\right)\|f\|_{p}
$$

Here $C$ is a constant that is independent of both $p$ and the dimension $d$.
To state and prove our main results we need several auxiliary objects. Firstly, we let

$$
\begin{equation*}
\mathfrak{d}_{i}=p_{i} \partial_{x_{i}} \tag{2-4}
\end{equation*}
$$

That is, $\mathfrak{d}_{i}$ is the "differential" part of $\delta_{i}$. In many (though not all) of our applications we will have $q_{i} \equiv 0$ and thus $\delta_{i} \equiv \mathfrak{d}_{i}$. The formal adjoint of $\mathfrak{d}_{i}$ on $L^{2}\left(X_{i}, \mu_{i}\right)$ is

$$
\begin{equation*}
\mathfrak{d}_{i}^{*} f=-\frac{1}{w_{i}} \partial_{x_{i}}\left(p_{i} w_{i} f\right), \quad f \in C_{c}^{\infty}\left(X_{i}\right) \tag{2-5}
\end{equation*}
$$

A computation shows that $L_{i}=\mathfrak{d}_{i}^{*} \mathfrak{d}_{i}+r_{i}$, with

$$
\begin{equation*}
r_{i}=a_{i}+\left(q_{i}^{2}-p_{i} q_{i}^{\prime}-p_{i}^{\prime} q_{i}-p_{i} q_{i} \frac{w_{i}^{\prime}}{w_{i}}\right) \tag{2-6}
\end{equation*}
$$

We shall also need

$$
\widetilde{L}:=\sum_{i=1}^{d} \mathfrak{d}_{i}^{*} \mathfrak{d}_{i}=L-r, \quad \text { where } r:=\sum_{i=1}^{d} r_{i}
$$

Then $\widetilde{L}$ is the potential-free component of $L$ and the potential $r$ is a locally integrable function on $X$. We assume that

$$
\begin{equation*}
\text { there is a constant } K \geqslant 0 \text { such that } \sum_{i=1}^{d} q_{i}^{2}\left(x_{i}\right) \leq K \cdot r(x) \tag{A2}
\end{equation*}
$$

for all $x \in X$. This is our second (and last) crucial assumption. In many of our examples we shall have $q_{1}=\cdots=q_{d}=0$ and thus $r=A$ and (A2) holding with $K=0$.

Next we define

$$
M_{i}:=\sum_{j \neq i} \delta_{j}^{*} \delta_{j}+\delta_{i} \delta_{i}^{*}=L+\left[\delta_{i}, \delta_{i}^{*}\right]=L+v_{i}
$$

see [Nowak and Stempak 2006, (5.1)], and set

$$
c_{k}^{i}=\left\|\delta_{i} \varphi_{k}\right\|_{2}^{-1}
$$

if $\delta_{i} \varphi_{k} \neq 0$ and $c_{k}^{i}=0$ in the other case. Then $\left\{c_{k}^{i} \delta_{i} \varphi_{k}\right\}_{k \in \mathbb{N}^{d}}$ (excluding those $c_{k}^{i} \delta_{i} \varphi_{k}$ which vanish) is an orthonormal system of eigenvectors of $M_{i}$ such that $M_{i}\left(c_{k}^{i} \delta_{i} \varphi_{k}\right)$ equals $\lambda_{k} c_{k}^{i} \delta_{i} \varphi_{k}$.

We define

$$
\mathcal{D}=\operatorname{lin}\left\{\varphi_{k}: k \in \mathbb{N}^{d}\right\}, \quad \mathcal{D}_{i}=\delta_{i}[\mathcal{D}]=\operatorname{lin}\left\{\delta_{i} \varphi_{k}: k \in \mathbb{N}^{d}\right\}
$$

and make the technical assumption that

$$
\begin{equation*}
\text { both } \mathcal{D} \text { and } \mathcal{D}_{i}, i=1, \ldots, d \text {, are dense subspaces of } L^{p}, 1 \leq p<\infty \tag{T3}
\end{equation*}
$$

In most of our applications the condition (T3) will follow from [Forzani et al. 2015, Lemma 7.5], which is itself a consequence of [Berg and Christensen 1981, Theorem 5].

Lemma 1 [Forzani et al. 2015, Lemma 7.5]. Assume that $v$ is a measure on $X$ such that, for some $\varepsilon>0$, we have

$$
\int_{X} \exp \left(\varepsilon \sum_{i=1}^{d}\left|y_{i}\right|\right) d \nu(y)<\infty
$$

Then, for each $1 \leqslant p<\infty$, multivariable polynomials on $X$ are dense in $L^{p}(X, v)$.
In what follows we consider the self-adjoint extension of $M_{i}$ given by

$$
M_{i} f=\sum_{k \in \mathbb{N}^{d}} \lambda_{k}\left\langle f, c_{k}^{i} \delta_{i} \varphi_{k}\right\rangle_{L^{2}} c_{k}^{i} \delta_{i} \varphi_{k},
$$

on the domain

$$
\operatorname{Dom}\left(M_{i}\right)=\left\{f \in L^{2}: \sum_{k \in \mathbb{N}^{d}}\left|\lambda_{k}\right|^{2}\left|\left\langle f, c_{k}^{i} \delta_{i} \varphi_{k}\right\rangle_{L^{2}}\right|^{2}<\infty\right\}
$$

Keeping the symbol $M_{i}$ for this self-adjoint extension is a slight abuse of notation, which however will not lead to any confusion. Finally, we shall need the semigroups

$$
P_{t}:=e^{-t L^{1 / 2}} \quad \text { and } \quad Q_{t}^{i}:=e^{-t M_{i}^{1 / 2}}
$$

These are formally defined on $L^{2}$ as

$$
P_{t} f=\sum_{k \in \mathbb{N}^{d}} e^{-t \lambda_{k}^{1 / 2}}\left\langle f, \varphi_{k}\right\rangle_{L^{2}} \varphi_{k}, \quad Q_{t}^{i} f=\sum_{k \in \mathbb{N}^{d}} e^{-t \lambda_{k}^{1 / 2}}\left\langle f, c_{k}^{i} \delta_{i} \varphi_{k}\right\rangle_{L^{2}} c_{k}^{i} \delta_{i} \varphi_{k}
$$

Note that for $t>0$ we have $P_{t}[\mathcal{D}] \subseteq \mathcal{D}$ and $Q_{t}^{i}\left[\mathcal{D}_{i}\right] \subseteq \mathcal{D}_{i}, i=1, \ldots, d$.

## 3. General results for Riesz transforms

Recall that we are in the setting of the previous section. In particular the assumptions (A1), (A2), and the technical assumptions (T1), (T2), (T3), are in force. The following is the main result of our paper.

Theorem 1. For each $1<p<\infty$ we have

$$
\|\boldsymbol{R} f\|_{p} \leqslant 24(1+\sqrt{K})\left(p^{*}-1\right)\|\Pi f\|_{L^{p}}, \quad f \in L^{p}
$$

Remark. In all the examples we consider in Section 5 the projection $\Pi$ satisfies $\|\Pi\|_{p \rightarrow p} \leqslant 2,1 \leqslant p \leqslant \infty$. In fact in many of the examples $\Pi$ equals the identity operator.

In order to prove Theorem 1 we need two ingredients. The first of these ingredients is a bilinear formula that relates the Riesz transform with an integral in which both $P_{t}$ and $Q_{t}^{i}$ are present.
Proposition 2. Let $i=1, \ldots, d$. Then the formula

$$
\begin{equation*}
\left\langle R_{i} f, g\right\rangle_{L^{2}}=-4 \int_{0}^{\infty}\left\langle\delta_{i} P_{t} \Pi f, \partial_{t} Q_{t}^{i} g\right\rangle_{L^{2}} t d t \tag{3-1}
\end{equation*}
$$

holds for $f \in \mathcal{D}$ and $g \in \mathcal{D}_{i}$.
Before proving the proposition let us make two remarks.
Remark 1. Formulas similar to (3-1) were proved before, though, depending on the context, they may have involved spectral multipliers of the operator $L$. However, treating these spectral multipliers appropriately was achieved with variable success. A way of avoiding multipliers was first devised in [Carbonaro and Dragičević 2013] for Riesz transforms on manifolds. In such a setting, the above formula is a special case of the identity (3) there. The approach in [Carbonaro and Dragičević 2013] was adapted in [Mauceri and Spinelli 2015] to the case of Hodge-Laguerre operators. In the case of Laguerre polynomial expansions (see Section 5.2) the formula (3-1) is a special case of [Mauceri and Spinelli 2015, (5.1)]. We note that both in [Carbonaro and Dragičević 2013] and [Mauceri and Spinelli 2015] the authors needed to consider the Riesz transform as well as the formula (3-1) for differential forms; this is not needed in our approach.
Remark 2. Note that if the operators $\delta_{i}$ and $\delta_{i}^{*}$ commute, then $Q_{t}^{i}=P_{t}$ and the formula (3-1) can be formally obtained via the spectral theorem. The problem is that often these operators do not commute. A way to overcome this noncommutativity problem was devised by Nowak and Stempak [2013]. They introduced a symmetrization $T_{i}$ of $\delta_{i}$ that does commute with its adjoint; in fact $T_{i}^{*}=-T_{i}$. This symmetrization is defined on $L^{2}(\tilde{X})$, where

$$
\tilde{X}=\left(X_{1} \cup\left(-X_{1}\right)\right) \times \cdots \times\left(X_{d} \cup\left(-X_{d}\right)\right) .
$$

Set $T=-\sum_{i=1}^{d} T_{i}^{2}$ and let $S_{t}=e^{-t T^{1 / 2}}$. The formula (3-1) for $T_{i}$ is then formally

$$
\begin{equation*}
\left\langle T_{i} T^{-1 / 2} f, g\right\rangle_{L^{2}(\tilde{X})}=-4 \int_{0}^{\infty}\left\langle T_{i} S_{t} f, \partial_{t} S_{t} g\right\rangle_{L^{2}(\tilde{X})} t d t \tag{3-2}
\end{equation*}
$$

This leads to a proof of (3-1) different from the one presented in our paper. Namely, a computation shows that applying (3-2) to functions $f: \widetilde{X} \rightarrow \mathbb{R}$ and $g: \widetilde{X} \rightarrow \mathbb{R}$, which are both even in all the variables, we arrive at (3-1).

Proof of Proposition 2. We start with proving (3-1) for $f=\varphi_{k}$ and $g=\delta_{i} \varphi_{n}$, with some $k, n \in \mathbb{N}^{d}$. If $k=0$ and $\Lambda_{0}=0$ then both sides of (3-1) vanish. Thus we can assume that $\lambda_{k}>0$. A computation shows that

$$
\left\langle\delta_{i} L^{-1 / 2} f, g\right\rangle_{L^{2}}=\lambda_{k}^{-1 / 2}\left\langle\delta_{i} f, g\right\rangle_{L^{2}}
$$

and

$$
\begin{aligned}
-4 \int_{0}^{\infty}\left\langle\delta_{i} P_{t} f, \partial_{t} Q_{t}^{i} g\right\rangle_{L^{2}} t d t & =-4 \int_{0}^{\infty}\left\langle e^{-t \lambda_{k}^{1 / 2}} \delta_{i} f,-\lambda_{n}^{1 / 2} e^{-t \lambda_{n}^{1 / 2}} g\right\rangle_{L^{2}} t d t \\
& =4 \lambda_{n}^{1 / 2} \int_{0}^{\infty} e^{-t\left(\lambda_{k}^{1 / 2}+\lambda_{n}^{1 / 2}\right)} t d t \cdot\left\langle\delta_{i} f, g\right\rangle_{L^{2}}=\frac{4 \lambda_{n}^{1 / 2}}{\left(\lambda_{k}^{1 / 2}+\lambda_{n}^{1 / 2}\right)^{2}} \cdot\left\langle\delta_{i} f, g\right\rangle_{L^{2}} ;
\end{aligned}
$$

hence

$$
\begin{equation*}
\left\langle\delta_{i} L^{-1 / 2} f, g\right\rangle+4 \int_{0}^{\infty}\left\langle\delta_{i} P_{t} f, \partial_{t} Q_{t}^{i} g\right\rangle_{L^{2}} t d t=\left(\lambda_{k}^{-1 / 2}-\frac{4 \lambda_{n}^{1 / 2}}{\left(\lambda_{k}^{1 / 2}+\lambda_{n}^{1 / 2}\right)^{2}}\right) \cdot\left\langle\delta_{i} f, g\right\rangle_{L^{2}} \tag{3-3}
\end{equation*}
$$

Now $\delta_{i} f$ is also an $L^{2}$ eigenvector for $M_{i}$ corresponding to the eigenvalue $\lambda_{k}$. Consequently, since eigenspaces for $M_{i}$ corresponding to different eigenvalues are orthogonal, $\left\langle\delta_{i} f, g\right\rangle$ is nonzero only if $\lambda_{n}=\lambda_{k}$. Coming back to (3-3) we obtain (3-1) for $f=\varphi_{k}$ and $g=\delta_{i} \varphi_{n}$.

Finally, by linearity (3-1) holds also for $f \in \mathcal{D}$ and $g \in \mathcal{D}_{i}$.
The second ingredient we need to prove Theorem 1 is a bilinear embedding, as was the case in [Carbonaro and Dragičević 2013; Dragičević and Volberg 2006; 2012; Mauceri and Spinelli 2014; 2015]. For $N \in \mathbb{N}$ (the cases interesting to us being $N=1$ and $N=d$ ) we take $F=\left(f_{1}, \ldots, f_{N}\right): X \times(0, \infty) \rightarrow \mathbb{R}^{N}$ and set

$$
\begin{equation*}
|F|_{*}^{2}:=r|F|^{2}+\left|\partial_{t} F\right|^{2}+\sum_{i=1}^{d}\left|\mathfrak{d}_{i} F\right|^{2} \tag{3-4}
\end{equation*}
$$

The absolute values $|\cdot|$ in (3-4) denote the Euclidean norms on $\mathbb{R}^{N}$ of the vectors $F(x, t), \partial_{t} F(x, t)=$ $\left(\partial_{t} f_{1}(x, t), \ldots, \partial_{t} f_{N}(x, t)\right)$, and $\mathfrak{d}_{i} F(x, t)=\left(\mathfrak{d}_{i} f_{1}(x, t), \ldots, \mathfrak{d}_{i} f_{N}(x, t)\right)$, where $(x, t) \in X \times(0, \infty)$. Below we only state our bilinear embedding. The proof of it is presented in the next section.

Theorem 3. Let $f: X \rightarrow \mathbb{R}$ and $\boldsymbol{g}=\left(g_{1}, \ldots, g_{d}\right): X^{d} \rightarrow \mathbb{R}^{d}$ and assume that $f \in \mathcal{D}$ and $g_{i} \in \mathcal{D}_{i}$ for $i=1, \ldots, d$. Define

$$
F(x, t)=P_{t} \Pi f(x) \quad \text { and } \quad G(x, t)=Q_{t} g=\left(Q_{t}^{1} g_{1}, \ldots, Q_{t}^{d} g_{d}\right)
$$

Then

$$
\begin{equation*}
\int_{0}^{\infty} \int_{X}|F(x, t)|_{*}|G(x, t)|_{*} d \mu(x) t d t \leqslant 6\left(p^{*}-1\right)\|\Pi f\|_{p}\|\boldsymbol{g}\|_{q} \tag{3-5}
\end{equation*}
$$

Remark. The theorem can be slightly generalized, at least at a formal level. Namely in Theorem 3, we do not need that $v_{i}=\left[\delta_{i}, \delta_{i}^{*}\right]$. It is enough to have any $v_{i} \geq 0$ and take $Q_{t}=e^{-t M_{i}}$ with $M_{i}=L+v_{i}$.

Our main theorem is an immediate corollary of Proposition 2 and Theorem 3.

Proof of Theorem 1. It is enough to prove that for each $f \in L^{p}$ and $g_{i} \in L^{q}, i=1, \ldots, d$, the absolute value of $\sum_{i=1}^{d}\left\langle R_{i} f, g_{i}\right\rangle$ does not exceed

$$
24(1+\sqrt{K})\left(p^{*}-1\right)\|\Pi f\|_{p}\left\|\left(\sum_{i=1}^{d}\left|g_{i}\right|^{2}\right)^{1 / 2}\right\|_{q}
$$

A density argument based on the assumption (T3) allows us to take $f \in \mathcal{D}$ and $g_{i} \in \mathcal{D}_{i}, i=1, \ldots, d$. From Proposition 2 we have

$$
-\frac{1}{4}\left\langle R_{i} f, g_{i}\right\rangle_{L^{2}}=\int_{0}^{\infty}\left\langle\mathfrak{o}_{i} P_{t} \Pi f, \partial_{t} Q_{t}^{i} g_{i}\right\rangle_{L^{2}} t d t+\int_{0}^{\infty}\left\langle q_{i} P_{t} \Pi f, \partial_{t} Q_{t}^{i} g_{i}\right\rangle_{L^{2}} t d t
$$

and thus, assumption (A2) gives

$$
\begin{aligned}
\left|\sum_{i=1}^{d}\left\langle R_{i} f, g_{i}\right\rangle_{L^{2}}\right| & \leqslant 4 \int_{0}^{\infty} \int_{X}\left(\left(\sum_{i=1}^{d}\left|\mathfrak{d}_{i} P_{t} \Pi f(x)\right|^{2}\right)^{1 / 2}+\sqrt{K} \sqrt{r(x)}\left|P_{t} \Pi f(x)\right|\right)|G(x, t)|_{*} d \mu(x) t d t \\
& \leq 4(1+\sqrt{K}) \int_{0}^{\infty} \int_{X}|F(x, t)|_{*}|G(x, t)|_{*} d \mu(x) t d t
\end{aligned}
$$

Now, Theorem 3 completes the proof.

## 4. Bilinear embedding theorem

This section is devoted to the proof of our embedding theorem, Theorem 3. We shall follow closely the reasoning from [Carbonaro and Dragičević 2013; Mauceri and Spinelli 2015].
4.1. The Bellman function. Before proceeding to the proof of Theorem 3 we need to introduce its most important ingredient: the Bellman function.

Choose $p \geqslant 2$. Let $q=p /(p-1)$,

$$
\gamma=\gamma(p)=\frac{1}{8} q(q-1),
$$

and define $\beta_{p}:[0, \infty)^{2} \rightarrow[0, \infty]$ by

$$
\beta_{p}\left(s_{1}, s_{2}\right)=s_{1}^{p}+s_{2}^{q}+\gamma \begin{cases}s_{1}^{2} s_{2}^{2-q}, & s_{1}^{p} \leqslant s_{2}^{q} \\ (2 / p) s_{1}^{p}+(2 / q-1) s_{2}^{q}, & s_{1}^{p} \geqslant s_{2}^{q}\end{cases}
$$

For $m=\left(m_{1}, m_{2}\right) \in \mathbb{N}^{2}$, the Nazarov-Treil Bellman function corresponding to $p, m$ is the function

$$
B=B_{p, m}: \mathbb{R}^{m_{1}} \times \mathbb{R}^{m_{2}} \rightarrow[0, \infty)
$$

given, for any $\zeta \in \mathbb{R}^{m_{1}}$ and $\eta \in \mathbb{R}^{m_{2}}$, by

$$
B_{p, m}(\zeta, \eta)=\frac{1}{2} \beta_{p}(|\zeta|,|\eta|)
$$

The function $B$ originated in an article by F. Nazarov and S. Treil [1996]. It was employed (and simplified) in [Carbonaro and Dragičević 2013; 2017; Dragičević and Volberg 2006; 2011; 2012]. Note that $B$ is
$C^{1}\left(\mathbb{R}^{m_{1}+m_{2}}\right)$ and is $C^{2}$ everywhere except on the set

$$
\left\{(\zeta, \eta) \in \mathbb{R}^{m_{1}} \times \mathbb{R}^{m_{2}}: \eta=0 \text { or }|\zeta|^{p}=|\eta|^{q}\right\} .
$$

To remedy the nonsmoothness of $B$ we consider the regularization

$$
B_{\kappa, p, m}=B_{\kappa}:=B *_{\mathbb{R}^{m_{1}}+m_{2}} \psi_{\kappa},
$$

where

$$
\psi(x)=c_{m} e^{-1 /\left(1-|x|^{2}\right)} \chi_{B^{m_{1}+m_{2}}}(x) \quad \text { and } \quad \psi_{\kappa}(x)=\frac{1}{\kappa^{m_{1}+m_{2}}} \psi(x / \kappa),
$$

with $c_{m}$ such that $\int_{\mathbb{R}^{m_{1}+m_{2}}} \psi_{\kappa}(x) d x=1$. Here $\chi_{B^{m_{1}+m_{2}}}$ stands for the characteristic function of the ( $m_{1}+m_{2}$ )-dimensional Euclidean ball centered at the origin and of radius 1 . Since both $B$ and $\psi_{\kappa}$ are biradial also $B_{\kappa}$ is biradial. Hence, there is $\beta_{\kappa}=\beta_{\kappa, p}$ acting from $[0, \infty)^{2}$ to $\mathbb{R}$ such that

$$
B_{\kappa}(\zeta, \eta)=\frac{1}{2} \beta_{\kappa}(|\zeta|,|\eta|), \quad \zeta \in \mathbb{R}^{m_{1}}, \eta \in \mathbb{R}^{m_{2}} .
$$

We shall need some properties of $\beta_{\kappa}$ and $B_{\kappa}$ that were essentially proved in [Carbonaro and Dragičević 2013; Dragičević and Volberg 2012; Mauceri and Spinelli 2014; 2015].

Proposition 4. Let $\kappa \in(0,1)$. Then, for $s_{i}>0, i=1,2$, we have
(i) $0 \leqslant \beta_{\kappa}\left(s_{1}, s_{2}\right) \leqslant(1+\gamma(p))\left(\left(s_{1}+\kappa\right)^{p}+\left(s_{2}+\kappa\right)^{q}\right)$,
(ii) $0 \leqslant \partial_{s_{1}} \beta_{\kappa}(s) \leqslant C_{p} \max \left(\left(s_{1}+\kappa\right)^{p-1}, s_{2}+\kappa\right)$ and $0 \leqslant \partial_{s_{2}} \beta_{\kappa}(s) \leqslant C_{p}\left(s_{2}+\kappa\right)^{q-1}$, with $C_{p}$ a positive constant.

The function $B_{\kappa}$ belongs to $C^{\infty}\left(\mathbb{R}^{m_{1}+m_{2}}\right)$, and for any $\xi=(\zeta, \eta) \in \mathbb{R}^{m_{1}+m_{2}}$ there exists a positive $\tau_{\kappa}=\tau_{\kappa}(|\zeta|,|\eta|)$ such that for $\omega=\left(\omega_{1}, \omega_{2}\right) \in \mathbb{R}^{m_{1}+m_{2}}$ we have
(iii) $\left\langle\operatorname{Hess}\left(B_{\kappa}\right)(\xi) \omega, \omega\right\rangle \geqslant \frac{1}{2} \gamma(p)\left(\tau_{\kappa}\left|\omega_{1}\right|^{2}+\tau_{\kappa}^{-1}\left|\omega_{2}\right|^{2}\right)$.

Moreover, there is a continuous function $E_{\kappa}: \mathbb{R}^{m_{1}+m_{2}} \rightarrow \mathbb{R}$ for which
(iv) $\left\langle\left(\nabla B_{\kappa}\right)(\xi), \xi\right\rangle \geqslant \frac{1}{2} \gamma(p)\left(\tau_{\kappa}|\zeta|^{2}+\tau_{\kappa}^{-1}|\eta|^{2}\right)-\kappa E_{\kappa}(\xi)$,
(v) $\left|E_{\kappa}(\xi)\right| \leq C_{m, p}\left(|\zeta|^{p-1}+|\eta|+|\eta|^{q-1}+\kappa^{q-1}\right)$.

Proof (sketch). Let $\tau=\tau(|\zeta|,|\eta|)$ be the function from [Carbonaro and Dragičević 2013, Theorem 3] and define $\tau_{\kappa}=\tau *_{\mathbb{R}^{d+1}} \psi_{\kappa}$. With exactly this $\tau_{\kappa}$, items (i), (ii), and (iii) were proved in [Mauceri and Spinelli 2014, Proposition 6.3].

Let

$$
E_{\kappa}(\xi)=-\int_{\mathbb{R}^{m_{1}+m_{2}}}\langle\nabla B(\xi-\kappa s), s\rangle \psi_{\kappa}(s) d s, \quad \xi \in \mathbb{R}^{m_{1}+m_{2}}
$$

see [Dragičević and Volberg 2012, (2.10)]. Item (iv) (with these $\tau_{\kappa}$ and $E_{\kappa}$ ) follows from [Dragičević and Volberg 2012, Theorem 4(iii')], together with the observation from [Carbonaro and Dragičević 2013; Dragičević and Volberg 2012] that

$$
\left(\tau * \psi_{\kappa}\right)(\xi)\left(\tau^{-1} * \psi_{\kappa}\right)(\xi) \geq\left(\int_{\mathbb{R}^{d+1}}\left(\tau(y) \psi_{\kappa}(x-y)\right)^{1 / 2}\left(\tau^{-1}(y) \psi_{\kappa}(x-y)\right)^{1 / 2} d y\right)^{1 / 2}=1
$$

Item (v) is proved in [Dragičević and Volberg 2012, p. 207]. Note that, our Bellman function $B_{\kappa}$ coincides with $-\frac{1}{2} Q_{\kappa}$ from that paper (when $Q_{\kappa}$ is restricted to real arguments).

We remark that in [Dragičević and Volberg 2012, Theorem 4 (iii')] a stronger statement is proved with an additional negative term $-B_{\kappa}(\xi)$ on the left-hand side of (iv).
4.2. Proof of Theorem 3. Define $u: X \times(0, \infty) \rightarrow \mathbb{R} \times \mathbb{R}^{d}$ by

$$
u=u(x, t)=\left(P_{t} \Pi f(x), Q_{t} g(x)\right)=\left(P_{t} \Pi f(x), Q_{t}^{1} g_{1}(x), \ldots, Q_{t}^{d} g_{d}(x)\right)
$$

Assume first that $p \geq 2$ and set

$$
b_{\kappa}=B_{\kappa} \circ u: X \times(0, \infty) \rightarrow[0, \infty)
$$

Here $B_{\kappa}=B_{\kappa, d, p}$ is the Bellman function from Proposition 4 with $m_{1}=1$ and $m_{2}=d$. For each $i=1, \ldots, d$, we fix a sequence $\left\{\sigma_{i}^{n}\right\}_{n \in \mathbb{N}}$ which converges to $\sigma_{i}$, and a sequence $\left\{\Sigma_{i}^{n}\right\}_{n \in \mathbb{N}}$ which converges to $\Sigma_{i}$. We also impose that $\sigma_{i}<\sigma_{i}^{n}<\Sigma_{i}^{n}<\Sigma_{i}$ for $i=1, \ldots, d, n \in \mathbb{N}$. Defining

$$
X_{i}^{n}=\left[\sigma_{i}^{n}, \Sigma_{i}^{n}\right] \quad \text { and } \quad X_{n}=X_{1}^{n} \times \cdots \times X_{d}^{n}
$$

where $n \in \mathbb{N}$, we see that $\left\{X_{n}\right\}_{n \in \mathbb{N}}$ is an increasing family of compact subsets of $X$ such that $X=\bigcup_{n} X_{n}$. We shall estimate the integral

$$
\begin{equation*}
I(n, \varepsilon):=\int_{0}^{\infty} \int_{X_{n}}\left(\partial_{t}^{2}-\widetilde{L}\right)\left(b_{\kappa(n)}\right)(x, t) d \mu(x) t e^{-\varepsilon t} d t \tag{4-1}
\end{equation*}
$$

from below and above and then, first let $n \rightarrow \infty$ and then $\varepsilon \rightarrow 0^{+}$. Here $\kappa(n)$ is a small quantity depending on $n$ which will be determined in the proof. Since $X_{n}$ is compact, $f \in \mathcal{D}$ and $g_{i} \in \mathcal{D}_{i}, i=1, \ldots, d$, the integral (4-1) is in fact absolutely convergent. In what follows we will often briefly write $\kappa$ instead of $\kappa(n)$.

The lower estimate of (4-1) for $\boldsymbol{p} \geq \mathbf{2}$. The key result here is Proposition 5 below. Its proof hinges on the assumption (A1).
Proposition 5. For $x \in X$ and $t>0$ it holds that

$$
\begin{equation*}
\left(\left(\partial_{t}^{2}-\widetilde{L}\right) b_{\kappa}\right)(x, t) \geqslant \gamma|F(x, t)|_{*}|G(x, t)|_{*}-\kappa r(x) E_{\kappa}(u(x, t)) . \tag{4-2}
\end{equation*}
$$

Proof. Set $\mathfrak{d}_{0}:=\partial_{t}$. To justify (4-2) we shall need the pointwise equality

$$
\begin{equation*}
\left(\partial_{t}^{2}-\widetilde{L}\right) b_{\kappa}=r\left\langle\nabla B_{\kappa}(u), u\right\rangle+\sum_{i=1}^{d} v_{i} \cdot\left(\partial_{\eta_{i}} B_{\kappa}(u) \cdot Q_{t}^{i} g_{i}\right)+\sum_{i=0}^{d}\left\langle\operatorname{Hess}\left(B_{\kappa}\right)\left(\mathfrak{d}_{i} u\right), \mathfrak{d}_{i} u\right\rangle . \tag{4-3}
\end{equation*}
$$

We first we focus on proving (4-3).
From the chain rule we have $\mathfrak{d}_{i} b_{\kappa}=p_{i}\left\langle\nabla B_{\kappa}(u), \partial_{x_{i}} u\right\rangle$. Moreover, a computation shows that, for $i=1, \ldots, d$,

$$
\mathfrak{d}_{i}^{*}=-p_{i} \partial_{x_{i}}-p_{i} \frac{w_{i}^{\prime}}{w_{i}}-p_{i}^{\prime} \quad \text { and } \quad \mathfrak{d}_{i}^{*} \mathfrak{d}_{i}=-p_{i}^{2} \partial_{x_{i}}^{2}-\left(p_{i} \frac{w_{i}^{\prime}}{w_{i}}+2 p_{i}^{\prime}\right) p_{i} \partial_{x_{i}}
$$

Consequently, applying once again the chain rule we obtain for $i=0, \ldots, d$,

$$
\begin{aligned}
\mathfrak{d}_{i}^{*} \mathfrak{d}_{i} b_{\kappa} & =-p_{i} \partial_{x_{i}}\left(p_{i}\left\langle\nabla B_{\kappa}(u), \partial_{x_{i}} u\right\rangle\right)-\left(p_{i} \frac{w_{i}^{\prime}}{w_{i}}+p_{i}^{\prime}\right) p_{i}\left\langle\nabla B_{\kappa}(u), \partial_{x_{i}} u\right\rangle \\
& =-p_{i}^{2} \partial_{x_{i}}\left(\left\langle\nabla B_{\kappa}(u), \partial_{x_{i}} u\right\rangle\right)-p_{i} p_{i}^{\prime}\left\langle\nabla B_{\kappa}(u), \partial_{x_{i}} u\right\rangle-\left(p_{i} \frac{w_{i}^{\prime}}{w_{i}}+p_{i}^{\prime}\right) p_{i}\left\langle\nabla B_{\kappa}(u), \partial_{x_{i}} u\right\rangle \\
& =\left\langle\nabla B_{\kappa}(u),-p_{i}^{2} \partial_{x_{i}}^{2} u-\left(p_{i} \frac{w_{i}^{\prime}}{w_{i}}+2 p_{i}^{\prime}\right) p_{i} \partial_{x_{i}} u\right\rangle-p_{i}^{2}\left\langle\operatorname{Hess}\left(B_{\kappa}\right)\left(\partial_{x_{i}} u\right), \partial_{x_{i}} u\right\rangle \\
& =\left\langle\nabla B_{\kappa}(u), \mathfrak{d}_{i}^{*} \mathfrak{d}_{i} u\right\rangle-\left\langle\operatorname{Hess}\left(B_{\kappa}\right)\left(\mathfrak{d}_{i} u\right), \mathfrak{d}_{i} u\right\rangle .
\end{aligned}
$$

Now, summing the above formula in $i=0, \ldots, d$, we obtain

$$
\begin{equation*}
\left(\mathfrak{d}_{0}^{2}-\widetilde{L}\right) b_{\kappa}=\left\langle\nabla B_{\kappa}(u),\left(\mathfrak{d}_{0}^{2}-\widetilde{L}\right) u\right\rangle+\sum_{i=0}^{d}\left\langle\operatorname{Hess}\left(B_{\kappa}\right)(u)\left(\mathfrak{d}_{i} u\right), \mathfrak{d}_{i} u\right\rangle . \tag{4-4}
\end{equation*}
$$

The formula (4-4) implies (4-3). Indeed we have

$$
\left(\partial_{t}^{2}-L\right) u=\left(\left(\partial_{t}^{2}-L\right) P_{t} f,\left(\partial_{t}^{2}-L\right) Q_{t} \boldsymbol{g}\right),
$$

where

$$
\left(\partial_{t}^{2}-L\right) P_{t} f=0
$$

and

$$
\left(\partial_{t}^{2}-L\right) Q_{t} \boldsymbol{g}=\left(\left(\partial_{t}^{2}-L\right) Q_{t}^{1} g_{1}, \ldots,\left(\partial_{t}^{2}-L\right) Q_{t}^{d} g_{d}\right)
$$

Moreover,

$$
\left(\partial_{t}^{2}-L\right) Q_{t}^{i} g_{i}=\left(\partial_{t}^{2}-M_{i}\right) Q_{t}^{i} g_{i}+v_{i} \cdot Q_{t}^{i} g_{i}=v_{i} \cdot Q_{t}^{i} g_{i}
$$

and using (4-4) the equation (4-3) follows.
Having demonstrated (4-3) we pass to the proof of (4-2). Proposition 4(ii) implies ( $\left.\partial_{\eta_{i}} B_{\kappa}(u) \cdot Q_{t}^{i} g_{i}\right) \geq 0$. Thus (4-3) together with the assumption (A1) produce

$$
\begin{equation*}
\left(\partial_{t}^{2}-\widetilde{L}\right) b_{\kappa} \geqslant r\left\langle\nabla B_{\kappa}(u), u\right\rangle+\sum_{i=0}^{d}\left\langle\operatorname{Hess}\left(B_{\kappa}\right)\left(\mathfrak{d}_{i} u\right), \mathfrak{o}_{i} u\right\rangle . \tag{4-5}
\end{equation*}
$$

Finally, (4-2) is a consequence of (4-5), items (iii) and (iv) from Proposition 4, and the inequality between the arithmetic and geometric mean.

Coming back to the proof of the lower estimate in (4-1) we now take $\{\kappa(n)\}_{n \in \mathbb{N}}$ such that $|\kappa(n)| \leqslant 1$, $\lim _{n} \kappa(n)=0$ and

$$
\begin{equation*}
|\kappa(n)|^{1 / 2} \int_{X_{n}}\left|r(x) E_{\kappa(n)}(u(x, t))\right| d \mu(x) \leqslant 1 . \tag{4-6}
\end{equation*}
$$

To see that such a sequence exists we use Proposition 4(v) and the fact that $P_{t} f \in \mathcal{D}$ and $Q_{t}^{i} g_{i} \in \mathcal{D}_{i}$ (hence also $P_{t} f \in C^{\infty}(X)$ and $Q_{t}^{i} g_{i} \in C^{\infty}(X)$ ). Next, (4-2), together with (4-6), leads to

$$
\liminf _{n \rightarrow \infty} I(n, \varepsilon) \geq \gamma \int_{0}^{\infty} \int_{X}|F(x, t)|_{*}|G(x, t)|_{*} d \mu(x) t e^{-\varepsilon t} d t,
$$

and, consequently, by the monotone convergence theorem

$$
\begin{equation*}
\liminf _{\varepsilon \rightarrow 0^{+}} \liminf _{n \rightarrow \infty} I(n, \varepsilon) \geq \gamma(p) \int_{0}^{\infty} \int_{X}|F(x, t)|_{*}|G(x, t)|_{*} d \mu(x) t d t \tag{4-7}
\end{equation*}
$$

This is our lower estimate of (4-1).
The upper estimate of (4-1) for $\boldsymbol{p} \geq \mathbf{2}$. The main ingredients here are the technical assumptions (T2) and (T3). We split the integral in (4-1) as

$$
\begin{aligned}
I(n, \varepsilon) & =I_{1}(n, \varepsilon)-I_{2}(n, \varepsilon) \\
& :=\int_{0}^{\infty} \int_{X_{n}} \partial_{t}^{2}\left(b_{\kappa(n)}\right)(x, t) d \mu(x) t e^{-\varepsilon t} d t-\int_{0}^{\infty} \int_{X_{n}} \widetilde{L}\left(b_{\kappa(n)}\right)(x, t) d \mu(x) t e^{-\varepsilon t} d t .
\end{aligned}
$$

First we prove that

$$
\begin{equation*}
\lim _{n \rightarrow \infty} I_{2}(n, \varepsilon)=0 \tag{4-8}
\end{equation*}
$$

To see this we recall that $\widetilde{L}=\sum_{i=1}^{d} \mathfrak{d}_{i}^{*} \mathfrak{d}_{i}$ with $\mathfrak{d}_{i}$ given by (2-4) and $\mathfrak{d}_{i}^{*}$ being the formal adjoint of $\mathfrak{d}_{i}$ on $L^{2}$. Then,

$$
I_{2}(n, \varepsilon)=\sum_{i=1}^{d} I_{2}^{i}(n, \varepsilon):=\sum_{i=1}^{d} \int_{0}^{\infty} \int_{X_{n}}\left(\mathfrak{d}_{i}^{*} \mathfrak{d}_{i}\right)\left(b_{\kappa(n)}\right)(x, t) d \mu(x) t e^{-\varepsilon t} d t
$$

and it is enough to prove that each of the integrals $I_{2}^{i}(n, \varepsilon)$ goes to 0 as $n \rightarrow \infty$. As the reasoning is symmetric in $i=1, \ldots, d$, we present it only for $I_{2}^{1}(n, \varepsilon)$. Define

$$
X^{(1)}=X_{2} \times \cdots \times X_{d}, \quad x^{(1)}=\left(x_{2}, \ldots, x_{d}\right), \quad \text { and } \quad \mu^{(1)}=\mu_{2} \otimes \cdots \otimes \mu_{d} .
$$

Formula (2-5) together with integration by parts in the $x_{1}$-variable produces

$$
I_{2}^{1}(n, \varepsilon)=\int_{0}^{\infty} \int_{X^{(1)}}\left(\left(p_{1}^{2} w_{1} \partial_{x_{1}} b_{\kappa}\right)\left(\Sigma_{1}^{n}, x^{(1)}\right)-\left(p_{1}^{2} w_{1} \partial_{x_{1}} b_{\kappa}\right)\left(\sigma_{1}^{n}, x^{(1)}\right)\right) d \mu^{(1)}\left(x^{1}\right) t e^{-\varepsilon t} d t
$$

Call $z_{1}^{n}$ either of the quantities $\sigma_{1}^{n}$ or $\Sigma_{1}^{n}$. Then the chain rule gives

$$
\begin{align*}
\left(p_{1}^{2} w_{1} \partial_{x_{1}} b_{\kappa}\right)\left(z_{1}^{n}, x^{(1)}\right)= & p_{1}^{2}\left(z_{1}^{n}\right) w_{1}\left(z_{1}^{n}\right) \partial_{x_{1}} P_{t} f\left(z_{1}^{n}, x^{(1)}\right) \partial_{\zeta} B_{\kappa}\left(P_{t} f\left(z_{1}^{n}, x^{(1)}\right), Q_{t} \boldsymbol{g}\left(z_{1}^{n}, x^{(1)}\right)\right) \\
& +p_{1}^{2}\left(z_{1}^{n}\right) w_{1}\left(z_{1}^{n}\right)\left\langle\partial_{x_{1}} Q_{t} \boldsymbol{g}\left(z_{1}^{n}, x^{(1)}\right), \nabla_{\eta} B_{\kappa}\left(P_{t} f\left(z_{1}^{n}, x^{(1)}\right), Q_{t} \boldsymbol{g}\left(z_{1}^{n}, x^{(1)}\right)\right)\right\rangle . \tag{4-9}
\end{align*}
$$

Since $f \in \mathcal{D}$ and $g_{i} \in \mathcal{D}_{i}$ we have that $P_{t} f \in \mathcal{D}$ and $Q_{t} \boldsymbol{g} \in \mathcal{D}_{1} \otimes \cdots \otimes \mathcal{D}_{d}$. Recall that $\varphi_{k}$ is defined by (2-2), while $\delta_{i} \varphi_{k}, i=1, \ldots, d$, are given by (2-3). Now, Proposition 4(ii) implies

$$
\begin{equation*}
\left|\nabla_{\zeta, \eta} B_{\kappa}(\zeta, \eta)\right| \leqslant C_{p, q}\left(|\zeta|^{p-1}+|\eta|^{q-1}+|\eta|+\kappa^{q-1}\right) \tag{4-10}
\end{equation*}
$$

Therefore, since $|\kappa(n)| \leqslant 1$, a calculation based on (4-9) together with the assumptions (T2), (T3), and Hölder's inequality produces $\lim _{n} I_{2}^{1}(n, \varepsilon)=0$.

Now we focus on $I_{1}(n, \varepsilon)$. Since $f \in \mathcal{D}, g_{i} \in \mathcal{D}_{i}, i=1, \ldots, d, B_{\kappa} \in C^{\infty}\left(\mathbb{R}^{d+1}\right)$ and we integrate over $x \in X_{n}$, the double integral is absolutely convergent. Thus Fubini's theorem gives

$$
I_{1}(n, \varepsilon)=\int_{X_{n}} \int_{0}^{\infty} \partial_{t}^{2}\left(b_{\kappa(n)}\right)(x, t) t e^{-\varepsilon t} d t d \mu(x)
$$

Integrating by parts in the inner integral twice we obtain

$$
\begin{aligned}
I_{1}(n, \varepsilon) & =-\int_{X_{n}} \int_{0}^{\infty} \partial_{t}\left(b_{\kappa(n)}\right)(x, t)(1-\varepsilon t) e^{-\varepsilon t} d t d \mu(x) \\
& =\int_{X_{n}} b_{\kappa(n)}(x, 0) d \mu(x)+\varepsilon^{2} \int_{X_{n}} \int_{0}^{\infty} b_{\kappa(n)}(x, t) t e^{-\varepsilon t} d t d \mu(x)-2 \varepsilon \int_{X_{n}} \int_{0}^{\infty} b_{\kappa(n)}(x, t) e^{-\varepsilon t} d t d \mu(x) \\
& \leqslant \int_{X_{n}} b_{\kappa(n)}(x, 0) d \mu(x)+\varepsilon^{2} \int_{X_{n}} \int_{0}^{\infty} b_{\kappa(n)}(x, t) t e^{-\varepsilon t} d t d \mu(x) \\
& :=I_{1}^{1}(n)+I_{1}^{2}(n, \varepsilon) .
\end{aligned}
$$

In the first two equalities above we neglected the boundary terms by using the chain rule together with (4-10).

First we treat $I_{1}^{2}(n, \varepsilon)$. Proposition 4(i) gives

$$
I_{1}^{2}(n, \varepsilon) \leqslant \varepsilon^{2} C_{p} \int_{X_{n}} \int_{0}^{\infty}\left(\left|P_{t} \Pi f(x)\right|^{p}+\left|Q_{t} g(x)\right|^{q}+\max \left(\kappa(n)^{p}, \kappa(n)^{q}\right)\right) t e^{-\varepsilon t} d t d \mu(x) .
$$

Take $\kappa(n)$ which satisfies (4-6) and

$$
\begin{equation*}
\max \left(\kappa(n)^{p-1 / 2}, \kappa(n)^{q-1 / 2}\right) \mu\left(X_{n}\right) \leqslant 1 . \tag{4-11}
\end{equation*}
$$

Then, since $f \in \mathcal{D}$ and $g_{i} \in \mathcal{D}_{i}, i=1, \ldots, d$, we have

$$
\limsup _{n \rightarrow \infty} I_{1}^{2}(n, \varepsilon) \leqslant \varepsilon^{2} C_{p} \int_{X} \int_{0}^{\infty}\left|P_{t} \Pi f(x)\right|^{p}+\left|Q_{t} \boldsymbol{g}(x)\right|^{q} t d t d \mu(x) \leqslant C(p, f, g) \varepsilon^{2}
$$

and, consequently,

$$
\begin{equation*}
\limsup _{\varepsilon \rightarrow 0^{+}} \limsup _{n \rightarrow \infty} I_{1}^{2}(n, \varepsilon)=0 \tag{4-12}
\end{equation*}
$$

Coming back to $I_{1}^{1}(n)$ we use Proposition 4(i) to estimate

$$
I_{1}^{1}(n) \leqslant \frac{1}{2}(1+\gamma) \int_{X_{n}}(|\Pi f(x)|+\kappa(n))^{p} d \mu(x)+\frac{1}{2}(1+\gamma) \int_{X_{n}}(|\boldsymbol{g}(x)|+\kappa(n))^{q} d \mu(x) .
$$

Now for each $\varepsilon>0$ we split the first integral onto $\int_{|\kappa(n)| \leqslant \varepsilon|\Pi f(x)|}$ and $\int_{|\kappa(n)|>\varepsilon|\Pi f(x)|}$ and the second integral onto $\int_{|\kappa(n)| \leqslant \varepsilon|g(x)|}$ and $\int_{|\kappa(n)|>\varepsilon|g(x)|}$. Then we obtain

$$
\begin{aligned}
& I_{1}^{1}(n) \leqslant \frac{1}{2}(1+\gamma)\left((1+\varepsilon)^{p}\|\Pi f\|_{p}^{p}+(1+\varepsilon)^{q}\|\boldsymbol{g}\|_{q}^{q}\right) \\
& \\
& \quad+\frac{1}{2}(1+\gamma)\left(\left(1+\varepsilon^{-1}\right)^{p} \kappa(n)^{p} \mu\left(X_{n}\right)+\left(1+\varepsilon^{-1}\right)^{q} \kappa(n)^{q} \mu\left(X_{n}\right)\right) .
\end{aligned}
$$

Since $\kappa(n)$ satisfies (4-6) and (4-11) we arrive at

$$
\limsup _{\varepsilon \rightarrow 0^{+}} \limsup _{n \rightarrow \infty} I_{1}^{1}(n) \leqslant \frac{1}{2}(1+\gamma)\left(\|\Pi f\|_{p}^{p}+\|\boldsymbol{g}\|_{q}^{q}\right) .
$$

Recalling (4-8) and (4-12) we thus proved

$$
\begin{equation*}
\limsup _{\varepsilon \rightarrow 0^{+}} \limsup _{n \rightarrow \infty} I(n, \varepsilon) \leqslant \frac{1}{2}(1+\gamma(p))\left(\|\Pi f\|_{p}^{p}+\|\boldsymbol{g}\|_{q}^{q}\right) \tag{4-13}
\end{equation*}
$$

which is the upper estimate of (4-1) we need.
Completion of the proof of the bilinear embedding. Consider first $p \geq 2$. Combining the lower estimate (4-7) and the upper estimate (4-13) we obtain

$$
\begin{equation*}
\int_{0}^{\infty} \int_{X}|F(x, t)|_{*}|G(x, t)|_{*} d \mu(x) t d t \leqslant \frac{1+\gamma(p)}{2 \gamma(p)}\left(\|\Pi f\|_{p}^{p}+\|\boldsymbol{g}\|_{q}^{q}\right) . \tag{4-14}
\end{equation*}
$$

Finally, a polarization argument finishes the proof. More precisely, for $s>0$ we replace $f$ with $s f$ and $\boldsymbol{g}$ with $s^{-1} \boldsymbol{g}$ on both sides of (4-14). Then, the left-hand side is unchanged, while minimizing the right-hand side over $s>0$ we obtain

$$
\begin{equation*}
\int_{0}^{\infty} \int_{X}|F(x, t)|_{*}|G(x, t)|_{*} d \mu(x) t d t \leqslant \frac{1+\gamma(p)}{2 \gamma(p)}\left(\left(\frac{p}{q}\right)^{1 / p}+\left(\frac{q}{p}\right)^{1 / q}\right)\|\Pi f\|_{p}\|\boldsymbol{g}\|_{q} \tag{4-15}
\end{equation*}
$$

Using the above inequality, a calculation leads to (3-5). We sketch the argument below.
Note that for $p \geq 2$ we have $p^{*}=p$ and recall that $\gamma(p)=\frac{1}{8} q(q-1)$. Thus, for $1<q \leq 2$ we obtain

$$
\begin{align*}
\frac{1+\gamma(p)}{2 \gamma(p)}\left(\left(\frac{p}{q}\right)^{1 / p}+\left(\frac{q}{p}\right)^{1 / q}\right) & =\frac{1}{2}(8+q(q-1))(q-1)^{1 / q-1}(p-1) \\
& \leq(q+3)(q-1)^{1 / q-1}\left(p^{*}-1\right) \tag{4-16}
\end{align*}
$$

Setting $s=q-1$ we need to maximize the function $H(s):=(s+4) s^{-s /(s+1)}$ for $s \in(0,1]$. Let

$$
h(s)=\log (s+4)-\frac{s \log s}{s+1}
$$

so that $H(s)=e^{h(s)}$. Then we have

$$
h^{\prime}(s)=\frac{1}{s+4}-\frac{\log s}{(s+1)^{2}}-\frac{1}{s+1} \quad \text { and } \quad h^{\prime \prime}(s)=-\frac{1}{(s+4)^{2}}+\frac{2 \log s}{(s+1)^{3}}+\frac{s-1}{s(s+1)^{2}}
$$

consequently, $h^{\prime \prime}(s)<0$ for $s \in(0,1)$. Observe that $h^{\prime}\left(\frac{7}{20}\right)>0$ and $h^{\prime}\left(\frac{2}{5}\right)<0$. Therefore $h^{\prime}$ has a unique zero inside the interval $\left(\frac{7}{20}, \frac{2}{5}\right)$ and $h$ attains a global maximum there. Obviously, the same is true for $H=e^{h}$. Now it is easy to see that

$$
\max _{7 / 20 \leq s \leq 2 / 5} H(s)<\frac{22}{5} \cdot\left(\frac{7}{20}\right)^{-2 / 7}<6,
$$

and thus also $\sup _{0<s \leq 1} H(s)<6$. Hence, coming back to (4-16) we obtain

$$
\frac{1+\gamma(p)}{2 \gamma(p)}\left(\left(\frac{p}{q}\right)^{1 / p}+\left(\frac{q}{p}\right)^{1 / q}\right) \leq 6\left(p^{*}-1\right) .
$$

In view of (4-15) this implies (3-5) and completes the proof of Theorem 3 for $p \geq 2$.

The proof of Theorem 3 for $p \leq 2$ proceeds analogously once we switch $p$ with $q$ and $P_{t} f$ with $Q_{t} g$ in the definition of $b_{\kappa}$. Namely, we consider $\tilde{b}_{\kappa}(x, t)=\widetilde{B}_{\kappa}\left(Q_{t} \boldsymbol{g}, P_{t} f\right)$, where $\widetilde{B}_{\kappa}(\zeta, \eta)=B_{\kappa, q,(d, 1)}(\zeta, \eta)$, $\zeta \in \mathbb{R}^{d}, \eta \in \mathbb{R}$. Here $B_{\kappa, q,(d, 1)}$ is the function from Proposition 4 with $m_{1}=d$ and $m_{2}=1$. Then we repeat the argument used for $p \geq 2$. The function $\widetilde{B}_{\kappa}$ satisfies items (iii)-(v) of Proposition 4 with $p$ replaced by $q$. Therefore both the lower estimate (4-7) and the upper estimate (4-13) hold with $\gamma(p)$ replaced by $\gamma(q)$.

## 5. Examples

Throughout this section we apply Theorem 1 to the examples of orthogonal systems considered in [Nowak and Stempak 2006, Section 7]. This is possible for all of these systems except for the Fourier-Bessel expansions [Nowak and Stempak 2006, Section 7.8]. In this case the condition (T2) fails. Despite this failure we think that it might be possible to treat also the Fourier-Bessel expansions by the methods of the present paper. It might be also interesting to try to apply the methods of our paper to the Riesz transforms considered in [Nowak and Sjögren 2012] (in the case of Jacobi trigonometric polynomial expansions).

In all of the examples we present, for more details the reader is kindly referred to [Nowak and Stempak 2006, Sections 7.1-7.7]. The formulas for $v_{i}$ and $r=\sum_{i=1}^{d} r_{i}$ in the examples below follow directly from (2-1) and (2-6). Recall that

$$
\mu=\mu_{1} \otimes \cdots \otimes \mu_{d}, \quad X=X_{1} \times \cdots \times X_{d}, \quad L^{p}=L^{p}(X, \mu), \quad\|\cdot\|_{p}=\|\cdot\|_{L^{p}}
$$

and

$$
p^{*}=\max \left(p, \frac{p}{p-1}\right) .
$$

5.1. Ornstein-Uhlenbeck operator: Hermite polynomial expansions. Here we consider

$$
p_{i}=1, \quad q_{i}=0, \quad a_{i}=0, \quad w_{i}\left(x_{i}\right)=\pi^{-1 / 2} e^{-x_{i}^{2}}, \quad d \mu_{i}\left(x_{i}\right)=w_{i}\left(x_{i}\right) d x_{i}
$$

on $X_{i}=\mathbb{R}$. Then

$$
\begin{equation*}
\delta_{i}=\mathfrak{d}_{i}=\partial_{x_{i}}, \quad \delta_{i}^{*}=-\partial_{x_{i}}+2 x_{i}, \quad v_{i}=\left[\delta_{i}, \delta_{i}^{*}\right]=2, \quad r=0, \tag{5-1}
\end{equation*}
$$

and

$$
L=\sum_{i=1}^{d} L_{i}=-\Delta+2\langle x, \nabla\rangle
$$

is the Ornstein-Uhlenbeck operator on $X=\mathbb{R}^{d}$. The operator $L$ is essentially self-adjoint on $C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ with the self-adjoint extension given by

$$
L f=\sum_{k \in \mathbb{N}^{d}}|k|\left\langle f, \widetilde{H}_{k}\right\rangle_{L^{2}} \widetilde{H}_{k} .
$$

In the formula above $|k|=k_{1}+\cdots+k_{d}$, the symbol $L^{2}$ stands for $L^{2}=L^{2}\left(\mathbb{R}^{d}, \mu\right)$, while $\left\{\tilde{H}_{k}\right\}_{k \in \mathbb{N}^{d}}$ is the system of $L^{2}$ normalized Hermite polynomials; see [Nowak and Stempak 2006, Section 7.1; Lebedev 1972, p. 60]. In this section we take

$$
\varphi_{k}=\widetilde{H}_{k}, \quad k \in \mathbb{N}^{d} .
$$

Note that $\mu$ is a probability measure in this setting. The projection $\Pi$ becomes

$$
\Pi f=\sum_{k \in \mathbb{N}^{d}, k \neq 0}\left\langle f, \widetilde{H}_{k}\right\rangle_{L^{2}} \widetilde{H}_{k}, \quad f \in L^{2}
$$

Then

$$
(I-\Pi) f=\left\langle f, \widetilde{H}_{0}\right\rangle_{L^{2}} \widetilde{H}_{0}
$$

and, since $\widetilde{H}_{0}=1$, the operator $I-\Pi$ is the projection onto the constants given by

$$
(I-\Pi) f(x)=\int_{X} f(y) d \mu(y), \quad x \in X
$$

Hence, by Holder's inequality $\|(I-\Pi) f\|_{p} \leqslant\|f\|_{p}$, and, consequently,

$$
\begin{equation*}
\|\Pi f\|_{p} \leqslant 2\|f\|_{p}, \quad 1 \leqslant p \leqslant \infty \tag{5-2}
\end{equation*}
$$

Next

$$
\begin{equation*}
\delta_{i} \widetilde{H}_{k}=\sqrt{2 k_{j}} \tilde{H}_{k-e_{j}}, \tag{5-3}
\end{equation*}
$$

where, by convention $\widetilde{H}_{k-e_{j}}=0$ if $k_{j}=0$. This convention is also used for the examples presented in the next sections. The Riesz transform is defined by

$$
R_{i} f=\sum_{k \in \mathbb{N}^{d}, k \neq 0}\left(\frac{k_{j}}{|k|}\right)^{1 / 2}\left\langle f, \widetilde{H}_{k}\right\rangle_{L^{2}} \widetilde{H}_{k-e_{i}}, \quad f \in L^{2}
$$

Dimension-free estimates for the vector $\boldsymbol{R} f=\left(R_{1} f, \ldots, R_{d} f\right)$ were proved by Meyer [1984]; see also [Gundy 1986; Gutiérrez 1994; Pisier 1988] for different proofs. Later Dragičević and Volberg [2006, Corollary 0.4] found a proof which uses the Bellman function method. The best result in terms of the size of the constants is due to Arcozzi [1998, Corollary 2.4] who proved that $\|\boldsymbol{R} f\|_{p} \leqslant 2\left(p^{*}-1\right)\|f\|_{p}$, $1<p<\infty$. An application of Theorem 1 produces similar, though weaker, bounds.

Theorem 6. Fix $1<p<\infty$. Then, for $f \in L^{p}$ such that $\int_{X} f(y) d \mu(y)=0$, we have

$$
\begin{equation*}
\|\boldsymbol{R} f\|_{p} \leqslant 24\left(p^{*}-1\right)\|f\|_{p} \tag{5-4}
\end{equation*}
$$

Remark. Using (5-2) we may extend the bound (5-4) to all $f \in L^{p}$ with 24 being replaced by 48 .
Proof. We apply Theorem 1. In order to do so we need to check that its assumptions are satisfied.
By ( $5-1$ ) we see that (A1) and (A2) (with $K=0$ ) hold. Condition (T1) is proved by an easy calculation based on integration by parts. The assumption (T2) is also straightforward. Finally, (T3) follows from Lemma 1 and (5-3).

Now, if $\int_{X} f(y) d \mu(y)=0$ then $\Pi f=f$. Thus, an application of Theorem 1 completes the proof.
5.2. Laguerre operator: Laguerre polynomial expansions. Here, for a parameter $\alpha \in(-1, \infty)^{d}$, we consider

$$
p_{i}=\sqrt{x_{i}}, \quad q_{i}=0, \quad a_{i}=0, \quad w_{i}\left(x_{i}\right)=\frac{1}{\Gamma\left(\alpha_{i}+1\right)} x_{i}^{\alpha_{i}} e^{-x_{i}} d x_{i}, \quad d \mu_{i}\left(x_{i}\right)=w_{i}\left(x_{i}\right) d x_{i}
$$

on $X_{i}=(0, \infty)$. Then $\delta_{i}=\mathfrak{d}_{i}=\sqrt{x_{i}} \partial_{x_{i}}$, and thus

$$
\begin{equation*}
\delta_{i}^{*}=-\sqrt{x_{i}} \partial_{x_{i}}-\frac{\alpha_{i}+\frac{1}{2}}{\sqrt{x_{i}}}+\sqrt{x_{i}}, \quad v_{i}=\left[\delta_{i}, \delta_{i}^{*}\right]=\frac{\alpha_{i}+\frac{1}{2}+x_{i}}{2 x_{i}}, \quad r=0 . \tag{5-5}
\end{equation*}
$$

In this case

$$
L=\sum_{i=1}^{d} L_{i}=\sum_{i=1}^{d}-x_{i} \partial_{x_{i}}^{2}-\left(\alpha_{i}+1-x_{i}\right) \partial_{x_{i}}
$$

is the Laguerre operator on $X=(0, \infty)^{d}$. It is symmetric on $C_{c}^{\infty}\left((0, \infty)^{d}\right)$ and has a self-adjoint extension

$$
L f=\sum_{k \in \mathbb{N}^{d}}|k|\left\langle f, \widetilde{L}_{k}^{\alpha}\right\rangle_{L^{2}} \widetilde{L}_{k}^{\alpha} .
$$

Here $L^{2}=L^{2}\left((0, \infty)^{d}, \mu\right)$, while $\left\{\widetilde{L}_{k}^{\alpha}\right\}_{k \in \mathbb{N}^{d}}$ is the system of $L^{2}$ normalized Laguerre polynomials; see [Nowak and Stempak 2006, Section 7.2; Lebedev 1972, p. 76]. These Laguerre polynomials are our functions $\varphi_{k}$ in this section; namely

$$
\varphi_{k}=\widetilde{L}_{k}^{\alpha}, \quad k \in \mathbb{N}^{d} .
$$

Next we have

$$
\begin{equation*}
\delta_{i} \widetilde{L}_{k}^{\alpha}=\sqrt{k_{j}} \sqrt{x_{i}} \widetilde{L}_{k-e_{i}}^{\alpha+e_{i}}, \tag{5-6}
\end{equation*}
$$

while the projection $\Pi$ becomes

$$
\Pi f=\sum_{k \in \mathbb{N}^{d}, k \neq 0}\left\langle f, \widetilde{L}_{k}^{\alpha}\right\rangle_{L^{2}} \widetilde{L}_{k}^{\alpha}, \quad f \in L^{2}
$$

A repetition of the argument from the previous section shows that $\Pi f=f$ if and only if $\int_{X} f(y) d \mu(y)=0$ and

$$
\begin{equation*}
\|\Pi f\|_{p} \leqslant 2\|f\|_{p}, \quad 1 \leqslant p \leqslant \infty \tag{5-7}
\end{equation*}
$$

The Riesz transform is then given by

$$
R_{i} f=\sum_{k \in \mathbb{N}^{d}, k \neq 0}\left(\frac{k_{j}}{|k|}\right)^{1 / 2}\left\langle f, \widetilde{L}_{k}^{\alpha}\right\rangle_{L^{2}} \sqrt{x_{i}} \widetilde{L}_{k-e_{i}}^{\alpha+e_{i}}, \quad f \in L^{2}
$$

Dimension-free bounds for single Riesz transforms $R_{i}$ were first studied by Gutiérrez, Incognito and Torrea [Gutiérrez et al. 2001] for half-integer multi-indices, and generalized ${ }^{\dagger}$ by Nowak [2004] to multiindices $\alpha \in\left[-\frac{1}{2}, \infty\right)^{d}$. Moreover in [Graczyk et al. 2005], Graczyk, Loeb, López, Nowak, and Urbina proved dimension-free estimates on $L^{p}$ for the vector of Riesz-Laguerre transforms and half-integer multi-indices $\alpha$. Recently, the author [Wróbel 2014, Theorem 4.1(b)] obtained dimension-free bounds on $L^{p}$ for scalar Riesz transforms and general parameters $\alpha \in(-1, \infty)^{d}$, while Mauceri and Spinelli [2015, Theorem 5.2] proved a dimension-free bound for the vectorial Riesz transforms $\boldsymbol{R} f=\left(R_{1} f, \ldots, R_{d} f\right)$, and $\alpha \in\left[-\frac{1}{2}, \infty\right)^{d}$. All the bounds mentioned in this paragraph are also independent of the parameter $\alpha$

[^2](appropriately restricted). Moreover, the estimate from [Mauceri and Spinelli 2015, Theorem 5.2] is also linear in $p^{*}$.

By using Theorem 1 we obtain a result which coincides with [Mauceri and Spinelli 2015, Theorem 5.2] in the case of Riesz transforms acting on functions.
Theorem 7. Fix $\alpha \in\left[-\frac{1}{2}, \infty\right)^{d}$ and $1<p<\infty$. Then, for $f \in L^{p}$ which satisfy $\int_{X} f(y) d \mu(y)=0$, we have

$$
\|\boldsymbol{R} f\|_{p} \leqslant 24\left(p^{*}-1\right)\|f\|_{p}
$$

Remark. By (5-7) we have the same bound for general $f \in L^{p}$ with the constant being twice as large.
Proof. We are going to apply Theorem 1, so we need to verify its assumptions.
By (5-5) we see that if $\alpha \in\left[-\frac{1}{2}, \infty\right)^{d}$, then (A1) and (A2) (with $K=0$ ) are satisfied. Moreover, the assumptions (T1) and (T2) follow from a direct calculation. Next, for such $\alpha$ the condition (T3) can be deduced from Lemma 1 together with (5-6).

Now, if $\int_{X} f(y) d \mu(y)=0$ then $\Pi f=f$. Therefore, using Theorem 1 we complete the proof of Theorem 7.
5.3. Jacobi operator: Jacobi polynomial expansions. In this section for parameters $\alpha, \beta \in(-1, \infty)^{d}$ we consider

$$
\begin{gathered}
p_{i}=\sqrt{1-x_{i}^{2}}, \quad q_{i}=0, \quad a_{i}=0 \\
w_{i}\left(x_{i}\right)=\frac{1}{C\left(\alpha_{i}, \beta_{i}\right)}\left(1-x_{i}\right)^{\alpha_{i}}\left(1+x_{i}\right)^{\beta_{i}} d x_{i}, \quad d \mu_{i}\left(x_{i}\right)=w_{i}\left(x_{i}\right) d x_{i}, \quad X_{i}=(-1,1),
\end{gathered}
$$

where $C\left(\alpha_{i}, \beta_{i}\right)$ is such that $\mu_{i}\left(X_{i}\right)=1$. Then $\delta_{i}=\mathfrak{d}_{i}=\sqrt{1-x_{i}^{2}} \partial_{x_{i}}$, and

$$
\begin{gather*}
\delta_{i}^{*}=-\sqrt{1-x_{i}^{2}} \partial_{x_{i}}+\left(\alpha_{i}+\frac{1}{2}\right) \sqrt{\frac{1+x_{i}}{1-x_{i}}}-\left(\beta_{i}+\frac{1}{2}\right) \sqrt{\frac{1-x_{i}}{1+x_{i}}}  \tag{5-8}\\
v_{i}=\left[\delta_{i}, \delta_{i}^{*}\right]=\frac{\alpha_{i}+\frac{1}{2}}{1-x_{i}}+\frac{\beta_{i}+\frac{1}{2}}{1+x_{i}}, \quad r=0 .
\end{gather*}
$$

Here

$$
L=\sum_{i=1}^{d} L_{i}=\sum_{i=1}^{d}-\left(1-x_{i}^{2}\right) \partial_{x_{i}}^{2}-\left(\beta_{i}-\alpha_{i}-\left(\alpha_{i}+\beta_{i}+2\right) x_{i}\right) \partial_{x_{i}}
$$

is the Jacobi operator on $X=(-1,1)^{d}$. Let $L^{2}=L^{2}\left((-1,1)^{d}, \mu\right)$, and denote by $\left\{\widetilde{P}_{k}^{\alpha, \beta}\right\}_{k \in \mathbb{N}^{d}}$ the system of $L^{2}$ normalized Jacobi polynomials; see [Nowak and Stempak 2006, Section 7.1; Szegő 1975, Chapter 4]. These Jacobi polynomials are our functions $\varphi_{k}$ in this section; namely

$$
\varphi_{k}=\widetilde{P}_{k}^{\alpha, \beta}, \quad k \in \mathbb{N} .
$$

The Jacobi operator is symmetric on $C_{c}^{\infty}\left((-1,1)^{d}\right)$ and has a self-adjoint extension

$$
L f=\sum_{k \in \mathbb{N}^{d}} \lambda_{k}\left\langle f, \widetilde{P}_{k}^{\alpha, \beta}\right\rangle_{L^{2}} \widetilde{P}_{k}^{\alpha, \beta},
$$

where $\lambda_{k}=\sum_{i=1}^{d} \lambda_{k_{i}}^{i}$ with $\lambda_{k_{i}}^{i}=k_{i}\left(k_{i}+\alpha_{i}+\beta_{i}+1\right), i=1, \ldots, d$. Similarly to the previous two sections the projection $\Pi$ is

$$
\Pi f=\sum_{k \in \mathbb{N}^{d}, k \neq 0}\left\langle f, \widetilde{P}_{k}^{\alpha, \beta}\right\rangle_{L^{2}} \widetilde{P}_{k}^{\alpha, \beta}, \quad f \in L^{2}
$$

Moreover, $\Pi f=f$ precisely when $\int_{X} f(y) d \mu(y)=0$ and we have

$$
\begin{equation*}
\|\Pi f\|_{p} \leqslant 2\|f\|_{p}, \quad 1 \leqslant p \leqslant \infty \tag{5-9}
\end{equation*}
$$

The action of $\delta_{i}$ on Jacobi polynomials is given by

$$
\begin{equation*}
\delta_{i} \widetilde{P}_{k}^{\alpha, \beta}=\sqrt{k_{i}\left(k_{i}+\alpha_{i}+\beta_{i}+1\right)} \sqrt{1-x_{i}^{2}} \widetilde{P}_{k-e_{i}}^{\alpha+e_{i}, \beta+e_{i}}, \tag{5-10}
\end{equation*}
$$

and the Riesz transform becomes

$$
R_{i} f=\sum_{k \in \mathbb{N}^{d}, k \neq 0}\left(\frac{\lambda_{k_{i}}^{i}}{\lambda_{k}}\right)^{1 / 2}\left\langle f, \widetilde{P}_{k}^{\alpha, \beta}\right\rangle_{L^{2}} \sqrt{1-x_{i}^{2}} \widetilde{P}_{k-e_{i}}^{\alpha+e_{i}, \beta+e_{i}}, \quad f \in L^{2}
$$

Dimension- and parameter-free estimates for single Riesz transforms $R_{i}$ are due to Nowak and Sjögren [2008], who proved them for $\alpha, \beta \in\left[-\frac{1}{2}, \infty\right)^{d}$.

An application of Theorem 1 generalizes [Nowak and Sjögren 2008, Theorem 5.1] to the vectorial Riesz transforms $\boldsymbol{R} f=\left(R_{1} f, \ldots, R_{d} f\right)$. This result is new according to our knowledge. Moreover, we obtain an explicit estimate which is linear in $p^{*}$.

Theorem 8. Fix $\alpha, \beta \in\left[-\frac{1}{2}, \infty\right)^{d}$ and $1<p<\infty$. Then, for $f \in L^{p}$ which satisfy $\int_{X} f(y) d \mu(y)=0$, we have

$$
\begin{equation*}
\|\boldsymbol{R} f\|_{p} \leqslant 24\left(p^{*}-1\right)\|f\|_{p}, \quad f \in L^{p} . \tag{5-11}
\end{equation*}
$$

Remark. As in the previous two sections (5-11) holds for all $f \in L^{p}$ with 48 ( $p^{*}-1$ ) in place of 24 ( $p^{*}-1$ ). This follows from (5-9).

Proof. We are going to apply Theorem 1, so we need to verify its assumptions for parameters $\alpha, \beta \in$ $\left[-\frac{1}{2}, \infty\right)^{d}$.

By (5-8) we see that if $\alpha, \beta \in\left[-\frac{1}{2}, \infty\right)^{d}$, then (A1) and (A2) (with $K=0$ ) are satisfied. Similarly, using (5-10) one can see that, for such $\alpha$ and $\beta$, the conditions (T1) and (T2) also hold. The assumption (T3) follows from Lemma 1 together with (5-10).

Now, since $\int_{X} f(y) d \mu(y)=0$ implies $\Pi f=f$, an application of Theorem 1 completes the proof of Theorem 8.
5.4. Harmonic oscillator: Hermite function expansions. Here we take

$$
p_{i}=1, \quad q_{i}=x_{i}, \quad a_{i}=1, \quad w_{i}\left(x_{i}\right)=1, \quad d \mu_{i}\left(x_{i}\right)=d x_{i}, \quad X_{i}=\mathbb{R}
$$

so that

$$
\begin{equation*}
\delta_{i}=\partial_{x_{i}}+x_{i}, \quad \mathfrak{d}_{i}=\partial_{x_{i}}, \quad \delta_{i}^{*}=-\partial_{x_{i}}+x_{i}, \quad v_{i}=\left[\delta_{i}, \delta_{i}^{*}\right]=2, \quad r(x)=|x|^{2}, \tag{5-12}
\end{equation*}
$$

and $L$ is the harmonic oscillator

$$
L=\sum_{i=1}^{d} L_{i}=-\Delta+|x|^{2}
$$

It is well known that $L$ is essentially self-adjoint on $C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ with the self-adjoint extension given by

$$
L f=\sum_{k \in \mathbb{N}^{d}}(2|k|+d)\left\langle f, h_{k}\right\rangle_{L^{2}} h_{k} .
$$

Here $L^{2}=L^{2}\left(\mathbb{R}^{d}, d x\right)$, while $\left\{h_{k}\right\}_{k \in \mathbb{N}^{d}}$ is the system of $L^{2}$ normalized Hermite functions; see [Nowak and Stempak 2006, Section 7.4]. The functions $h_{k}$ are our $\varphi_{k}$ in this section. They are of the form $h_{k}=h_{k_{1}} \otimes \cdots \otimes h_{k_{d}}$, where

$$
\begin{equation*}
h_{k_{i}}\left(x_{i}\right)=\widetilde{H}_{k_{i}}\left(x_{i}\right) e^{-x_{i}^{2} / 2}, \quad x_{i} \in \mathbb{R}, \tag{5-13}
\end{equation*}
$$

with $\widetilde{H}_{k_{i}}$ being the Hermite polynomial from Section 5.1. Note that as 0 is not an $L^{2}$ eigenvalue of $L$, the projection $\Pi$ equals the identity operator.

Next

$$
\begin{equation*}
\delta_{i} h_{k}=\sqrt{2 k_{j}} h_{k-e_{j}}, \tag{5-14}
\end{equation*}
$$

and thus the Riesz transform is

$$
R_{i} f=\sum_{k \in \mathbb{N}^{d}, k \neq 0}\left(\frac{2 k_{j}}{2|k|+d}\right)^{1 / 2}\left\langle f, h_{k}\right\rangle_{L^{2}} h_{k-e_{i}}, \quad f \in L^{2}
$$

Here dimension-free bounds for the vector of Riesz transforms can be deduced, by means of transference, from the paper of Coulhon, Müller, and Zienkiewicz [Coulhon et al. 1996]; see also [Harboure et al. 2004; Lust-Piquard 2006] for different proofs. Moreover, a dimension-free bound for the vector of Riesz transforms which is additionally linear in $p^{*}$ was proved in [Dragičević and Volberg 2012, Proposition 4].

Using Theorem 1 we are able to obtain a more explicit estimate for the vector $\boldsymbol{R} f$ than in [Dragičević and Volberg 2012]. However, contrary to that paper, our method says nothing about the vector of "adjoint" transforms $\boldsymbol{R}^{*} f=\left(\delta_{1}^{*} L^{-1 / 2} f, \ldots, \delta_{d}^{*} L^{-1 / 2} f\right)$.

Theorem 9. For $1<p<\infty$ we have

$$
\|\boldsymbol{R} f\|_{p} \leqslant 48\left(p^{*}-1\right)\|f\|_{p}, \quad f \in L^{p} .
$$

Proof. We apply Theorem 1. In order to do so we need to check that its assumptions are satisfied.
The equation (5-12) gives (A1) and (A2) with $K=1$. Condition (T1) is straightforward. The assumption (T2) holds since, by (5-13), Hermite functions $h_{k_{i}}$ vanish rapidly at $\pm \infty$. Finally, (T3) follows from (5-14) and the (well-known) density of Hermite functions in $L^{p}, 1 \leqslant p<\infty$.

Thus, an application of Theorem 1 is justified and the proof of Theorem 9 is completed.
5.5. Laguerre operator: Laguerre function expansions of Hermite type. For a parameter $\alpha \in(-1, \infty)^{d}$ we consider

$$
p_{i}=1, \quad q_{i}=x_{i}-\frac{\alpha_{i}+\frac{1}{2}}{x_{i}}, \quad a_{i}=1, \quad w_{i}\left(x_{i}\right)=1, \quad d \mu_{i}\left(x_{i}\right)=d x_{i}, \quad X_{i}=(0, \infty),
$$

so that

$$
\begin{gather*}
\delta_{i}=\partial_{x_{i}}+x_{i}-\frac{\alpha_{i}+\frac{1}{2}}{x_{i}}, \quad \mathfrak{d}_{i}=\partial_{x_{i}}, \quad \delta_{i}^{*}=-\partial_{x_{i}}+x_{i}-\frac{\alpha_{i}+\frac{1}{2}}{x_{i}}, \\
v_{i}=\left[\delta_{i}, \delta_{i}^{*}\right]=2, \quad r(x)=|x|^{2}+\sum_{i=1}^{d} \frac{\alpha_{i}^{2}-\frac{1}{4}}{x_{i}^{2}} . \tag{5-15}
\end{gather*}
$$

Here $L$ is the Laguerre operator

$$
L=\sum_{i=1}^{d} L_{i}=-\Delta+|x|^{2}+\sum_{i=1}^{d} \frac{\alpha_{i}^{2}-\frac{1}{4}}{x_{i}^{2}} .
$$

Then $L$ is symmetric on $C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ and has a self-adjoint extension given by

$$
L f=\sum_{k \in \mathbb{N}^{d}}(4|k|+2 d+2|\alpha|)\left\langle f, \varphi_{k}^{\alpha}\right\rangle_{L^{2}} \varphi_{k}^{\alpha} .
$$

In the above formula we set $|k|=k_{1}+\cdots+k_{d}$ and $|\alpha|=\alpha_{1}+\cdots+\alpha_{d}$; note that $|\alpha|$ may be negative. By $L^{2}$ we mean $L^{2}\left((0, \infty)^{d}, d x\right)$, while $\left\{\varphi_{k}^{\alpha}\right\}_{k \in \mathbb{N}^{d}}$ stands for the system of $L^{2}$ normalized Laguerre functions of Hermite type; see [Nowak and Stempak 2006, Section 7.5]. The functions $\varphi_{k}^{\alpha}$ are the tensor products $\varphi_{k}^{\alpha}=\varphi_{k_{1}}^{\alpha_{1}} \otimes \cdots \otimes \varphi_{k_{d}}^{\alpha_{d}}$ with

$$
\begin{equation*}
\varphi_{k_{i}}^{\alpha_{i}}\left(x_{i}\right)=\sqrt{2} \widetilde{L}_{k_{i}}^{\alpha_{i}}\left(x_{i}^{2}\right) x_{i}^{\alpha_{i}+1 / 2} e^{-x_{i}^{2} / 2}, \quad x_{i}>0, \tag{5-16}
\end{equation*}
$$

and $\widetilde{L}_{k_{i}}^{\alpha_{i}}$ the Laguerre polynomials from Section 5.2. In this section we take

$$
\varphi_{k}=\varphi_{k}^{\alpha} .
$$

As 0 is not an $L^{2}$ eigenvalue of $L$, the projection $\Pi$ equals the identity operator.
Next

$$
\begin{equation*}
\delta_{i} \varphi_{k}^{\alpha}=-2 \sqrt{k_{j}} \varphi_{k-e_{j}}^{\alpha+e_{j}}, \tag{5-17}
\end{equation*}
$$

and thus the Riesz transform is

$$
R_{i} f=-\sum_{k \in \mathbb{N}^{d}, k \neq 0}\left(\frac{4 k_{i}}{4|k|+2|\alpha|+2 d}\right)^{1 / 2}\left\langle f, \varphi_{k}^{\alpha}\right\rangle_{L^{2}} \varphi_{k-e_{j}}^{\alpha+e_{j}}, \quad f \in L^{2}
$$

Dimension-free bounds for single Riesz transforms $R_{i}$ were obtained by Stempak and the author [Stempak and Wróbel 2013, Theorem 5.1] for a certain restricted range of the parameter $\alpha$.

In this section, for $\alpha \in\left(\frac{1}{2}, \infty\right)^{d}$ we define

$$
C(\alpha)=\max _{i=1, \ldots, d} \frac{\alpha_{i}+\frac{1}{2}}{\alpha_{i}-\frac{1}{2}}
$$

By using Theorem 1 we obtain the following strengthening of [Stempak and Wróbel 2013, Theorem 5.1] in the case $\alpha \in\left(\frac{1}{2}, \infty\right)^{d}$.

Theorem 10. Let $\alpha \in\left(\frac{1}{2}, \infty\right)^{d}$. Then, for $1<p<\infty$, we have

$$
\|\boldsymbol{R} f\|_{p} \leqslant 24(1+\sqrt{C(\alpha)})\left(p^{*}-1\right)\|f\|_{p}, \quad f \in L^{p}
$$

Proof. We apply Theorem 1. In order to do so we need to check that its assumptions are satisfied.
The formula (5-15) gives (A1) and (A2) for $\alpha \in\left(\frac{1}{2}, \infty\right)^{d}$ with $K=C(\alpha)$. Conditions (T1) and (T2) follow from (5-16) and (5-17). Finally, (T3) follows from [Nowak 2003, Lemma 5.2] and (5-17).

Thus, an application of Theorem 1 is justified and the proof of Theorem 10 is completed.
5.6. Laguerre operator: Laguerre function expansions of convolution type. For a parameter $\alpha \in$ $(-1, \infty)^{d}$ we consider

$$
\begin{gathered}
p_{i}=1, \quad q_{i}=x_{i}, \quad a_{i}=2 \alpha_{i}+2, \\
w_{i}\left(x_{i}\right)=x_{i}^{2 \alpha_{i}+1}, \quad d \mu_{i}\left(x_{i}\right)=w_{i}\left(x_{i}\right) d x_{i}, \quad X_{i}=(0, \infty),
\end{gathered}
$$

so that

$$
\begin{gather*}
\delta_{i}=\partial_{x_{i}}+x_{i}, \quad \mathfrak{d}_{i}=\partial_{x_{i}}, \quad \delta_{i}^{*}=-\partial_{x_{i}}+x_{i}-\frac{2 \alpha_{i}+1}{x_{i}}, \\
v_{i}=\left[\delta_{i}, \delta_{i}^{*}\right]=2+\frac{2 \alpha+1}{x_{i}^{2}}, \quad r(x)=|x|^{2} \tag{5-18}
\end{gather*}
$$

Here $L$ is the Laguerre operator

$$
L=\sum_{i=1}^{d} L_{i}=-\Delta+|x|^{2}-\sum_{i=1}^{d} \frac{2 \alpha_{i}+1}{x_{i}} \partial_{x_{i}}
$$

Then $L$ is symmetric on $C_{c}^{\infty}\left((0, \infty)^{d}\right)$ and has a self-adjoint extension given by

$$
L f=\sum_{k \in \mathbb{N}^{d}}(4|k|+2 d+2|\alpha|)\left\langle f, \ell_{k}^{\alpha}\right\rangle_{L^{2}} \ell_{k}^{\alpha}
$$

Here $L^{2}=L^{2}\left((0, \infty)^{d}, w(x) d x\right)$, while $\left\{\ell_{k}^{\alpha}\right\}_{k \in \mathbb{N}^{d}}$ is the system of $L^{2}$ normalized Laguerre functions of convolution type; see [Nowak and Stempak 2006, Section 7.6]. The functions $\ell_{k}^{\alpha}$ are of the form $\ell_{k}^{\alpha}=\ell_{k_{1}}^{\alpha_{1}} \otimes \cdots \otimes \ell_{k_{d}}^{\alpha_{d}}$ with

$$
\begin{equation*}
\ell_{k_{i}}^{\alpha_{i}}\left(x_{i}\right)=\sqrt{2} \widetilde{L}_{k_{i}}^{\alpha_{i}}\left(x_{i}^{2}\right) e^{-x_{i}^{2} / 2}, \quad x_{i}>0 \tag{5-19}
\end{equation*}
$$

and $\widetilde{L}_{k_{i}}^{\alpha_{i}}$ the Laguerre polynomials from Section 5.2. In this section we take

$$
\varphi_{k}=\ell_{k}^{\alpha} .
$$

Also here, as 0 is not an $L^{2}$ eigenvalue of $L$, the projection $\Pi$ equals the identity operator.
Next

$$
\begin{equation*}
\delta_{i} \ell_{k}^{\alpha}=-2 \sqrt{k_{i}} x_{i} \ell_{k-e_{i}}^{\alpha+e_{i}} \tag{5-20}
\end{equation*}
$$

and thus the Riesz transform is

$$
R_{i} f=-\sum_{k \in \mathbb{N}^{d}, k \neq 0}\left(\frac{4 k_{i}}{4|k|+2|\alpha|+2 d}\right)^{1 / 2}\left\langle f, \ell_{k}^{\alpha}\right\rangle_{L^{2}} \ell_{k-e_{i}}^{\alpha+e_{i}}, \quad f \in L^{2}
$$

The boundedness of these Riesz transforms on $L^{p}$ was proved by Nowak and Stempak [2007, Theorem 3.4]. Later Nowak and Szarek [2012, Theorem 4.1] enlarged the range of admitted parameters $\alpha$. In both of these papers Calderón-Zygmund theory was used; thus the $L^{p}$ bounds depended on the dimension $d$. Applying Theorem 1 we obtain a dimension-free bound for the vectorial Riesz transform $\boldsymbol{R} f=\left(R_{1} f, \ldots, R_{d} f\right)$.
Theorem 11. Let $\alpha \in\left[-\frac{1}{2}, \infty\right)^{d}$. Then, for $1<p<\infty$, we have

$$
\begin{equation*}
\|\boldsymbol{R} f\|_{p} \leqslant 48\left(p^{*}-1\right)\|f\|_{p}, \quad f \in L^{p} . \tag{5-21}
\end{equation*}
$$

Proof. A continuity argument based on (5-19) and (5-20) shows that it suffices to prove (5-21) for $\alpha \in\left(-\frac{1}{2}, \infty\right)^{d}$. We are going to apply Theorem 1. In order to do so we need to check that its assumptions are satisfied.

The formula (5-18) gives (A1) and (A2) with $K=1$. Conditions (T1) and (T2) follow from (5-19) and (5-20). It remains to prove (T3). For the space $\mathcal{D}$ this condition follows from [Nowak 2003, Lemma 4.3]. In the case of $\mathcal{D}_{i}, i=1, \ldots, d$, the assumption (T3) can be deduced from (T3) for $\mathcal{D}$ together with (5-20).

Thus, an application of Theorem 1 is justified and the proof of Theorem 11 is completed.
5.7. Jacobi operator: Jacobi function expansions. For parameters $\alpha, \beta \in(-1, \infty)^{d}$ we consider

$$
\begin{gathered}
p_{i}=1, \quad q_{i}=-\frac{1}{4}\left(2 \alpha_{i}+1\right) \cot \left(\frac{1}{2} x_{i}\right)+\frac{1}{4}\left(2 \beta_{i}+1\right) \tan \left(\frac{1}{2} x_{i}\right), \quad a_{i}=\frac{1}{4}\left(\alpha_{i}+\beta_{1}+1\right)^{2}, \\
w_{i}\left(x_{i}\right)=1, \quad d \mu_{i}\left(x_{i}\right)=d x_{i}, \quad X_{i}=(0, \pi),
\end{gathered}
$$

so that

$$
\begin{gather*}
\delta_{i}=\partial_{x_{i}}-\frac{1}{4}\left(2 \alpha_{i}+1\right) \cot \left(\frac{1}{2} x_{i}\right)+\frac{1}{4}\left(2 \beta_{i}+1\right) \tan \left(\frac{1}{2} x_{i}\right), \quad \mathfrak{d}_{i}=\partial_{x_{i}}, \\
\delta_{i}^{*}=-\partial_{x_{i}}-\frac{1}{4}\left(2 \alpha_{i}+1\right) \cot \left(\frac{1}{2} x_{i}\right)+\frac{1}{4}\left(2 \beta_{i}+1\right) \tan \left(\frac{1}{2} x_{i}\right), \quad v_{i}=\left[\delta_{i}, \delta_{i}^{*}\right]=\frac{2 \alpha_{i}+1}{8 \cos ^{2}\left(\frac{1}{2} x_{i}\right)}+\frac{2 \beta_{i}+1}{8 \sin ^{2}\left(\frac{1}{2} x_{i}\right)},  \tag{5-22}\\
r(x)=\sum_{i=1}^{d} \frac{1}{16}\left(2 \alpha_{i}+1\right)^{2} \cot ^{2}\left(\frac{1}{2} x_{i}\right)+\frac{1}{16}\left(2 \beta_{i}+1\right)^{2} \tan ^{2}\left(\frac{1}{2} x_{i}\right)+\frac{1}{16}\left(\left(\alpha_{i}+\beta_{i}+1\right)^{2}-\left(2 \alpha_{i}+1\right)\left(2 \beta_{i}+1\right)\right) .
\end{gather*}
$$

Here $L$ is the Jacobi operator

$$
L=\sum_{i=1}^{d} L_{i}=-\Delta+\sum_{i=1}^{d}\left(\frac{4 \alpha_{i}^{2}-1}{16 \sin ^{2}\left(\frac{1}{2} x_{i}\right)}+\frac{4 \beta_{i}^{2}-1}{16 \cos ^{2}\left(\frac{1}{2} x_{i}\right)}\right) .
$$

Then $L$ is symmetric on $C_{c}^{\infty}\left((0, \pi)^{d}\right)$ and has a self-adjoint extension given by

$$
L f=\sum_{k \in \mathbb{N}^{d}} \lambda_{k}\left\langle f, \phi_{k}^{\alpha, \beta}\right\rangle_{L^{2}} \phi_{k}^{\alpha, \beta} .
$$

Here $\lambda_{k}=\sum_{i=1}^{d} \lambda_{k_{i}}^{i}$ with $\lambda_{k_{i}}^{i}=\left(k_{i}+\frac{1}{2}\left(\alpha_{i}+\beta_{i}+1\right)\right)^{2}, L^{2}=L^{2}\left((0, \pi)^{d}, d x\right)$, while $\left\{\phi_{k}^{\alpha, \beta}\right\}_{k \in \mathbb{N}^{d}}$ is the system of $L^{2}$ normalized Jacobi functions; see [Nowak and Stempak 2006, Section 7.7]. These Jacobi functions have the tensor product form $\phi_{k}^{\alpha, \beta}=\phi_{k_{1}}^{\alpha_{1}, \beta_{1}} \otimes \cdots \otimes \phi_{k_{d}}^{\alpha_{d}, \beta_{d}}$ with

$$
\begin{equation*}
\phi_{k_{i}}^{\alpha_{i}, \beta_{i}}\left(x_{i}\right)=2^{\left(\alpha_{i}+\beta_{i}+1\right) / 2} \widetilde{P}_{k_{i}}^{\alpha_{i}, \beta_{i}}\left(\cos x_{i}\right)\left(\sin \left(\frac{1}{2} x_{i}\right)\right)^{\alpha_{i}+1 / 2}\left(\cos \left(\frac{1}{2} x_{i}\right)\right)^{\beta_{i}+1 / 2} \tag{5-23}
\end{equation*}
$$

for $x_{i} \in(0, \pi)$, and $\widetilde{P}_{k_{i}}^{\alpha_{i}, \beta_{i}}$ being the Jacobi polynomials from Section 5.3. In this section we take

$$
\varphi_{k}=\phi_{k}^{\alpha, \beta} .
$$

In the case when $\alpha, \beta \in\left[\frac{1}{2}, \infty\right)^{d}$, the $L^{2}$ kernel of $L$ is trivial, and thus the projection $\Pi$ equals the identity operator.

Next

$$
\begin{equation*}
\delta_{i} \phi_{k}^{\alpha, \beta}=-\sqrt{k_{i}\left(k_{i}+\alpha_{i}+\beta_{i}+1\right)} \phi_{k-e_{i}}^{\alpha+e_{i}, \beta+e_{i}}, \tag{5-24}
\end{equation*}
$$

and thus the Riesz transform is

$$
R_{i} f=-\sum_{k \in \mathbb{N}^{d}, k_{i} \neq 0}\left(\frac{k_{i}\left(k_{i}+\alpha_{i}+\beta_{i}+1\right)}{\lambda_{k}}\right)^{1 / 2}\left\langle f, \phi_{k}^{\alpha, \beta}\right\rangle_{L^{2}} \phi_{k-e_{i}}^{\alpha+e_{i}, \beta+e_{i}}, \quad f \in L^{2} .
$$

In the case $d=1$ the $L^{p}$ boundedness of these Riesz transforms was proved by Stempak [2007]. Using Theorem 1 we obtain the following multidimensional bounds.
Theorem 12. Let $\alpha, \beta \in\left[\frac{1}{2}, \infty\right)^{d}$. Then, for $1<p<\infty$, we have

$$
\|\boldsymbol{R} f\|_{p} \leqslant 48\left(p^{*}-1\right)\|f\|_{p}, \quad f \in L^{p} .
$$

Proof. A continuity argument based on (5-23) and (5-24) allows us to focus on $\alpha, \beta \in\left(\frac{1}{2}, \infty\right)^{d}$. We are going to apply Theorem 1 for such parameters $\alpha$ and $\beta$. In order to do so we need to check that its assumptions are satisfied.

The formula (5-22) gives (A1) and (A2) (with $K=1$ ). Conditions (T1) and (T2) follow from (5-23) and (5-24), while (T3) can be deduced from the density of polynomials in $C((-1,1))$ together with (5-23) and (5-24).

Thus, an application of Theorem 1 is permitted and the proof of Theorem 12 is completed.

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[^0]:    MSC2010: 42C10, 42A50, 33C50.
    Keywords: Riesz transform, Bellman function, orthogonal expansion.

[^1]:    ${ }^{\dagger}$ The hypoellipticity of $L_{i}$ is not necessary for the theory from [Nowak and Stempak 2006] to work (Nowak, personal communication, 2017). When not having this property, one has to add instead some extra assumptions (much weaker than smoothness) on the regularity of the eigenfunctions $\varphi_{k}$.

[^2]:    ${ }^{\dagger}$ In [Nowak 2004, Theorem 13] the author also states an estimate on $L^{p}$ for the vector of Riesz-Laguerre transforms that is dimension-free for certain values of $\alpha$. Unfortunately this result is not properly proved there (Nowak, personal communication, 2017). This is due to a problem in the proof of the vectorial $g$-function bound from [Nowak 2004, Theorem 7(b)].

