# Martingales, Singular Integrals, and Fourier Multipliers 

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## PURDUE UNIVERSITY <br> GRADUATE SCHOOL Thesis/Dissertation Acceptance

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MARTINGALES, SINGULAR INTEGRALS, AND FOURIER MULTIPLIERS

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Is approved by the final examining committee:

Dr. Rodrigo Banuelos
Chair
Dr. Fabrice Baudoin
Dr. Jonathon Peterson

Dr. Aaron Nung Kwan Yip

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Approved by Major Professor(s): Rodrigo Banuelos
Approved by: Dr. David Goldberg $\quad 6 / 21 / 2015$

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To my mother Elizabeth Ann Perlmutter and to my late father Franklin Lewis Perlmutter.

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#### Abstract

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Many probabilistic constructions have been created to study the $L^{p}$-boundedness, $1<p<\infty$, of singular integrals and Fourier multipliers. We will use a combination of analytic and probabilistic methods to study analytic properties of these constructions and obtain results which cannot be obtained using probability alone.

In particular, we will show that a large class of operators, including many that are obtained as the projection of martingale transforms with respect to the background radiation process of Gundy and Varapolous or with respect to space-time Brownian motion, satisfy the assumptions of Calderón-Zygmund theory and therefore boundedly map $L^{1}$ to weak- $L^{1}$.

We will also use a method of rotations to study the $L^{p}$ boundedness, $1<p<\infty$, of Fourier multipliers which are obtained as the projections of martingale transforms with respect to symmetric $\alpha$-stable processes, $0<\alpha<2$. Our proof does not use the fact that $0<\alpha<2$ and therefore allows us to obtain a larger class of multipliers, indexed by a parameter, $0<r<\infty$, which are bounded on $L^{p}$. As in the case of the multipliers which arise as the projection of martingale transforms, these new multipliers also have potential applications to the study of the Beurling-Ahlfors transform and are related to the celebrated conjecture of $T$. Iwaniec concerning its exact $L^{p}$ norm.


## 1. Introduction

### 1.1 Overview

Martingale inequality methods provide a powerful tool to study the $L^{p}$ boundedness, $1<p<\infty$, of the basic Calderón-Zygmund singular integral operators and other Fourier multipliers on $\mathbb{R}^{n}$. An advantage of these techniques is that they give very good information on the size of these $L^{p}$ bounds and, in particular, provide constants that are independent of the dimension. These same arguments can be used to extend results from $\mathbb{R}^{n}$ to manifolds and to the Ornstein-Uhleneck case. For some applications of these methods we refer the reader to [3], [7], [11], [8], [6], [16], [25], [35], [24], and the many references provided there. However, as powerful as these techniques are, weak-type martingale inequalities cannot be directly transferred to singular integral operators. For example, while Burkholder's celebrated $L^{p}$ inequalities, $1<p<\infty$, for martingale transforms [18], with his famous bound " $\left(p^{*}-1\right)$ ", gives the same $L^{p}$ bound for many singular integral operators, his weak-type martingale transform bound " 2 " provides no information for the weak-type inequalities of those operators. This is due to the fact that the probabilistic representation of such operators involves the use of conditional expectation which does not preserve weak-type inequalities. When viewed as analytic objects, many of the operators which are obtained as the projections of martingale transforms have natural generalizations which cannot be studied by purely probabilistic methods. Therefore, in chapters 2 and 3 we will use a combination of analytic and probabilistic techniques to study these operators.

The main result of chapter 2 is that a very general class of operators, including many of the operators considered in [11] and [8], which arise as the projections of martingale transforms, are in fact Calderón-Zygmund operators. This class includes operators which are not, in general, of convolution type. Once we know that these are

Calderón-Zygmund operators, they then satisfy all the properties of such operators, including their weak-type boundedness. This does not, of course, answer an important question that has been of interest to many people for many years, originally raised by Stein in [40] in the case of the Riesz transforms: do these operators have weak-type bounds independent of the dimension? An affirmative answer would in turn raise an even more important question: do weak-type inequalities hold for Riesz transforms on Wiener space? After nearly 35 years and the efforts of many, these questions remain completely open. For a more precise formulation of these questions, see [4].

The purpose of chapter 3 is to study the $L^{p}$ boundedness of a class of Fourier multipliers which are closely related to multipliers obtained as the projection of martingale transform of $\alpha$-stable processes, $0<\alpha<2$, in [9] and [5]. Using analytic methods, we are able to obtain a family of operators indexed by $r, 0<r<\infty$, that are bounded on $L^{p}\left(\mathbb{R}^{n}\right)$. When $0<r<2$, these operators coincide with the operators from [9] and [5]. However, for $r \geq 2$ these are a new family of operators whose $L^{p}$ boundedness have not been previously studied (except in the trivial case that $p=2$ ). These problems are motivated by a celebrated 1982 conjecture of Tadeusz Iwaniec [29] which asserts that the $L^{p}$ norm of the Beurling-Ahlfors operator is the same as the $L^{p}$ norm of martingale transforms, namely $\left(p^{*}-1\right)$. Although great progress has been made on this conjecture, it too remains open.

### 1.2 Calderón-Zygmund Operators

Following standard terminology (see for example [26, p.175]), we will say that an operator $T$ acting on the Schwartz space of rapidly decreasing functions on $\mathbb{R}^{n}$ is a Calderón-Zygmund (CZ) operator if it admits a bounded extension to $L^{2}\left(\mathbb{R}^{n}\right)$ and is of the form

$$
\begin{equation*}
T f(x)=\lim _{\epsilon \searrow 0} \int_{|x-\tilde{x}|>\epsilon} K(x, \tilde{x}) f(\tilde{x}) d \tilde{x} \tag{1.1}
\end{equation*}
$$

where the kernel $K$ is defined on the set $\{x \neq \tilde{x}\}$ and satisfies the following conditions

$$
\begin{align*}
|K(x, \tilde{x})| & \leq \frac{\kappa}{|x-\tilde{x}|^{n}}  \tag{1.2}\\
\left|\nabla_{x} K(x, \tilde{x})\right| & \leq \frac{\kappa}{|x-\tilde{x}|^{n+1}}  \tag{1.3}\\
\left|\nabla_{\tilde{x}} K(x, \tilde{x})\right| & \leq \frac{\kappa}{|x-\tilde{x}|^{n+1}}, \tag{1.4}
\end{align*}
$$

for some universal constant $\kappa$. Integrals as in (1.1) are referred to as principal value integrals. For the rest of this thesis, we will assume that all integrals, where the integrand has an isolated singularity, are to be interpreted as principal value integrals. If there exists a function $\bar{K}$, defined on $\mathbb{R}^{n} \backslash\{0\}$, so that $\bar{K}(x-\tilde{x})=K(x, \tilde{x})$ for all $x \neq \tilde{x}$, then we say that $T$ is of convolution type. The Hilbert, Riesz, and BeurlingAhlfors transforms discussed below are basic examples of CZ operators of convolution type which give rise to interesting Fourier multipliers. It is well known (see for example [26, p.183]) that CZ operators are strong-type ( $p, p$ ) for $1<p<\infty$ and are weak-type ( 1,1 ). More precisely, there exists universal constants $C_{p, n, \kappa}$, depending only on $p, n$, and $\kappa$, such that

$$
\begin{equation*}
\|T f\|_{p} \leq C_{p, n, \kappa}\|f\|_{p}, \quad 1<p<\infty \tag{1.5}
\end{equation*}
$$

and

$$
\begin{equation*}
|\{x:|T f(x)|>\lambda\}| \leq \frac{C_{1, n, \kappa}}{\lambda}\|f\|_{1} \tag{1.6}
\end{equation*}
$$

where here and below we use $|A|$ to denote the Lebesgue measure of a set $A$.
We note that CZ operators do not, in general, map $L^{1}\left(\mathbb{R}^{n}\right)$ into $L^{1}\left(\mathbb{R}^{n}\right)$, nor do they map $L^{\infty}\left(\mathbb{R}^{n}\right)$ into $L^{\infty}\left(\mathbb{R}^{n}\right)$. They do, however, map the Hardy space $H^{1}\left(\mathbb{R}^{n}\right)$, an important subset of $L^{1}\left(\mathbb{R}^{n}\right)$, into $L^{1}\left(\mathbb{R}^{n}\right)$ and map $L^{\infty}\left(\mathbb{R}^{n}\right)$ into the set of functions with bounded mean oscillation. This topic will be discussed further in section 2.2.

### 1.2.1 The Hilbert Transform

The Hilbert transform is the prototypical example of a Calderón-Zygmund operator of convolution type. For a rapidly decreasing function, $f: \mathbb{R} \rightarrow \mathbb{R}$, we define

$$
H f(x)=\frac{1}{\pi} \int_{\mathbb{R}} \frac{f(y)}{x-y} d y
$$

In other words, the Hilbert transform is the operator given by convolving a function with the singular kernel $\frac{1}{\pi x}$. It is important to note that this integral must be interpreted as a principal value integral since otherwise it may not converge. (Take, for instance, $x=0$ and $f(y)=e^{-y^{2}}$.) An interesting property of the Hilbert transform is that if we let $u(x, y)$ and $v(x, y)$ be the extensions of $f(x)$ and $H f(x)$ to the upper half-space by convolution against the Poisson kernel $p_{y}(x)$, then $u(x, y)+i v(x, y)$ is holomorphic on the upper half-space.

### 1.2.2 The Riesz Transforms

The natural generalizations of the Hilbert transform to higher dimensions are known as the Riesz transforms. For $f \in L^{p}\left(\mathbb{R}^{n}\right)$, we define the Riesz transform in direction $j, 1 \leq j \leq n$, by

$$
R_{j} f(x)=\frac{\Gamma\left(\frac{n+1}{2}\right)}{\pi^{(n+1) / 2}} \int_{\mathbb{R}^{n}} \frac{x_{j}-\tilde{x}_{j}}{|x-\tilde{x}|^{n+1}} f(\tilde{x}) d \tilde{x} .
$$

When $n=1$, this reduces to the Hilbert transform. $R_{j}$ is a Fourier multiplier with

$$
\widehat{R_{j} f}(\xi)=\frac{i \xi_{j}}{|\xi|} \widehat{f}(\xi) .
$$

In [40], Stein showed that we may take the constant, $C_{p, n}$, to be independent of $n$ in (1.5) for the Riesz transforms. Whether or not the constant in (1.6) can be taken independent of n is unknown with the best known result being that the constant is at worst $O(\sqrt{n})$ as $n \rightarrow \infty$ [31]. Gundy and Varopoulos showed in [27] that the Riesz transforms could be interpreted probabilistically as projections of martingale transforms, and from this it again follows that the constant may be taken to be
independent of dimension. See [4] for more on this topic. These techniques were further explored by Bañuelos and Wang in [8] to prove the sharp inequalities

$$
\left\|R_{j} f\right\|_{p} \leq C_{p}\|f\|_{p} \quad \text { and } \quad\left\|\left(\left(R_{j} f\right)^{2}+f^{2}\right)^{1 / 2}\right\|_{p} \leq \sqrt{C_{p}^{2}+1}\|f\|_{p}
$$

where

$$
p^{*}=\max \left\{p, \frac{p}{p-1}\right\}, \quad \text { and } \quad C_{p}=\cot \left(\frac{\pi}{2 p^{*}}\right) .
$$

The first inequality had been proved earlier by Iwaniec and Martin in [30] using the method of rotations.

### 1.2.3 The Beurling-Ahlfors Transfrom

For $f \in L^{p}(\mathbb{C})$, we define the Beurling-Ahlfors operator by

$$
B f(z)=-\frac{1}{\pi} \int_{\mathbb{C}} \frac{f(w)}{(z-w)^{2}} d w
$$

$B$ is a Fourier multiplier with

$$
\widehat{B f}(\xi)=\frac{\xi_{1}^{2}-\xi_{2}^{2}-2 i \xi_{1} \xi_{2}}{|\xi|^{2}} \widehat{f}(\xi)
$$

Therefore, we can decompose the Beurling-Ahlfors transform into a linear combination of second order Riesz transforms,

$$
\begin{equation*}
B=R_{2}^{2}-R_{1}^{2}+2 i R_{1} R_{2} . \tag{1.7}
\end{equation*}
$$

Because of its many connections to quasiconformal mappings and other topics in complex analysis (see for example [2]) there has been a lot of interest for many years in finding its operator norm on $L^{p}(\mathbb{C}), 1<p<\infty$, which we denote $\|B\|_{p}$. In [33], Lehto showed that $\|B\|_{p} \geq\left(p^{*}-1\right)$. A long standing conjecture of Iwaniec [29] is that $\|B\|_{p}=\left(p^{*}-1\right)$. Despite the efforts of many researchers, Iwaniec's conjecture remains open. There are, however, many partial results, and the techniques developed in these efforts have lead to many other interesting questions and applications. In particular, there are a number of probabilistic constructions which provide upper bounds for $\|B\|_{p}$.

In [8], Bañuelos and Wang showed that $\|B\|_{p} \leq 4\left(p^{*}-1\right)$. This constant was reduced to $2\left(p^{*}-1\right)$ by Nazarov and Volberg in [35] using a Littlewood-Paley inequality proved using Bellman functions techniques. The Bellman function in [35] is itself constructed from Burkholder martingale inequalities. In [11] the martingale techniques from [8] were applied to space-time Brownian motion to reproduce the bound $2\left(p^{*}-1\right)$. The methods of [11] were refined in [10] to reduce this constant to $1.575\left(p^{*}-1\right)$, which is the best known bound as of now valid for all $1<p<\infty$. We do point out that for $1000<p<\infty$, this bound was improved to $1.4\left(p^{*}-1\right)$ in [16].

### 1.3 Multiplier Theorems

Two important tools for studying the $L^{p}\left(\mathbb{R}^{n}\right)$ boundedness of Fourier multipliers, which we will use in chapter 3, are the Marcinkiewicz mutliplier theorem and the Hörmander-Mikhlin multiplier theorem which we state below for convenience. For proofs of these results see [26] or [39].

Theorem 1.3.1 (Marcinkiewicz). Let $m \in L^{\infty}\left(\mathbb{R}^{n}\right)$ with $\|m\|_{\infty} \leq K$ for some $0<$ $K<\infty$. Supposed that $m(\xi)$ is $n$-times continuously differentiable on the subset of $\mathbb{R}^{n}$ where none of the $\xi_{i}$ are zero. For $j \in \mathbb{Z}$, let $I_{j}$ denote the dyadic interval $\left(-2^{j+1},-2^{j}\right] \cup\left[2^{j}, 2^{j+1}\right)$. Suppose that for all $1 \leq k \leq n$, for all subsets $\left\{i_{1}, \ldots, i_{k}\right\}$ of $\{1, \ldots, n\}$ of order $k$, and for all integers $l_{i_{1}}, \ldots l_{i_{k}}$, we have that

$$
\begin{equation*}
\int_{I_{l_{i_{1}}}} \ldots \int_{I_{l_{i_{k}}}}\left|\partial_{i_{1}} \ldots \partial_{i_{k}} m(\xi)\right| d \xi_{i_{k}} \ldots d \xi_{i_{1}} \leq K<\infty \tag{1.8}
\end{equation*}
$$

whenever $\xi_{j} \neq 0$ for all $j \notin\left\{i_{1}, \ldots, i_{k}\right\}$. Then $m(\xi)$ is a bounded Fourier multiplier on $L^{p}\left(\mathbb{R}^{n}\right)$ for all $1<p<\infty$ and

$$
\left\|T_{m} f\right\|_{p} \leq C_{n} K\left(p^{*}-1\right)^{6 n}\|f\|_{p} \quad \text { for all } f \in L^{p}\left(\mathbb{R}^{n}\right)
$$

where $C_{n}$ is a constant depending only on $n$.

Theorem 1.3.2 (Hörmander-Mikhlin). Let $n_{0}=\left\lfloor\frac{n}{2}\right\rfloor+1$, and let $m(\xi)$ be $n_{0}$-times differentiable on $\mathbb{R}^{n} \backslash\{0\}$. Suppose there exists $0<K<\infty$ such that $\|m\|_{\infty} \leq K$ and that also

$$
\begin{equation*}
\sup _{R>0} R^{-n+2|\beta|} \int_{R<|\xi|<2 R}\left|\partial^{\beta} m(\xi)\right|^{2} d \xi<K^{2} \tag{1.9}
\end{equation*}
$$

for all multi-indices such that $|\beta| \leq n_{0}$. Then $m(\xi)$ is a bounded Fourier multiplier on $L^{p}\left(\mathbb{R}^{n}\right)$ for all $1<p<\infty$ and there exists $C_{n}$ depending only on $n$ such that

$$
\left\|T_{m} f\right\|_{p} \leq C_{n} K\left(p^{*}-1\right)\|f\|_{p} .
$$

### 1.4 Lévy Processes

A Lévy process on $\mathbb{R}^{n}$ is an $\mathbb{R}^{n}$-valued stochastic process, $\left(X_{t}\right)_{t \geq 0}$, which almost surely starts at the origin, has stationary, independent increments, and satisfies the stochastic continuity condition $\lim _{t \searrow 0} \mathbb{P}\left(|X|_{t}>\epsilon\right)=0$ for all $\epsilon>0$. The famous LévyKhintchine formula states that there exists a point $b \in \mathbb{R}^{n}$, a non-negative symmetric $n \times n$ matrix $B$, and a measure $\nu$ such that $\nu(\{0\})=0$ and

$$
\int_{\mathbb{R}^{n}} \min \left\{|z|^{2}, 1\right\} d \nu(z)<\infty,
$$

such that the characteristic function of $X_{t}$ is given by $\mathbb{E}\left(e^{i \xi \cdot X_{t}}\right)=e^{t \rho(\xi)}$ where

$$
\rho(\xi)=i b \cdot \xi-\frac{1}{2} B \xi \cdot \xi+\int_{\mathbb{R}^{n}}\left[e^{i \xi \cdot z}-1-i(\xi \cdot z) \mathbb{I}_{(|z|<1)}\right] \nu(d z) .
$$

$(b, B, \nu)$ is referred to as the Lévy triple of $X_{t}$. The triple $(b, 0,0)$ corresponds to a drift process $X_{t}=b t ;(0, B, 0)$ corresponds to a centered Gaussian process with whose covariance is given by $\left[X_{s}^{i}, X_{t}^{j}\right]=B_{i, j} \min \{s, t\}$; and $(0,0, \nu)$ corresponds to a "purejump" process. If $X_{t}$ and $Y_{t}$ are independent Lévy processes with triples ( $b_{X}, B_{X}, \nu_{X}$ ) and $\left(b_{Y}, B_{Y}, \nu_{Y}\right)$, then $X_{t}+Y_{t}$ is a Lévy process with the triple $\left(b_{X}+b_{Y}, B_{X}+\right.$ $\left.B_{Y}, \nu_{X}+\nu_{Y}\right)$. Therefore, the Lévy-Khinchtine formula says that any Lévy process can be decomposed into the sum of three independent Lévy processes, a drift process, a centered Gaussian process, and a pure-jump process. For further background on Lévy processes see [13], [14], and [38].

### 1.4.1 $\alpha$-stable Processes

For $0<\alpha \leq 2$, the symmetric $\alpha$-stable process is a Lévy processes, $\left(X_{t}\right)_{t \geq 0}$ with $\rho(\xi)=-|\xi|^{\alpha}$. In the case that $\alpha=2,(b, B, \nu)=(0, I, 0),\left(X_{t}\right)_{t \geq 0}$ is Brownian motion (running at twice the usual speed), and density of $X_{t}$ is given by the Gaussian heat kernel

$$
\begin{equation*}
\frac{1}{(4 \pi t)^{n / 2}} e^{-|x|^{2} / 4 t} \tag{1.10}
\end{equation*}
$$

For $0<\alpha \leq 2$, we have that $(b, B, \nu)=\left(0,0, d \nu(z)=C_{n, \alpha} \frac{1}{z z n^{n+\alpha}} d z\right)$. If $\alpha=1$, then $\left(X_{t}\right)_{t \geq 0}$ is the Cauchy process and the density of $X_{t}$ is given by the Poisson kernel

$$
\begin{equation*}
\frac{\Gamma\left(\frac{n+1}{2}\right)}{\pi^{(n+1) / 2}} \frac{t}{\left(|x|^{2}+t^{2}\right)^{(n+1) / 2}} \tag{1.11}
\end{equation*}
$$

Except for in the cases $\alpha=1$ and $\alpha=2$, we do not have a simple analytic expression for the density of $X_{t}$ as in (1.10) and (1.11). However, there are a number of integral representations which are available for any $\alpha$.

### 1.5 Martingale Transforms

The study of martingale transforms and their boundedness on $L^{p}$ dates back to D . L. Burkholder's 1966 paper [17]. Since that time, martingale transforms have been extensively studied for both their theoretical importance in probability theory and their applications to finance. As alluded to in subsection 1.1, they have also been widely used to study the boundedness of singular integrals and Fourier multipliers.

### 1.5.1 Discrete Martingales

If $\left(f_{n}\right)_{n \geq 0}$ is a discrete-time martingale defined on a probability space $\left(\Omega, \mathbb{F}_{\infty}, \mathbb{P}\right)$, with filtration $\mathbb{F}=\left(\mathbb{F}_{n}\right)_{n \geq 0}$, then we may define a difference sequence,

$$
d_{k}=f_{k}-f_{k-1}
$$

for $k \geq 1$, and $d_{0}=f_{0}$ so that

$$
f_{n}=\sum_{k=0}^{n} d_{k} .
$$

If $v_{k}$ is a predictable sequence of random variables, in the sense that $v_{0}=0$ and $v_{k} \in \mathbb{F}_{k-1}$ for $k \geq 1$, such that $\left|v_{k}\right| \leq 1$ a.s. for all $k$, then the martingale transform of $f$ by $v$ is defined by

$$
(v * f)_{n}=\sum_{k=0}^{n} v_{k} d_{k}
$$

It is straight forward to check that $(v * f)_{n}$ is a martingale. The primary result of [17] was that the mapping $f \rightarrow v * f$ is bounded on $L^{p}, 1<p<\infty$, and weak-type $(1,1)$. That is, there exists constants $C_{p}$ and $C_{1}$ such that

$$
\begin{equation*}
\left.\left(\sup _{n} \mathbb{E}\left|(v * f)_{n}\right|^{p}\right)^{1 / p} \leq C_{p}\left(\sup _{n} \mathbb{E}\left|f_{n}\right|^{p}\right)\right)^{1 / p} \tag{1.12}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathbb{P}\left\{\sup _{n}\left|(v * f)_{n}\right|>\lambda\right\} \leq \frac{C_{1}}{\lambda}\|f\|_{1}, \quad \text { for all } \lambda>0 \tag{1.13}
\end{equation*}
$$

Eighteen years later in [18], Burkholder was able to show that, for all $1<p<\infty$, the best possible value of $C_{p}$ in (1.12) is $p^{*}-1$ and that the best possible constant in (1.13) is 2.

In [19], Burkholder introduced a condition called differential subordination, which allows for a much simpler proof of the fact that (1.12) holds with constant $p^{*}-1$. Furthermore, this construction allows us to extend this result to martingale transforms defined with respect to stochastic integrals. Let $\left(f_{n}\right)_{n \geq 0}$ and $\left(g_{n}\right)_{n \geq 0}$ be martingales taking values in a separable Hilbert space and let $d_{k}$ and $e_{k}$ be their difference sequences (so that $f_{n}=\sum_{k=0}^{n} d_{k}$ and $g_{n}=\sum_{k=0}^{n} e_{k}$ ). We say that $g_{n}$ is differentially subordinate to $f_{n}$ if $\left|e_{k}\right| \leq\left|d_{k}\right|$ a.s. for all $k$. Burkholder showed that if $g_{n}$ is differentially subordinate to $f_{n}$, then

$$
\begin{equation*}
\left.\left(\sup _{n} \mathbb{E}\left|g_{n}\right|^{p}\right)^{1 / p} \leq\left(p^{*}-1\right)\left(\sup _{n} \mathbb{E}\left|f_{n}\right|^{p}\right)\right)^{1 / p} \tag{1.14}
\end{equation*}
$$

Note that $p^{*}-1$ is the same as the constant that appears in Iwaniec's conjecture regarding the Beurling-Ahlfors transform.

### 1.5.2 Stochastic Integrals and Continuous-Time Martingales

Let $\left(X_{t}\right)_{t \geq 0}$ and $\left(Y_{t}\right)_{t \geq 0}$ be martingales on a probability space $\left(\Omega, \mathbb{F}_{\infty}, \mathbb{P}\right)$, with a common filtration, $\mathbb{F}=\left(\mathbb{F}_{t}\right)_{t \geq 0}$ that take values in a separable Hilbert space. Assume that $\mathbb{F}$ is right continuous and $\mathbb{F}_{0}$ contains all events of probability zero. Let $[X]_{t}$ and $[Y]_{t}$ denote quadratic variations of $X_{t}$ and $Y_{t}$ respectively. $Y_{t}$ is said to be differentially subordinate to $X_{t}$ if $\left|Y_{0}\right| \leq\left|X_{0}\right|$ and the process $[X]_{t}-[Y]_{t}$ is non-decreasing. Note that the quadratic variation of a discrete-time martingale $f_{n}=\sum_{k=0}^{n} d_{k}$ is given by $[f]_{n}=\sum_{k=0}^{n}\left|d_{k}\right|^{2}$. Therefore, this condition is the natural generalization of the differential subordination condition for discrete martingales. In [8], Bañuelos and Wang, showed that if $X_{t}$ and $Y_{t}$ have continuous sample paths, and $Y_{t}$ is differentially subordinate to $X_{t}$, then

$$
\begin{equation*}
\left(\sup _{t} \mathbb{E}\left|Y_{t}\right|^{p}\right)^{1 / p} \leq\left(p^{*}-1\right)\left(\sup _{t} \mathbb{E}\left|X_{t}\right|^{p}\right)^{1 / p} \tag{1.15}
\end{equation*}
$$

In the case that $X_{t}$ and $Y_{t}$ are orthogonal, in the sense that their quadratic covariation is zero, the constant $p^{*}-1$ can be improved to $\cot \left(\frac{\pi}{2 p^{*}}\right)$. We note that in [41], Wang showed (1.15) holds even if the assumption of continuous sample paths is removed. This fact will be important when we consider martingale transforms with respect to general Lévy processes.

A particularly important class of examples are martingales of the form

$$
X_{t}=\int_{0}^{t} H_{s} \cdot d B_{s}
$$

where $B_{t}$ is $n$-dimensional Brownian motion and $H_{s}$ is a $\mathbb{R}^{n}$-valued predictable process. If $A_{s}$ is a predictable, matrix-valued process such that for all $s>0$, and all $v \in \mathbb{R}^{n},\left|A_{s} v\right| \leq|v|$, then

$$
(A * X)_{t}=\int_{0}^{t} A_{s} H_{s} \cdot d B_{s}
$$

is called the martingale transform of $X$ by $A$. Similarly to the discrete martingale transforms in the previous subsection, $(A * X)_{t}$ is a martingale that is differentially subordinate to $X_{t}$. We remark that if for all $s>0$ and all $v \in \mathbb{R}^{n}, A v \cdot v=0$, then $X_{t}$ and $(A * X)_{t}$ are orthogonal.

### 1.6 Martingales Transforms and Harmonic Analysis

As alluded to in section 1.1, there are several constructions which use martingale transforms to study the $L^{p}$ boundedness of the classical Calderón-Zygmund singular integrals mentioned in section 1.2 and other Fourier multipliers. In many of these constructions, the method is based on the same fundamental idea. For a function $f$ in $L^{p}\left(\mathbb{R}^{n}\right)$, we construct a martingale $M(f)_{t}$ such that $\sup _{t}\left\|M(f)_{t}\right\|_{p}=\|f\|_{p}$. Then we apply a martingale transform to get a new martingale, $N(f)_{t}$, such that $\sup _{t}\left\|N(f)_{t}\right\|_{p} \leq C_{p} \sup _{t}\left\|M(f)_{t}\right\|_{p}$. Finally, we project $N(f)_{t}$ onto $L^{p}\left(\mathbb{R}^{n}\right)$ using conditional expectation to get a new function which we denote by $S f(x)$. Conditional expectation is a contraction on $L^{p}\left(\mathbb{R}^{n}\right)$ so $\|S f\|_{p} \leq \sup _{t}\left\|N(f)_{t}\right\|_{p}$. Combining these three inequalities yields $\|S f\|_{p} \leq C_{p}\|f\|_{p}$. If appropriate choices are made at each step, this operator will coincide with an operator of classical interest in analysis such as the Riesz Transforms or Beurling-Ahlfors transform.

### 1.6.1 The Background Radiation Process

We first consider the construction developed by Gundy and Varopoulos in [27] and used by Bañuelos and Wang in [8]. The first step is to construct a martingale corresponding to each function $f \in L^{p}$. Let

$$
\begin{equation*}
p_{y}(x)=\frac{\Gamma\left(\frac{n+1}{2}\right)}{\pi^{(n+1) / 2}} \frac{y}{\left(|x|^{2}+y^{2}\right)^{(n+1) / 2}} \tag{1.16}
\end{equation*}
$$

be the Poisson kernel for the upper half-space, $\mathbb{R}_{+}^{n+1}$, and for $f(x) \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$, let $\left(p_{y} * f\right)(x)=u_{f}(x, y)$ be the Poisson extension of $f$. (Note that by (1.11) $p_{y}$ is the density of the Cauchy process at time $y$.) Background radiation is a "time-reversed Brownian motion," $\left(B_{t}\right)_{t \leq 0}$, taking values in $\mathbb{R}_{+}^{n+1}$ such that $B_{-\infty}$ has distribution given by the Lebesgue measure on $\mathbb{R}^{n} \times\{\infty\}$, and $B_{0}$ is distributed by the Lebesgue measure on $\mathbb{R}^{n} \times\{0\}$. We write $B_{t}=\left(X_{t}, Y_{t}\right)$ with $X_{t}$ taking values in $\mathbb{R}^{n}$ and $Y_{t}>0$. The standard rules of stochastic calculus, in particular Itô's formula, hold for the background radiation process. Therefore, $u_{f}\left(X_{t}, Y_{t}\right)$ is a martingale and

$$
f\left(X_{0}\right)=u_{f}\left(B_{0}\right)=\int_{-\infty}^{0} \nabla u_{f}\left(X_{s}, Y_{s}\right) \cdot d B_{s}
$$

where $\nabla=\left(\partial_{x_{1}}, \ldots, \partial_{x_{n}}, \partial_{y}\right)$. If $A(x, y)$ is an $(n+1) \times(n+1)$ matrix-valued function such that

$$
\|A\|=\left\|\sup _{|v| \leq 1}(|A(x, y) v|)\right\|_{L^{\infty}\left(\mathbb{R}^{n} \times[0, \infty)\right)}<\infty
$$

we define the martingale transform of $f$ by $A$ as

$$
\begin{equation*}
(A * f)=\int_{-\infty}^{0} A\left(X_{s}, Y_{s}\right) \nabla u_{f}\left(X_{s}, Y_{s}\right) \cdot d B_{s} . \tag{1.17}
\end{equation*}
$$

The random variable $A * f$ is not a function of the endpoint, $X_{0}$. This motivates us to define a projection operator by averaging the integral in (1.17) over all paths ending at $x$, that is,

$$
T_{A} f(x)=\mathbb{E}\left(\int_{-\infty}^{0} A\left(X_{s}, Y_{s}\right) \nabla u_{f}\left(X_{s}, Y_{s}\right) \cdot d B_{s} \mid X_{0}=x\right)
$$

It is known (see [4]) that $\mathbb{E}\left(\left.\left|\left(\left.f\left(B_{0}\right)\right|^{p}\right)=\int_{\mathbb{R}^{n}}\right| f(x)\right|^{p} d x\right.$, which implies

$$
\sup _{t \geq 0}\left\|u_{f}\left(B_{t}\right)\right\|_{p}=\|f\|_{p}
$$

since $\left|u_{f}\left(B_{t}\right)\right|^{p}$ is a submartingale. In other words, lifting $f \in L^{p}\left(\mathbb{R}^{n}\right)$ to the space of martingales does not change its norm. Combining this with the fact that conditional expectation is a contraction in $L^{p}\left(\mathbb{R}^{n}\right)$, we see that the operator norm of $T_{A}$ is the same as the operator norm of the martingale transform $X \rightarrow A * X$. Thus, we have

$$
\left\|T_{A} f(x)\right\|_{p} \leq\left(p^{*}-1\right)\|A\|\|f\|_{p}
$$

It is known (see for example [18]) that martingale transforms are weak-type $(1,1)$ and in fact we have the sharp inequality

$$
P\{|A * X|>\lambda\} \leq \frac{2\|A\|}{\lambda}\|X\|_{1} .
$$

Unfortunately, this does not give us information about the weak-type behavior of $T_{A}$ because weak-type inequalities are not preserved under conditional expectation. However, we can represent $T_{A}$ analytically by finding a kernel $K_{A}(x, \tilde{x})$ such that

$$
T_{A} f(x)=\int_{\mathbb{R}^{n}} K_{A}(x, \tilde{x}) f(\tilde{x}) d \tilde{x} .
$$

Let $f, g \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ and note that

$$
g\left(B_{0}\right)=\int_{-\infty}^{0} \nabla u_{g}\left(B_{s}\right) \cdot d B_{s}
$$

by Itô's formula. Therefore, using basic facts about the covariation of stochastic integrals and the occupation time formula for the background radiation process, (see [23, p. 31 and 57] and [27])

$$
\begin{align*}
\int_{\mathbb{R}^{n}} T_{A} f(x) g(x) d x & =\int_{\mathbb{R}^{n}} \mathbb{E}\left(\int_{-\infty}^{0} A\left(X_{s}, Y_{s}\right) \nabla u_{f}\left(X_{s}, Y_{s}\right) \cdot d B_{s} \mid X_{0}=x\right) g(x) d x \\
& =\mathbb{E}\left(\int_{-\infty}^{0} A\left(X_{s}, Y_{s}\right) \nabla u_{f}\left(X_{s}, Y_{s}\right) \cdot d B_{s} g\left(B_{0}\right)\right) \\
& =\mathbb{E}\left(\int_{-\infty}^{0} A\left(X_{s}, Y_{s}\right) \nabla u_{f}\left(X_{s}, Y_{s}\right) \cdot d B_{s} \int_{-\infty}^{0} \nabla u_{g}\left(B_{s}\right) \cdot d B_{s}\right) \\
& =\mathbb{E}\left(\int_{-\infty}^{0} A\left(X_{s}, Y_{s}\right) \nabla u_{f}\left(X_{s}, Y_{s}\right) \cdot \nabla u_{g}\left(B_{s}\right) d s\right) \\
& =\int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 y A(x, y) \nabla u_{f}(x, y) \cdot \nabla u_{g}(x, y) d x d y . \tag{1.18}
\end{align*}
$$

Using the fact that $\nabla u_{f}(x, y)=\left(\left(\nabla p_{y}\right) * f\right)(x)$ and applying Fubini's theorem, we see that we have

$$
K_{A}(x, \tilde{x})=\int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 y A(\bar{x}, y) \nabla p_{y}(\bar{x}-\tilde{x}) \cdot \nabla p_{y}(\bar{x}-x) d \bar{x} d y .
$$

This representation will be used in chapter 2 show that $T_{A}$ is, under mild assumptions,
a Calderón-Zygmund operator.
If we define $A_{j}=\left(a_{l, m}^{j}\right)$ by

$$
a_{l, m}^{j}=\left\{\begin{array}{lr}
1 & l=n+1, m=j \\
-1 & l=j, m=n+1 \\
0 & \text { otherwise }
\end{array}\right\}
$$

then plugging into (1.18) and Fourier transforming shows that $T_{A_{j}}=R_{j}$. Since $A_{j}$ satisfies the orthogonality condition $A_{j} v \cdot v=0$ for all $v \in \mathbb{R}^{n}$, it follows that

$$
\left\|R_{j}\right\|_{p} \leq \cot \left(\frac{\pi}{2 p^{*}}\right)\|f\|_{p}, \quad 1<p<\infty
$$

We can also define $A$ in such a way that $T_{A}=R_{i, j}$ and $\|A\|=1$. This implies,

$$
\left\|R_{i} R_{j}\right\|_{p} \leq\left(p^{*}-1\right)\|f\|_{p}, \quad 1<p<\infty .
$$

By (1.7), this also implies that $\|B\|_{p} \leq 4\left(p^{*}-1\right)$. If $A$ is any matrix with constant coefficients, $T_{A}$ will be a linear combination of the identity and first and second order Riesz transforms. Moreover, if $A(x, y)=A(y)$ is independent of $x$ and $\|A\|<\infty$, then $T_{A}$ is a Fourier multiplier. For more examples of multipliers corresponding to various choices of $A$, see [8] and [4].

### 1.6.2 Space-Time Brownian Motion

The approach of [11] is similar to the construction discussed in the previous subsection, but uses space-time Brownian motion and the heat kernel for the half Laplacian,

$$
\begin{equation*}
h_{t}(x)=\frac{1}{(2 \pi t)^{n / 2}} e^{-|x|^{2} / 2 t}, \tag{1.19}
\end{equation*}
$$

instead of background radiation and the Poisson kernel. (We remark that $h_{t}$ is the density of a standard Brownian motion at time $t$. Observe that this is, up to a simple time change, $t=2 s$, the density of the stable process given in (1.10).) Fix $T>0$, and let $Z_{t}=\left(B_{t}, T-t\right)$ for $0 \leq t \leq T$ where $B_{t}$ is Brownian motion on $\mathbb{R}^{n}$ with initial distribution given by the Lebesgue measure. Letting $u_{f}(x, t)$ denote the extension of $f$ to the upper half-space by convolution with $h_{t}$, Itô's formula shows that $u_{f}\left(Z_{t}\right)$ is a martingale and

$$
u_{f}\left(Z_{t}\right)=\int_{0}^{s} \nabla_{x} u_{f}\left(B_{s}, T-s\right) \cdot d B_{s} .
$$

For an $n \times n$ matrix-valued function, $A(x, t)$, such that $\|A\|<\infty$, we define a martingale transform and a projection operator by

$$
A * f=\int_{0}^{T} A\left(B_{s}, T-s\right) \nabla_{x} u_{f}\left(B_{s}, T-s\right) \cdot d B_{s}
$$

and

$$
S_{A}^{T} f(x)=\mathbb{E}\left(A * f \mid Z_{T}=(x, 0)\right)
$$

It is shown in [11] that $\lim _{T \rightarrow \infty} S_{A}^{T}=S_{A}$ exists in $L^{2}\left(\mathbb{R}^{n}\right)$. Moreover,

$$
\begin{equation*}
\left\|S_{A} f(x)\right\|_{p} \leq\left(p^{*}-1\right)\|A\|\|f\|_{p} \tag{1.20}
\end{equation*}
$$

If $A^{(i, j)}$ is defined by

$$
a_{l, m}^{(i, j)}=\left\{\begin{array}{lr}
-1 & l=i, m=j \\
0 & \text { otherwise }
\end{array}\right\}
$$

then $S_{A}$ is the second order Riesz transform, $R_{i} R_{j}$. By (1.7), this easily leads us to the conclusion that $\|B\|_{p} \leq 2\left(p^{*}-1\right)$. As with the operators arising from background radiation, if $A(x, y)=A(y)$ is independent of $x$, then $S_{A}$ is a Fourier multiplier. Furthermore, we may again find a kernel so that

$$
S_{A} f(x)=\int_{\mathbb{R}^{n}} \mathcal{K}_{A}(x, \tilde{x}) f(\tilde{x}) d \tilde{x},
$$

where

$$
\mathcal{K}_{A}(x, \tilde{x})=\int_{0}^{\infty} \int_{\mathbb{R}^{n}} A(\bar{x}, t) \nabla_{x} h_{t}(\bar{x}-\tilde{x}) \nabla_{x} h_{t}(\bar{x}-x) d \bar{x} d t .
$$

### 1.6.3 Martingale Transforms with respect to General Lévy Processes

In [9] and [5], the construction discussed in the previous subsection was generalized by replacing Brownian motion with more general Lévy processes. This results in a large class of Fourier multipliers, with formulas given in terms of the characteristics of the Lévy process, which are bounded on $L^{p}\left(\mathbb{R}^{n}\right), 1<p<\infty$. Let $\nu$ be a Lévy measure on $\mathbb{R}^{n}, \varphi$ a complex-valued function on $\mathbb{R}^{n}$ with $\|\varphi\|_{\infty} \leq 1, \mu$ a finite Borel measure on $\mathbb{S}^{n-1}$, and $\psi$ a complex-valued function on $\mathbb{S}^{n-1}$ with $\|\psi\|_{\infty} \leq 1$. Define $m_{\mu, \nu}(\xi)$ by

$$
\begin{equation*}
m_{\mu, \nu}(\xi)=\frac{\int_{\mathbb{R}^{n}}(\cos (\xi \cdot z)-1) \varphi(z) \nu(d z)+A \xi \cdot \xi}{\int_{\mathbb{R}^{n}}(\cos (\xi \cdot z)-1) \nu(d z)+B \xi \cdot \xi} \tag{1.21}
\end{equation*}
$$

where

$$
A=\left(\int_{\mathbb{S}^{n-1}} \theta_{i} \theta_{j} \psi(\theta) d \mu(\theta)\right)_{1 \leq i, j \leq n} \quad \text { and } \quad B=\left(\int_{\mathbb{S}^{n-1}} \theta_{i} \theta_{j} d \mu(\theta)\right)_{1 \leq i, j \leq n}
$$

Note that $(\cos (\xi \cdot z)-1)=\Re\left(e^{i \xi \cdot z}-1-i(\xi \cdot z) \mathbb{I}_{(|z|<1)}\right)$. Therefore, $m_{\mu, \nu}$ may be interpreted as a "modulation" of the real part of the Lévy exponent of some process, $X_{t}$, divided by the "unmodulated" real part of the Lévy exponent of $X_{t}$. The primary result of [5] is to show that $m_{\mu, \nu}$ a bounded multiplier on $L^{p}\left(\mathbb{R}^{n}\right)$ for all $1<p<\infty$ and

$$
\left\|T_{m_{\mu, \nu}} f\right\|_{p} \leq\left(p^{*}-1\right)\|f\|_{p} \quad \text { for all } f \in L^{p}\left(\mathbb{R}^{n}\right)
$$

We will now give a brief summary of how this multiplier is obtained in the case where $\mu=0$ and $\nu$ is symmetric and finite, which corresponds to $X_{t}$ being a compound Poisson process. (The general case can then be proved by symmetrization and approximation arguments. See [5] for details.) Similarly to [11], we fix $T>0$, let $\left(Z_{t}\right)_{0 \leq t \leq T}=\left(X_{t}, T-t\right)_{0 \leq t \leq T}$, and let $V_{f}(x, t)=P_{t} f(x)=\mathbb{E}^{T}\left(f\left(X_{t}+x\right)\right)$. It is shown in [9] that $V_{f}\left(Z_{t}\right)$ is a martingale, with $\sup _{t}\left\|V_{f}\left(Z_{t}\right)\right\|_{p}=\|f\|_{p}$ for all $1<p<\infty$, and by the generalized Itô's formula (see for example [36])

$$
V_{f}\left(Z_{t}\right)-V_{f}\left(Z_{0}\right)=\int_{0}^{t+} \int_{\mathbb{R}^{n}}\left[V_{f}\left(Z_{s-}+z\right)-V_{f}\left(Z_{s-}\right)\right] \tilde{N}(d s, d z),
$$

where $Z_{s-}=\lim _{u \nearrow_{s}} Z_{u}$, and $\tilde{N}$ is the so-called compensator, defined for each fixed $t>0$ on Borel sets of $\mathbb{R}^{n}$ by

$$
\tilde{N}(t, A)=N(t, A)-t \nu(A)
$$

where $N$ is the Poisson random measure that describes the jumps of $X_{t}$, i.e.

$$
N(t, A)=\left|\left\{s: 0 \leq s \leq t, X_{s}-X_{s-} \in A\right\}\right| .
$$

Therefore if $\varphi: \mathbb{R}^{n} \rightarrow \mathbb{C}$ with $\|\varphi\|_{\infty} \leq 1$, we can define the martingale transform of $V_{f}\left(Z_{t}\right)$ by $\varphi$ as

$$
\varphi * V_{f}\left(Z_{t}\right)=\int_{0}^{t+} \int_{\mathbb{R}^{n}}\left[V_{f}\left(Z_{s-}+z\right)-V_{f}\left(Z_{s-}\right)\right] \varphi(z) \tilde{N}(d s, d z)
$$

The quadratic variations of $V_{f}\left(Z_{t}\right)$ and $\varphi * V_{f}\left(Z_{t}\right)$ are given by

$$
\left[V_{f}(Z)\right]_{t}=\int_{0}^{t+} \int_{\mathbb{R}^{n}}\left|V_{f}\left(Z_{s-}+z\right)-V_{f}\left(Z_{s-}\right)\right|^{2} N(d s, d z)
$$

and

$$
\left[\varphi * V_{f}(Z)\right]_{t}=\int_{0}^{t+} \int_{\mathbb{R}^{n}}\left|V_{f}\left(Z_{s-}+z\right)-V_{f}\left(Z_{s-}\right)\right|^{2}|\varphi(z)|^{2} N(d s, d z)
$$

Therefore, $\varphi * V_{f}\left(Z_{t}\right)$ is differentially subordinate to $V_{f}\left(Z_{t}\right)$ and

$$
\sup _{t}\left\|\varphi * V_{f}\left(Z_{t}\right)\right\|_{p} \leq\left(p^{*}-1\right)\|f\|_{p}
$$

A projection operator can be defined by

$$
S_{\varphi}^{T} f(x)=\mathbb{E}^{T}\left(\varphi * V_{f}\left(Z_{T}\right) \mid Z_{T}=(x, 0)\right)
$$

and we again have that

$$
\left\|S_{\varphi}^{T} f(x)\right\|_{p} \leq\left(p^{*}-1\right)\|f\|_{p}
$$

It is shown that as $T \rightarrow \infty$, a limiting operator, $S_{\varphi}$, exists and satisfies the bound

$$
\left\|S_{\varphi} f(x)\right\|_{p} \leq\left(p^{*}-1\right)\|f\|_{p}
$$

Moreover, $S_{\varphi}$ is a Fourier multiplier and $\widehat{S_{\varphi} f}(\xi)=m_{\mu, \nu}(\xi) \widehat{f}(\xi)$.
A particularly interesting class of operators occurs when we take $X_{t}$ to be the symmetric $\alpha$-stable process with $0<\alpha<2$ and assume that $\varphi$ is homogeneous of order zero. In polar coordinates, we may write $d \nu(z)=C_{n, \alpha} r^{-1-\alpha} d r d \sigma(\theta)$ where $C_{n, \alpha}$ is a constant chosen so that

$$
\rho(\xi)=\int_{\mathbb{R}^{n}}(\cos (\xi \cdot z)-1) d \nu(z)=-|\xi|^{\alpha} .
$$

In this case, the numerator of (1.21) is given by

$$
\begin{aligned}
C_{n, \alpha} \int_{\mathbb{R}^{n}}(\cos (\xi \cdot z)-1) \varphi(z) d \nu(z) & =C_{n, \alpha} \int_{\mathbb{S}^{n-1}} \varphi(\theta) \int_{0}^{\infty} \cos (r \xi \cdot \theta) r^{-1-\alpha} d r d \sigma(\theta) \\
& =C_{n, \alpha} \int_{\mathbb{S}^{n-1}} \varphi(\theta)|\xi \cdot \theta|^{\alpha} \int_{0}^{\infty} \cos (s) s^{-1-\alpha} d s d \sigma(\theta) \\
& =C_{n, \alpha}^{\prime} \int_{\mathbb{S}^{n-1}} \varphi(\theta)|\xi \cdot \theta|^{\alpha} d \sigma(\theta)
\end{aligned}
$$

Therefore, the corresponding multiplier is given by

$$
m_{\alpha}(\xi)=\frac{\int_{\mathbb{S}^{n-1}}|\xi \cdot \theta|^{\alpha} \varphi(\theta) d \sigma(\theta)}{\int_{\mathbb{S}^{n-1}}|\xi \cdot \theta|^{\alpha} d \sigma(\theta)} .
$$

If we set $n=2$ and choose $\varphi(\theta)=e^{-2 i \arg \theta}$, then it is shown in [5] that $m_{\alpha}(\xi)=\frac{\alpha}{\alpha+2} \frac{\bar{\xi}}{\xi}$. Therefore, for all $0<\alpha<2$ and all $f \in L^{p}\left(\mathbb{R}^{n}\right)$

$$
\|B f\|_{p} \leq \frac{\alpha+2}{\alpha}\left(p^{*}-1\right)\|f\|_{p}
$$

Letting $\alpha \nearrow 2$, we recover the bound $\|B\|_{p} \leq 2\left(p^{*}-1\right)$.
The condition $0<\alpha<2$ is natural from a probabilistic prospective. Otherwise, the measure $d \nu(z)=\frac{C_{n, \alpha}}{|z|^{n+\alpha}}$ is not a Lévy measure on $\mathbb{R}^{n}$. However, for any $r>0$, the multiplier

$$
\begin{equation*}
m_{r}(\xi)=\frac{\int_{\mathbb{S}^{n-1}}|\xi \cdot \theta|^{r} \varphi(\theta) d \sigma(\theta)}{\int_{\mathbb{S}^{n-1}}|\xi \cdot \theta|^{r} d \sigma(\theta)} \tag{1.22}
\end{equation*}
$$

satisfies $\left\|m_{r}\right\|_{\infty} \leq 1$. Therefore, $T_{m_{r}}$ is a bounded operator on $L^{2}\left(\mathbb{R}^{n}\right)$. Furthermore, for any $r>0$, if we choose $\varphi(\theta)=e^{-2 i \arg \theta}$, the formula $T_{m_{r}} f(x)=\frac{r}{r+2} B f(x)$ is valid for all $f \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$. Therefore, if we could prove conjecture (1), stated below, then letting $r \rightarrow \infty$ it would follow that $\|B\|_{p} \leq p^{*}-1$, and therefore the celebrated conjecture of Iwaniec would be proved. This motivated the following conjecture of Bañuelos which first appeared in [4].

Conjecture 1 Let $n \geq 2,0<r<\infty, \varphi \in L^{\infty}\left(\mathbb{S}^{n-1}\right),\|\varphi\|_{\infty} \leq 1$, and let $m_{r}$ be defined as in (1.22). Then the corresponding operator, $T_{m_{r}}$, is bounded on $L^{p}\left(\mathbb{R}^{n}\right)$ for all $1<p<\infty$ and

$$
\left\|T_{m_{r}} f\right\|_{p} \leq\left(p^{*}-1\right)\|f\|_{p}, \quad \text { for all } f \in L^{p}\left(\mathbb{R}^{n}\right)
$$

This is a very strong conjecture since it includes Iwaniec's conjecture, which has remained unproved for over thirty years, as a special case.

### 1.7 Calderón-Zygmund Operators Arising from Martingale Transforms, Statement of Results

The constructions used in [8] and [11] give very good constants for the $L^{p}\left(\mathbb{R}^{n}\right)$ boundedness, $1<p<\infty$, of the operators which are constructed there as the projection of martingale transforms. However, using purely probabilistic methods, we are not able to get any information about the behavior of these operators on $L^{1}\left(\mathbb{R}^{n}\right)$. The purpose of chapter 2 is to show that these operators are Calderón-Zygmund operators and therefore are weak-type $(1,1)$. Specifically, we prove the following theorem. In the case that $\alpha=1$ or 2 , these operators are the conditional expectations of martingale transforms which were used in [8] and [11] respectively. (See subsections 1.6.1 and 1.6.2.)

Theorem 1.7.1 Let $0<\alpha \leq 2$. Let $\left(X_{t}\right)_{t>0}$ be a symmetric $\alpha$-stable process on $\mathbb{R}^{n}$ and let $\varphi$ denote the density of $X_{1}$. For $y \geq 0$, let $\varphi_{y}(x)=\frac{1}{y^{n}} \varphi\left(\frac{x}{y}\right)$. Let $A(x, y)=$ $\left(a^{i, j}(x, y)\right)$ be an $(n+1) \times(n+1)$ matrix-valued function with

$$
\begin{equation*}
\|A\|=\left\|\sup _{|v| \leq 1}(|A(x, y) v|)\right\|_{L^{\infty}\left(\mathbb{R}^{n} \times[0, \infty)\right)}<\infty . \tag{1.23}
\end{equation*}
$$

Further assume that $a^{i, j}(x, y)=a^{i, j}(y)$ is independent of $x$ whenever $i$ or $j=n+1$. Consider the kernel

$$
\begin{equation*}
K_{A}(x, \tilde{x})=\int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 y A(\bar{x}, y) \nabla \varphi_{y}(\bar{x}-\tilde{x}) \nabla \varphi_{y}(\bar{x}-x) d \bar{x} d y \tag{1.24}
\end{equation*}
$$

where $\nabla=\left(\partial_{x_{1}}, \ldots, \partial_{x_{n}}, \partial_{y}\right)$. Then the operator

$$
T_{A} f(x)=\int_{\mathbb{R}^{n}} K(x, \tilde{x}) f(\tilde{x}) d \tilde{x}
$$

is a CZ operator.
Remark 1 If we make the additional assumption that $a^{i, j}(y)=0$ whenever $i$ or $j=n+1$, we may also write our kernel in terms of the density of $X_{t}$, which we denote $\psi_{t}$. It is well known (see e.g. [13]) that $\psi_{t}$ obeys the scaling relation $\psi_{t}(x)=\frac{1}{t^{n / \alpha}} \psi\left(\frac{x}{t^{1 / \alpha}}\right)$ which implies $\varphi_{t^{1 / \alpha}}=\psi_{t}$. Therefore, after a simple change of variables we see that

$$
\begin{equation*}
K_{A}(x, \tilde{x})=\int_{0}^{\infty} \int_{\mathbb{R}^{n}} \frac{2}{\alpha} t^{\frac{2}{\alpha}-1} A\left(\bar{x}, t^{1 / \alpha}\right) \nabla \psi_{t}(\bar{x}-\tilde{x}) \nabla \psi_{t}(\bar{x}-x) d \bar{x} d t . \tag{1.25}
\end{equation*}
$$

The reason why we need the assumption that $a^{i, j}(y)=0$ whenever $i$ or $j=n+1$ is because these entries correspond to "vertical" derivatives with respect to the dilation parameter $t$, and the change of variables $y=t^{1 / \alpha}$ does not commute with the taking of vertical derivatives.

### 1.8 A Method of Rotations for Lévy Multipliers, Statement of Results

The main results of chapter 3 are two theorems which are partial solutions to Conjecture 1. The probabilistic methods used in [9] and [5] do not apply when $r \geq 2$. Instead, we will study $T_{m_{r}}$ by analytic methods which make use of the Marcinkiewicz mutliplier theorem and the Hörmander-Mikhlin multiplier theorem (see section 1.3).

Theorem 1.8.1 Let $n \geq 2,0<r<\infty, \varphi \in L^{\infty}\left(\mathbb{S}^{n-1}\right),\|\varphi\|_{\infty} \leq 1$, and let $m_{r}$ be defined as in (1.22). Then the corresponding operator, $T_{m_{r}}$, is bounded on $L^{p}\left(\mathbb{R}^{n}\right)$ for all $1<p<\infty$ and

$$
\left\|T_{m_{r}} f\right\|_{p} \leq C_{n}\left(p^{*}-1\right)^{6 n} \frac{\Gamma\left(\frac{r+n}{2}\right)}{\Gamma\left(\frac{r+1}{2}\right)}\|f\|_{p}, \quad \text { for all } f \in L^{p}\left(\mathbb{R}^{n}\right)
$$

where $C_{n}$ is a constant which depends only on $n$.
Remark 2 Sterling's formula implies that if $a>0$

$$
\frac{\Gamma(x+a)}{\Gamma(x)}=O\left(x^{a}\right) \quad \text { as } x \rightarrow \infty .
$$

Therefore,

$$
\frac{\Gamma\left(\frac{r+n}{2}\right)}{\Gamma\left(\frac{r+1}{2}\right)}=O\left(r^{(n-1) / 2}\right) \quad \text { as } r \rightarrow \infty .
$$

In the case that $r$ is sufficiently large, we can use the Hörmander-Mikhlin multiplier theorem to obtain estimates on the $L^{p}$ bounds of $T_{m_{r}}$ that are linear in $p$ as $p \rightarrow \infty$.

Theorem 1.8.2 Let $n \geq 2$ and define $n_{0}=\left\lfloor\frac{n}{2}\right\rfloor+1$. Let $n_{0} \leq r<\infty, \varphi \in L^{\infty}\left(\mathbb{S}^{n-1}\right)$, $\|\varphi\|_{\infty} \leq 1$, and let $m_{r}$ be defined as in (1.22). Then the corresponding operator, $T_{m_{r}}$, is bounded on $L^{p}\left(\mathbb{R}^{n}\right)$ for all $1<p<\infty$ and

$$
\left\|T_{m_{r}} f\right\|_{p} \leq C_{n} \max \left\{r^{n_{0}}, 1\right\}\left(p^{*}-1\right)\|f\|_{p}, \quad \text { for all } f \in L^{p}\left(\mathbb{R}^{n}\right),
$$

where $C_{n}$ is a constant depending only on $n$. Furthermore, $T_{m_{r}}$ is weak-type $(1,1)$ and

$$
\left|\left\{T_{m_{r}} f(x)>\lambda\right\}\right| \leq C_{n} \max \left\{r^{n_{0}}, 1\right\} \frac{\|f\|_{1}}{\lambda}
$$

Remark 3 Comparing the estimates in theorem 1.8 .1 and theorem 1.8.2, we see that each has some advantages over the other. The constants obtained in theorem 1.8.1 have slower growth as $r \rightarrow \infty$ than those obtained in theorem 1.8.2 and have the advantage of being valid for all $r>0$. On the other hand, theorem 1.8.2 gives estimates which are linear in $p$ as $p \rightarrow \infty$ and includes weak-type $(1,1)$ estimates which theorem 1.8.1 does not. This is because the proof of theorem 1.8.1 involves the method of rotations and the Marcinkiewicz multiplier theorem, neither of which give weak-type inequalities. We also remark that it is unknown if the operators which are obtained in [9] and [5] satisfy weak-type (1,1) inequalities. While it is true that martingale transforms do satisfy weak-type $(1,1)$ estimates, these estimates are not preserved under conditional expectation, as we already mentioned several times before.

## 2. Calderón-Zygmund Operators Arising from Martingale Transforms

### 2.1 The Proof of Theorem 1.7.1

Proof We need to verify that $T_{A}$ is bounded on $L^{2}\left(\mathbb{R}^{n}\right)$ and that $K_{A}$ satisfies the estimates (1.2), (1.3), and (1.4). From the definition of $T_{A}$, we observe that (1.3) and (1.4) are equivalent.

Lemma $1 T_{A}$ is bounded on $L^{2}\left(\mathbb{R}^{n}\right)$. In particular, there exists a constant $C_{n, \alpha}$, depending only on $n$ and $\alpha$, such that for all $f \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$

$$
\begin{equation*}
\left\|T_{A} f\right\|_{2} \leq C_{n, \alpha}\|A\|\|f\|_{2} \tag{2.1}
\end{equation*}
$$

Proof Let $f, g \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$. We will show that

$$
\left|\int_{\mathbb{R}^{n}} T_{A} f(x) g(x) d x\right| \leq C_{n, \alpha}\|A\|\|f\|_{2}\|g\|_{2} .
$$

Letting $u_{f}$ and $u_{g}$ denote $\varphi_{y} * f$ and $\varphi_{y} * g$ respectively,

$$
\begin{aligned}
& \left|\int_{\mathbb{R}^{n}} T_{A} f(x) g(x) d x\right| \\
& =\left|\int_{\mathbb{R}^{n}} \int_{\mathbb{R}^{n}} K_{A}(x, \tilde{x}) f(\tilde{x}) g(x) d \tilde{x} d x\right| \\
& =\left|\int_{\mathbb{R}^{n}} \int_{\mathbb{R}^{n}} \int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 y A(\bar{x}, y) \nabla \varphi_{y}(\bar{x}-\tilde{x}) \cdot \nabla \varphi_{y}(\bar{x}-x) f(\tilde{x}) g(x) d \bar{x} d y d \tilde{x} d x\right| \\
& =\left|\int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 y A(\bar{x}, y) \int_{\mathbb{R}^{n}} \nabla \varphi_{y}(\bar{x}-\tilde{x}) f(\tilde{x}) d \tilde{x} \cdot \int_{\mathbb{R}^{n}} \nabla \varphi_{y}(\bar{x}-x) g(x) d x d \bar{x} d y\right| \\
& =\left|\int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 y A(\bar{x}, y) \nabla u_{f}(\bar{x}, y) \cdot \nabla u_{g}(\bar{x}, y) d \bar{x} d y\right| \\
& \leq 2\|A\| \int_{0}^{\infty} \int_{\mathbb{R}^{n}} y^{1 / 2}\left|\nabla u_{f}(x, y)\right| y^{1 / 2}\left|\nabla u_{g}(x, y)\right| d x d y .
\end{aligned}
$$

Now by the Cauchy-Schwartz inequality and Holder's inequality,

$$
\begin{aligned}
& \int_{0}^{\infty} \int_{\mathbb{R}^{n}} y^{1 / 2}\left|\nabla u_{f}(x, y)\right| y^{1 / 2}\left|\nabla u_{g}(x, y)\right| d x d y \\
& \leq \int_{0}^{\infty}\left(\int_{\mathbb{R}^{n}} y\left|\nabla u_{f}(x, y)\right|^{2} d x\right)^{1 / 2}\left(\int_{\mathbb{R}^{n}} y\left|\nabla u_{g}(x, y)\right|^{2} d x\right)^{1 / 2} d y \\
& \leq\left(\int_{0}^{\infty} \int_{\mathbb{R}^{n}} y\left|\nabla u_{f}(x, y)\right|^{2} d x d y\right)^{1 / 2}\left(\int_{0}^{\infty} \int_{\mathbb{R}^{n}} y\left|\nabla u_{g}(x, y)\right|^{2} d x d y\right)^{1 / 2} .
\end{aligned}
$$

The proof will be complete once we show that

$$
\left(\int_{0}^{\infty} y \int_{\mathbb{R}^{n}}\left|\nabla u_{f}(x, y)\right|^{2} d x d y\right) \leq C_{n, \alpha}\|f\|_{2}^{2} .
$$

Since $\varphi$ is the density of $X_{1}$, which has characteristic function $\mathbb{E}\left(e^{i X_{1} \cdot \xi}\right)=e^{-|\xi|^{\alpha}}$, we have that $\widehat{\varphi(\xi)}=e^{-(2 \pi|\xi|)^{\alpha}}$. Therefore, we may apply Plancherel's theorem, use the scaling relation for the Fourier transform, and substitute $t=y|\xi|$, to see that

$$
\begin{aligned}
\int_{0}^{\infty} y \int_{\mathbb{R}^{n}}\left|\nabla_{x} u_{f}(x, y)\right|^{2} d x d y & =\int_{0}^{\infty} y \int_{\mathbb{R}^{n}} 4 \pi^{2}|\xi|^{2}\left|\widehat{\varphi_{y}}(\xi)\right|^{2}|\widehat{f}(\xi)|^{2} d \xi d y \\
& =C \int_{0}^{\infty} y \int_{\mathbb{R}^{n}}|\xi|^{2}|\widehat{\varphi}(\xi y)|^{2}|\widehat{f}(\xi)|^{2} d \xi d y \\
& =C \int_{0}^{\infty} t \int_{\mathbb{R}^{n}}\left|\widehat{\varphi}\left(\xi^{\prime} t\right)\right|^{2}|\widehat{f}(\xi)|^{2} d \xi d t \\
& =C \int_{\mathbb{R}^{n}}|\widehat{f}(\xi)|^{2} \int_{0}^{\infty} t e^{-2(2 \pi t)^{\alpha}} d t d \xi \\
& \leq C_{n, \alpha} \int_{\mathbb{R}^{n}}|\widehat{f}(\xi)|^{2} d \xi=C_{n, \alpha}\|f\|_{2}^{2}
\end{aligned}
$$

Likewise,

$$
\begin{aligned}
\int_{0}^{\infty} y \int_{\mathbb{R}^{n}}\left|\partial_{y} u_{f}(x, y)\right|^{2} d x d y & =\int_{0}^{\infty} y \int_{\mathbb{R}^{n}}\left|\partial_{y} \widehat{\varphi_{y}}(\xi)\right|^{2}|\widehat{f}(\xi)|^{2} d \xi d y \\
& =C \int_{0}^{\infty} y \int_{\mathbb{R}^{n}}\left|\partial_{y} \widehat{\varphi}(\xi y)\right|^{2}|\widehat{f}(\xi)|^{2} d \xi d y \\
& =C \int_{0}^{\infty} y \int_{\mathbb{R}^{n}}|\xi \cdot \nabla \widehat{\varphi}(\xi y)|^{2}|\widehat{f}(\xi)|^{2} d \xi d y \\
& \leq C \int_{0}^{\infty} y \int_{\mathbb{R}^{n}}|\xi|^{2}|\nabla \widehat{\varphi}(\xi y)|^{2}|\widehat{f}(\xi)|^{2} d \xi d y \\
& \leq C \int_{\mathbb{R}^{n}}|\xi|^{4} \int_{0}^{\infty} y^{3}|\widehat{\varphi}(\xi y)|^{2}|\widehat{f}(\xi)|^{2} d y d \xi \\
& =C \int_{\mathbb{R}^{n}}|\widehat{f}(\xi)|^{2} \int_{0}^{\infty} t\left|\widehat{\varphi}\left(\xi^{\prime} t\right)\right|^{2} d t d \xi \leq C_{n, \alpha}\|f\|_{2}^{2} .
\end{aligned}
$$

Now that we know $T_{A}$ is bounded on $L^{2}\left(\mathbb{R}^{n}\right)$, we will show that it is, in fact, a CZ operator. It suffices to show that $K_{A}^{i, j}$ satisfies (1.2) and (1.3) for $1 \leq i, j, \leq n+1$ where

$$
\begin{equation*}
K_{A}^{i, j}(x, \tilde{x})=\int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 y a^{i, j}(\bar{x}, y) \partial_{x_{i}} \varphi_{y}(\bar{x}-\tilde{x}) \partial_{x_{j}} \varphi_{y}(\bar{x}-x) d \bar{x} d y . \tag{2.2}
\end{equation*}
$$

The following lemma will be used to see that certain integrals converge.

Lemma 2 There exists a constant $C_{n, \alpha}$, depending only on $n$ and $\alpha$, such that for all $x \in \mathbb{R}^{n}, 1 \leq i, j \leq n$,

$$
\begin{align*}
|\varphi(x)| & \leq \frac{C_{n, \alpha}}{\left(1+|x|^{2}\right)^{(n+\alpha) / 2}}  \tag{2.3}\\
\left|\partial_{x_{i}} \varphi(x)\right| & \leq \frac{C_{n, \alpha}|x|}{\left(1+|x|^{2}\right)^{(n+2+\alpha) / 2}} \leq \frac{C_{n, \alpha}}{\left(1+|x|^{2}\right)^{(n+1+\alpha) / 2}} \tag{2.4}
\end{align*}
$$

and

$$
\begin{equation*}
\left|\partial_{x_{i}} \partial_{x_{j}} \varphi(x)\right| \leq \frac{C_{n, \alpha}}{\left(1+|x|^{2}\right)^{(n+2+\alpha) / 2}} . \tag{2.5}
\end{equation*}
$$

Proof Inverting the characteristic function of $X_{1}$ we see

$$
\begin{equation*}
\varphi(x)=\int_{\mathbb{R}^{n}} e^{-i x \cdot \xi} e^{-|\xi|^{\alpha}} d \xi . \tag{2.6}
\end{equation*}
$$

From this we readily see that $\varphi \in C^{\infty}\left(\mathbb{R}^{n}\right)$, so in order to show (2.4) it suffices to show that there exists a constant $C_{n, \alpha}$ so that

$$
\begin{equation*}
\left|\partial_{x_{i}} \varphi(x)\right| \leq C_{n, \alpha}|x| \tag{2.7}
\end{equation*}
$$

and

$$
\begin{equation*}
\left|\partial_{x_{i}} \varphi(x)\right| \leq \frac{C_{n, \alpha}}{|x|^{n+1+\alpha}} . \tag{2.8}
\end{equation*}
$$

Using the fact that

$$
\int_{\mathbb{R}^{n}} \xi_{i} e^{-|\xi|^{\alpha}} d \xi=0
$$

we see that

$$
\begin{aligned}
\left|\partial_{x_{i}} \varphi(x)\right| & =\left|\int_{\mathbb{R}^{n}} \xi_{i} e^{-i x \cdot \xi} e^{-|\xi|^{\alpha}} d \xi\right| \\
& =\left|\int_{\mathbb{R}^{n}} \xi_{i}\left(e^{-i x \cdot \xi}-1\right) e^{-|\xi|^{\alpha}} d \xi\right| \\
& \leq \int_{\mathbb{R}^{n}}|\xi|\left|e^{-i x \cdot \xi}-1\right| e^{-|\xi|^{\alpha}} d \xi \\
& \leq 2 \int_{\mathbb{R}^{n}}|\xi|^{2}|x| e^{-|\xi|^{\alpha}} d \xi \leq C_{n, \alpha}|x|,
\end{aligned}
$$

with the last inequality following because

$$
\left|e^{i x \cdot \xi}-1\right| \leq|\cos (x \cdot \xi)-1|+|\sin (x \cdot \xi)| \leq 2|x \cdot \xi| .
$$

Therefore (2.7) holds.
To show (2.8), we express $X_{t}$ as a process subordinated to Brownian motion. A subordinator is an a.s. increasing one-dimensional Lévy process. It is well known (see [14] for details) that there exists a subordinator, $T_{t}$, such that

$$
X_{t}=B_{T_{t}},
$$

where $B_{t}$ is a standard Brownian motion (run at twice the usual speed). By conditioning on $T_{t}$ we see that the density of $X_{t}$ is given by

$$
\psi_{t}(x)=\int_{0}^{\infty} \frac{1}{(4 \pi s)^{n / 2}} e^{-|x|^{2} / 4 s} \eta^{\alpha / 2}(t, s) d s,
$$

where $\eta^{\alpha / 2}(t, \cdot)$ is the density of $T_{t}$. Since $\varphi=\psi_{1}$, we see that

$$
\begin{aligned}
\partial_{x_{i}} \varphi(x) & =\int_{0}^{\infty} \frac{1}{(4 \pi s)^{n / 2}} \frac{x_{i}}{s} e^{-|x|^{2} / 4 s} \eta^{\alpha / 2}(1, s) d s \\
& =C_{n} \frac{x_{i}}{|x|^{n}} \int_{0}^{\infty} u^{\frac{n}{2}-1} e^{-u} \eta^{\alpha / 2}\left(1, \frac{|x|^{2}}{4 u}\right) d u .
\end{aligned}
$$

It is known (see e.g. [15]) that there exists a constant $C_{\alpha}$, depending only on $\alpha$, such that

$$
\begin{equation*}
\eta^{\alpha / 2}(t, s) \leq C_{\alpha} t s^{-1-\alpha / 2} . \tag{2.9}
\end{equation*}
$$

Therefore we have

$$
\left|\partial_{x_{i}} \varphi(x)\right| \leq \frac{C_{\alpha}}{|x|^{n+1+\alpha}} \int_{0}^{\infty} u^{(n+\alpha) / 2} e^{-u} d u
$$

so (2.8) holds. Similar computations show

$$
|\varphi(x)| \leq \frac{C_{\alpha}}{|x|^{n+\alpha}} \int_{0}^{\infty} u^{(n+\alpha-2) / 2} e^{-u} d u
$$

and

$$
\left|\partial_{x_{i}} \partial_{x_{j}} \varphi(x)\right| \leq \frac{C_{\alpha}}{|x|^{n+\alpha+2}} \int_{0}^{\infty} u^{(n+\alpha+2) / 2}(u+1) e^{-u} d u
$$

Moreover, since $\varphi$ is smooth, it and all of its all of its partial derivatives are bounded near the origin. Therefore $\varphi$ satisfies (2.3) and (2.5).

We are now poised to prove the theorem.
Case 1. Either $i$ or $j=n+1$ :
The fact that $a^{(i, j)}(x, y)=a^{(i, j)}(y)$ depends only on $y$ allows us to use the semigroup property of $\psi_{y}$. Note that

$$
\varphi_{y} * \varphi_{y}=\psi_{y^{\alpha}} * \psi_{y^{\alpha}}=\psi_{2 y^{\alpha}}=\varphi_{2^{1 / \alpha} y} .
$$

Therefore, substituting $w=\bar{x}-\tilde{x}$ we see that

$$
\begin{aligned}
\left|K^{(i, j)}(x, \tilde{x})\right| & =\left|\int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 y a^{(i, j)}(y) \partial_{x_{i}} \varphi_{y}(w) \partial_{x_{j}} \varphi_{y}(w-(x-\tilde{x})) d w d y\right| \\
& =\left|\int_{0}^{\infty} 2 y a^{(i, j)}(y) \int_{\mathbb{R}^{n}} \partial_{x_{i}} \varphi_{y}(w) \partial_{x_{j}} \varphi_{y}(w-(x-\tilde{x})) d w d y\right| \\
& =\left|\int_{0}^{\infty} 2 y a^{(i, j)}(y) \partial_{x_{i}} \partial_{x_{j}} \varphi_{2^{1 / \alpha}}(x-\tilde{x}) d y\right| \\
& \leq\left\|a^{(i, j)}\right\|_{\infty} \int_{0}^{\infty} 2 y\left|\partial_{x_{i}} \partial_{x_{j}} \varphi_{2^{1 / \alpha}}(x-\tilde{x})\right| d y .
\end{aligned}
$$

Likewise,

$$
\begin{aligned}
\left|\partial_{x_{k}} K^{(i, j)}(x, \tilde{x})\right| & =\left|\int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 y a^{(i, j)}(y) \partial_{x_{i}} \varphi_{y}(w) \partial_{x_{k}} \partial_{x_{j}} \varphi_{y}(w-(x-\tilde{x})) d w d y\right| \\
& \leq\left\|a^{(i, j)}\right\|_{\infty} \int_{0}^{\infty} 2 y\left|\partial_{x_{i}} \partial_{x_{j}} \partial_{x_{k}} \varphi_{2^{1 / \alpha}}(x-\tilde{x})\right| d y
\end{aligned}
$$

Therefore, it suffices to show that there exists a constant $C_{n, \alpha}$ so that $|K(x)| \leq$ $C_{n, \alpha} \frac{1}{|x|^{n}}$ and $\left|K^{\prime}(x)\right| \leq C_{n, \alpha} \frac{1}{|x|^{n+1}}$ for all $x \neq 0$, where

$$
K(x)=\int_{0}^{\infty} 2 y\left|\partial_{x_{i}} \partial_{x_{j}} \varphi_{2^{1 / \alpha} y}(x)\right| d y
$$

and

$$
K^{\prime}(x)=\int_{0}^{\infty} 2 y\left|\partial_{x_{i}} \partial_{x_{j}} \partial_{x_{k}} \varphi_{2^{1 / \alpha}}(x)\right| d y .
$$

$\varphi_{y}$ is homogeneous of order $-n$, so its $k$-th order partial derivatives are homogeneous of order $-n-k$. Therefore, if we make the substitution $y=|x| t$ we have

$$
\begin{aligned}
K(x) & =\int_{0}^{\infty} 2 y\left|\partial_{x_{i}} \partial_{x_{j}} \varphi_{2^{1 / \alpha}}(x)\right| d y \\
& =\int_{0}^{\infty} 2|x| t\left|\partial_{x_{i}} \partial_{x_{j}} \varphi_{2^{1 / \alpha}|x| t}\left(|x| x^{\prime}\right)\right||x| d t \\
& =\int_{0}^{\infty} 2|x| t \frac{1}{|x|^{n+2}}\left|\partial_{x_{i}} \partial_{x_{j}} \varphi_{2^{1 / \alpha} t}\left(x^{\prime}\right)\right||x| d t \\
& =\frac{1}{|x|^{n}} \int_{0}^{\infty} 2 t\left|\partial_{x_{i}} \partial_{x_{j}} \varphi_{2^{1 / \alpha} t}\left(x^{\prime}\right)\right| d t
\end{aligned}
$$

where $x^{\prime}=\frac{x}{|x|}$. Similarly,

$$
K^{\prime}(x)=\frac{1}{|x|^{n+1}} \int_{0}^{\infty} 2 t\left|\partial_{x_{i}} \partial_{x_{j}} \partial_{x_{k}} \varphi_{2^{1 / \alpha}}\left(x^{\prime}\right)\right| d t .
$$

The lemma will be proved as soon as we bound the above integrals. We have assumed that either $i$ or $j=n+1$, so we need to bound the following four integrals for any $1 \leq k, l \leq n$.

$$
\begin{array}{lr}
\int_{0}^{\infty} 2 t\left|\partial_{t} \partial_{t} \varphi_{2^{1 / \alpha}}\left(x^{\prime}\right)\right| d t, & \int_{0}^{\infty} 2 t\left|\partial_{t} \partial_{x_{k}} \varphi_{2^{1 / \alpha}}\left(x^{\prime}\right)\right| d t, \\
\int_{0}^{\infty} 2 t\left|\partial_{t} \partial_{x_{k}} \partial_{x_{l}} \varphi_{2^{1 / \alpha_{t}}}\left(x^{\prime}\right)\right| d t, & \int_{0}^{\infty} 2 t\left|\partial_{t} \partial_{t} \partial_{x_{k}} \varphi_{2^{1 / \alpha}}\left(x^{\prime}\right)\right| d t .
\end{array}
$$

We will show how to bound the first integral. The other three may be bounded by the exact same method. Recalling that $\varphi_{t}(x)=\frac{1}{t^{n}} \varphi\left(\frac{x}{t}\right)$, we see that

$$
\partial_{t} \partial_{t} \varphi_{2^{1 / \alpha} t}(x)=\frac{C_{n}^{(1)}}{t^{n+2}} \varphi\left(\frac{x}{t}\right)+\frac{C_{n}^{(2)}}{t^{n+3}} \sum_{i=1}^{n} x_{i} \partial_{x_{i}} \varphi\left(\frac{x}{t}\right)+\frac{C_{n}^{(3)}}{t^{n+4}} \sum_{i=1}^{n} \sum_{j=1}^{n} x_{i} x_{j} \partial_{x_{i}} \partial_{x_{j}} \varphi\left(\frac{x}{t}\right),
$$

where $C_{n}^{(1)}, C_{n}^{(2)}$ and $C_{n}^{(3)}$ are constants depending on $n$.

Therefore, it suffices to bound

$$
\int_{0}^{\infty} \frac{t}{t^{n+a}} \partial^{\beta} \varphi\left(\frac{x^{\prime}}{t}\right) d t
$$

when $a=2,3$, or 4 and $\beta$ is a multi-index with $|\beta|=a-2$. By (2.3), (2.4), and (2.5), we have

$$
\partial^{\beta} \varphi(x) \leq \frac{C_{n, \alpha}}{\left(1+|x|^{2}\right)^{(n+\alpha+|\beta|) / 2}},
$$

which implies

$$
\partial^{\beta} \varphi\left(\frac{x}{t}\right) \leq \frac{C_{n, \alpha} t^{n+\alpha+|\beta|}}{\left(t^{2}+|x|^{2}\right)^{(n+\alpha+|\beta|) / 2}} .
$$

Therefore,

$$
\begin{aligned}
\int_{0}^{\infty} \frac{t}{t^{n+a}} \partial^{\beta} \varphi\left(\frac{x^{\prime}}{t}\right) d t & \leq \int_{0}^{\infty} \frac{t}{t^{n+a}} \frac{t^{n+\alpha+|\beta|}}{\left(t^{2}+1\right)^{(n+\alpha+|\beta|) / 2}} d t \\
& =\int_{0}^{\infty} \frac{t^{\alpha-1}}{\left(1+t^{2}\right)^{(n+\alpha+|\beta|) / 2}} d t<\infty
\end{aligned}
$$

Case 2. $1 \leq i, j \leq n$ :
Since $a^{i, j}(x, y)$ depends on both $x$ and $y$, we are unable to use the semi-group property of $\psi_{y}$. We are, however, still able to pull out $\|A\|$ and use homogeneity. This again allows us to bound our kernel by the product of $\frac{1}{|x-\tilde{x}|^{n}}$ and an integral. As in case 1 , we start out by substituting $w=x-\tilde{x}$ to see

$$
\begin{aligned}
\left|K^{(i, j)}(x, \tilde{x})\right| & \leq\|A\| \int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 y\left|\partial_{x_{i}} \varphi_{y}(w) \| \partial_{x_{j}} \varphi_{y}(w-(\tilde{x}-x))\right| d w d y \text { and } \\
\left|\partial_{x_{k}} K^{(i, j)}(x, \tilde{x})\right| & \leq\|A\| \int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 y\left|\partial_{x_{i}} \varphi_{y}(w) \| \partial_{x_{k}} \partial_{x_{j}} \varphi_{y}(w-(\tilde{x}-x))\right| d w d y
\end{aligned}
$$

Therefore, we need to show that there exists a constant $C_{n, \alpha}$ so that $|K(x)| \leq C_{n, \alpha} \frac{1}{|x|^{n}}$ and $\left|K^{\prime}(x)\right| \leq C_{n, \alpha} \frac{1}{|x|^{n+1}}$ for all $x \neq 0$ where now

$$
K(x)=\int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 y\left|\partial_{x_{i}} \varphi_{y}(w) \| \partial_{x_{j}} \varphi_{y}(w-x)\right| d w d y
$$

and

$$
K^{\prime}(x)=\int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 y\left|\partial_{x_{i}} \varphi_{y}(w) \| \partial_{x_{k}} \partial_{x_{j}} \varphi_{y}(w-x)\right| d w d y
$$

Using homogeneity and substituting $y=t|x|$ and $w=|x| z$ we see that

$$
\begin{align*}
|K(x)| & =\int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 y\left|\partial_{x_{i}} \varphi_{y}(w) \| \partial_{x_{j}} \varphi_{y}(w-x)\right| d w d y \\
& =\int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 t|x|\left|\partial_{x_{i}} \varphi_{|x| t}(|x| z)\left\|\partial_{x_{j}} \varphi_{|x| t}\left(|x|\left(z-x^{\prime}\right)\right)\right\| x\right|^{n+1} d z d t \\
& =\frac{1}{|x|^{n}} \int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 t\left|\partial_{x_{i}} \varphi_{t}(z) \| \partial_{x_{j}} \varphi_{t}\left(z-x^{\prime}\right)\right| d z d t . \tag{2.10}
\end{align*}
$$

Similarly, we have

$$
\begin{equation*}
\left|K^{\prime}(x)\right|=\frac{1}{|x|^{n+1}} \int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 t\left|\partial_{x_{i}} \varphi_{t}(z) \| \partial_{x_{k}} \partial_{x_{j}} \varphi_{t}\left(z-x^{\prime}\right)\right| d z d t . \tag{2.11}
\end{equation*}
$$

Therefore, to complete the proof, we need to show that the integrals in (2.10) and (2.11) are convergent. (Note that a simple rotation of coordinates shows they do not depend on $x^{\prime}$.) We will show that the integral in (2.11) converges. The proof that the integral in (2.10) converges is similar.

By (2.4) and (2.5) we know that there exists a constant $C_{n, \alpha}$ so

$$
\begin{equation*}
\left|\partial_{x_{i}} \varphi(x)\right| \leq \frac{C_{n, \alpha}|x|}{\left(1+|x|^{2}\right)^{(n+2+\beta) / 2}} \tag{2.12}
\end{equation*}
$$

and

$$
\begin{equation*}
\left|\partial_{x_{i}} \partial_{x_{j}} \varphi(x)\right| \leq \frac{C_{n, \alpha}}{\left(1+|x|^{2}\right)^{(n+2+\beta) / 2}}, \tag{2.13}
\end{equation*}
$$

where $\beta=\min \left\{\alpha, \frac{1}{2}\right\}$. (The fact that $\beta \leq \frac{1}{2}$ will be used to see that a certain integral is convergent.) Therefore,

$$
\begin{equation*}
\left|\partial_{x_{i}} \varphi_{t}(x)\right| \leq \frac{C_{n} t^{\beta}|x|}{\left(t^{2}+|x|^{2}\right)^{(n+2+\beta) / 2}} \tag{2.14}
\end{equation*}
$$

and

$$
\begin{equation*}
\left|\partial_{x_{i}} \partial_{x_{j}} \varphi_{t}(x)\right| \leq \frac{C_{n} t^{\beta}}{\left(t^{2}+|x|^{2}\right)^{(n+2+\beta) / 2}} . \tag{2.15}
\end{equation*}
$$

This allows us to see that

$$
\begin{aligned}
& \int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 t\left|\partial_{x_{j}} \partial_{x_{i}} \varphi_{t}(z)\right|\left|\partial_{x_{j}} \varphi_{t}\left(z-x^{\prime}\right)\right| d z d t \\
& \leq \int_{\mathbb{R}^{n}} \int_{0}^{\infty} t \frac{t^{\beta}}{\left(|z|^{2}+t^{2}\right)^{(n+2+\beta) / 2}} \frac{\left|x^{\prime}-z\right| t^{\beta}}{\left(\left|x^{\prime}-z\right|^{2}+t^{2}\right)^{(n+2+\beta) / 2}} d t d z,
\end{aligned}
$$

so it suffices to show that $g_{1}(z)$ and $g_{2}(z)$ are integrable over $\mathbb{R}^{n}$ for

$$
\begin{aligned}
& g_{1}(z)=\int_{0}^{1} t^{1+2 \beta} \frac{1}{\left(|z|^{2}+t^{2}\right)^{(n+2+\beta) / 2}} \frac{\left|x^{\prime}-z\right|}{\left(\left|x^{\prime}-z\right|^{2}+t^{2}\right)^{(n+2+\beta) / 2}} d t \quad \text { and } \\
& g_{2}(z)=\int_{1}^{\infty} t^{1+2 \beta} \frac{1}{\left(|z|^{2}+t^{2}\right)^{(n+2+\beta) / 2}} \frac{\left|x^{\prime}-z\right|}{\left(\left|x^{\prime}-z\right|^{2}+t^{2}\right)^{(n+2+\beta) / 2}} d t .
\end{aligned}
$$

If $z_{n}$ is a sequence converging to $z$, then

$$
\begin{aligned}
\lim _{n \rightarrow \infty} g\left(z_{n}\right) & =\lim _{n \rightarrow \infty} \int_{1}^{\infty} t^{1+2 \beta} \frac{1}{\left(\left|z_{n}\right|^{2}+t^{2}\right)^{(n+\beta+2) / 2}} \frac{\left|x^{\prime}-z_{n}\right|}{\left(\left|x^{\prime}-z_{n}\right|^{2}+t^{2}\right)^{(n+\beta+2) / 2}} d t \\
& =\int_{1}^{\infty} t^{1+2 \beta} \frac{1}{\left(|z|^{2}+t^{2}\right)^{(n+2+\beta) / 2}} \frac{\left|x^{\prime}-z\right|}{\left(\left|x^{\prime}-z\right|^{2}+t^{2}\right)^{(n+2+\beta) / 2}} d t=g_{2}(z),
\end{aligned}
$$

with the middle inequality justified by the dominating convergence theorem applied to $\left|x^{\prime}-z\right| \frac{1}{t^{2 n+3}}$. Thus $g_{2}(z)$ is continuous on $\mathbb{R}^{n}$. Furthermore, for large $z$, substituting $t=|z| \tan (\theta)$ allows us to see that

$$
\begin{aligned}
g_{2}(z) & \leq C_{n, \beta} \int_{1}^{\infty} t^{1+2 \beta} \frac{|z|}{\left(|z|^{2}+t^{2}\right)^{n+\beta+2}} d t \\
& \leq C_{n, \beta} \int_{0}^{\pi / 2} \frac{|z|^{1+2 \beta} \tan ^{1+2 \beta}(\theta)|z||z| \sec ^{2}(\theta)}{|z|^{2 n+4+2 \beta} \sec ^{2 n+4+2 \beta}(\theta)} d \theta \\
& =C_{n, \beta} \frac{1}{|z|^{2 n+1}} \int_{0}^{\pi / 2} \sin ^{1+2 \beta}(\theta) \cos ^{2 n+1}(\theta) d \theta,
\end{aligned}
$$

so $g_{2}(z)$ is integrable.
Likewise, we can see that $g_{1}(z)$ is continuous on $\mathbb{R}^{n} \backslash\left\{0, x^{\prime}\right\}$ using by applying the dominating convergence theorem with $t^{2 \beta+1}|z|^{-n-\beta-2}\left|x^{\prime}-z\right|^{-n-\beta-1}$, and for large $z$ we have

$$
g_{1}(z) \leq C_{n, \beta} \frac{1}{|z|^{2 n+1}} \int_{0}^{\pi / 2} \sin ^{1+2 \beta}(\theta) \cos ^{2 n+1}(\theta) d \theta
$$

Therefore, it remains to show that $g_{1}(z)$ is integrable near 0 and $x^{\prime}$.
If $|z|<1 / 2$ and $0<t<1$, it is easy to see

$$
\frac{\left|x^{\prime}-z\right|}{\left(\left|x^{\prime}-z\right|^{2}+t^{2}\right)^{(n+2+\beta) / 2}} \leq C_{n, \beta}
$$

so again substituting $t=|z| \tan (\theta)$ we see

$$
\begin{aligned}
g_{1}(z) & \leq C_{n, \beta} \int_{0}^{1} \frac{t^{1+2 \beta}}{\left(|z|^{2}+t^{2}\right)^{(n+2+\beta) / 2}} d t \\
& \leq C_{n, \beta} \frac{1}{|z|^{n-\beta}} \int_{0}^{\pi / 2} \sin ^{2 \beta}(\theta) \cos ^{n-\beta-1}(\theta) d \theta .
\end{aligned}
$$

Since $\beta \leq \frac{1}{2}$, the last integral is finite, so $g_{1}(z)$ is integrable near 0 . A simple change of variables and a nearly identical computation shows that $g_{1}(z)$ is integrable near $x^{\prime}$, so therefore $g_{1}(z)$ is integrable on all of $\mathbb{R}^{n}$ which completes the proof.

We end this section by remarking that if $i$ or $j=n+1$, then the integral in (2.11) is divergent. This is why we need the assumption that $a^{i, j}(x, y)=a^{i, j}(y)$ in that case.

### 2.2 Remarks

Examining the proof of theorem 1.7.1, we see that the only facts we used were the homogeneity of $\varphi_{y}(x)$, the fact that $\widehat{\varphi}$ is "small enough" to cause $T_{A}$ to be bounded on $L^{2}\left(\mathbb{R}^{n}\right)$, and the bounds (2.3), (2.4), and (2.5). This immediately gives us the following corollary.

Corollary 1 Let $\phi \in C^{2}\left(\mathbb{R}^{n}\right)$ satisfy (2.3), (2.4), and (2.5), and for $y>0$, let $\phi_{y}=\frac{1}{y^{n}} \phi\left(\frac{x}{y}\right)$. Assume that there exists a constant $C$ such that for all $\xi^{\prime} \in S^{n-1}$

$$
\begin{equation*}
\int_{0}^{\infty} t \widehat{\phi}\left(t \xi^{\prime}\right)^{2} d t<C \tag{2.16}
\end{equation*}
$$

Let $A(x, y)$ be as in theorem 1.7.1. Consider the kernel

$$
K_{A}(x, \tilde{x})=\int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 y A(\bar{x}, y) \nabla \phi_{y}(\bar{x}-\tilde{x}) \nabla \phi_{y}(\bar{x}-x) d \bar{x} d y
$$

where $\nabla=\left(\partial_{x_{1}}, \ldots, \partial_{x_{n}}, \partial_{y}\right)$. Then the operator

$$
T_{A} f(x)=\int_{\mathbb{R}^{n}} K(x, \tilde{x}) f(\tilde{x}) d \tilde{x}
$$

is a CZ operator.

The key to proving lemma (2) was the fact that we could write the $\alpha$-stable process as $B_{T_{t}}$ where $T_{t}$ is the $\alpha / 2$ stable subordinator and $B_{t}$ is an independent Brownian motion (run at twice the usual speed). This motivates the following question. Let $T_{t}$ be a subordinator, let $B_{t}$ be an independent Brownian motion, and let $X_{t}=B_{T_{t}}$.

Under what conditions on $T_{t}$ does the density of $X_{1}$ satisfy the conditions of corollary 1 ?

If $X_{t}=B_{T_{t}}$ is any such process, called subordinate Brownian motion in the literature, it is well known (see for example [32]) that there exists a function $\Phi$ : $[0, \infty) \rightarrow[0, \infty)$, called the Laplace exponent of $T_{t}$, such that

$$
\begin{equation*}
e^{-\lambda T_{t}}=e^{-t \Phi(\lambda)} \tag{2.17}
\end{equation*}
$$

and that the Lévy symbol of $X_{t}$ is given by

$$
\rho(\xi)=-\Phi\left(|\xi|^{2}\right) .
$$

Inspecting the proofs of lemma 1 and lemma 2 , we see that in order for the density of $X_{1}$ to satisfy the conditions of corollary 1, it suffices to have a bound similar to (2.9) on the density of $T_{1}$, and for $\Phi(\lambda)$ to increase fast enough as $\lambda \rightarrow \infty$ for the integrals in the proofs to converge. We summarize this in the following corollary.

Corollary 2 Let $X_{t}=B_{T_{t}}$ where $T_{t}$ is a subordinator and $B_{t}$ is an independent Brownian motion run at twice the usual speed. Let $\phi$ denote the density of $X_{1}$, and for $y>0$, let $\phi_{y}(x)=\frac{1}{y^{n}} \phi\left(\frac{x}{y}\right)$. Let $\Phi$ be the Laplace exponent of $T_{t}$ and assume that there exists some $\delta>0$ so that

$$
\begin{equation*}
\Phi(\lambda) \geq O\left(\lambda^{\delta}\right), \quad \text { as } \lambda \rightarrow \infty \tag{2.18}
\end{equation*}
$$

Further assume that there exist a constants $C$ and $\gamma>0$ such that the density of $T_{1}$, $\eta(1, \cdot)$, satisfies

$$
\begin{equation*}
\eta(1, s) \leq C s^{-1-\gamma / 2} \tag{2.19}
\end{equation*}
$$

for all $s>0$. Let $A(x, y)$ be as in theorem 1.7.1 and consider the kernel

$$
K_{A}(x, \tilde{x})=\int_{0}^{\infty} \int_{\mathbb{R}^{n}} 2 y A(\bar{x}, y) \nabla \phi_{y}(\bar{x}-\tilde{x}) \nabla \phi_{y}(\bar{x}-x) d \bar{x} d y
$$

where $\nabla=\left(\partial_{x_{1}}, \ldots, \partial_{x_{n}}, \partial_{y}\right)$. Then the operator

$$
T_{A} f(x)=\int_{\mathbb{R}^{n}} K(x, \tilde{x}) f(\tilde{x}) d \tilde{x}
$$

is a CZ operator.

An interesting example of subordinate Brownian motion is provided by the so called relativistic $\alpha$-stable processes. For $0<\alpha<2, M>0$, there exists a Lévy process, $\left(X_{t}\right)_{t \geq 0}$ with symbol $\rho(\xi)=\left(|\xi|^{2}+M^{2 / \alpha}\right)^{\alpha / 2}-M$ and infinitesimal generator

$$
M-\left(-\Delta+M^{2 / \alpha}\right)^{\alpha / 2}
$$

When $\alpha=1$, this operator reduces to free-relativistic Hamiltonian which has been intensely studied because of its applications to relativistic quantum mechanics. For further background information on this process, we refer the reader to [20], [12], and the references provided in therein.

In [37] it is shown that $T_{t}$, the subordinator for $X_{t}^{m}$, has density

$$
\begin{equation*}
\eta^{m, \alpha / 2}(t, s)=e^{m t} e^{-m^{2 / \alpha} s} \eta^{\alpha / 2}(t, s), \tag{2.20}
\end{equation*}
$$

and Laplace exponent

$$
\begin{equation*}
\Phi(\lambda)=\left(\lambda+m^{2 / \alpha}\right)^{\alpha / 2}-m . \tag{2.21}
\end{equation*}
$$

Therefore, we readily see that the conditions of corollary 2 are satisfied.
The motivation of this chapter was to answer questions left open in [11] and [8]. Are the operators considered in those papers weak-type $(1,1)$ in addition to being strong-type $(p, p)$ for $1<p<\infty$ ? Proving that these operators are CZ shows that the answer to this question is, in fact, yes. However, CZ operators are also known to satisfy a number of other desirable properties. For example, they boundedly map the Hardy space $H^{1}\left(\mathbb{R}^{n}\right)$ to $L^{1}\left(\mathbb{R}^{n}\right)$ and $L^{\infty}\left(\mathbb{R}^{n}\right)$ to the space of functions with bounded mean oscillation (BMO). More precisely, if $T$ is any CZ operator, then there exist universal constants $C_{n}$ and $C_{n}^{\prime}$, which depend only on $n$, so that

$$
\|T\|_{H^{1} \rightarrow L^{1}} \leq C_{n}\left(\kappa+\|T\|_{L^{2} \rightarrow L^{2}}\right)
$$

and

$$
\|T\|_{L^{\infty} \rightarrow B M O} \leq C_{n}^{\prime}\left(\kappa+\|T\|_{L^{2} \rightarrow L^{2}}\right)
$$

where $\kappa$ is as in (1.2), (1.3) and (1.4). For details on this topic, see [26, ch. 8].

Another interesting property of CZ operators is that they are bounded on certain weighted $L^{p}$ spaces. A weight is a function $w \in L_{\mathrm{loc}}^{1}\left(\mathbb{R}^{n}\right)$ which is positive almost everywhere. The associated space $L^{p}(w), 1 \leq p<\infty$, is the collection of functions $f$ on $\mathbb{R}^{n}$ such that

$$
\|f\|_{L^{p}(w)}^{p}=\int_{\mathbb{R}^{m}}|f(x)|^{p} w(x) d x<\infty .
$$

The Muckenhoupt characteristic of $w$ is defined as

$$
\begin{equation*}
\|w\|_{A_{p}}=\sup _{Q} \frac{1}{|Q|} \int_{Q} w d x \cdot\left(\frac{1}{|Q|} \int_{Q} w^{-1 /(p-1)} d x\right)^{p-1} \tag{2.22}
\end{equation*}
$$

with the supremum taken over all cubes, $Q$. Note that when $p=2$ this becomes

$$
\|w\|_{A_{2}}=\sup _{Q} \frac{1}{|Q|} \int_{Q} w d x \cdot \frac{1}{|Q|} \int_{Q} w^{-1} d x .
$$

$w$ is said to be an $A_{p}$ weight if $\|w\|_{A_{p}}$ is finite. In this case, it is well known (see for example [26, ch. 9]) that if $T$ is a CZ operator, then there exists a constant $C_{n, p, T, w}$, depending on the $n, p, T$, and $w$, such that

$$
\|T f\|_{L^{p}(w)} \leq C_{n, p, T, w}\|f\|_{L^{p}(w)},
$$

for all $f \in L^{p}(w)$ when $1<p<\infty$. (A corresponding weak-type result holds when $p=1$.)

Recently, in [28], Hytönen proved the so called " $A_{2}$ conjecture," that $C_{n, 2, T, w}$ depends linearly on $\|w\|_{2}$, i.e., there exists a constant $C_{n, 2, T}$ such that

$$
\|T f\|_{L^{2}(w)} \leq C_{n, 2, T}\|w\|_{A_{2}}\|f\|_{L^{2}(w)},
$$

for all $f \in L^{2}(w)$. Combining this with a result of Dragičević, Grafakos, Pereyra, and Petermichl [21] shows that

$$
\|T f\|_{L^{p}(w)} \leq C_{n, p, T}\|w\|_{A_{p}}^{\max \left\{1, \frac{1}{p-1}\right\}}\|f\|_{L^{p}(w)}
$$

for all $f \in L^{p}(w)$. For more information weighted $L^{p}$ spaces and the $A_{2}$ conjecture, see [26, ch. 9] and [28].

The operators considered in [11] are generalized in [1] by taking the projections of martingales transforms involving more general Lévy processes in place of Brownian
motion. These more general operators are shown in these papers to obey the same " $p$ * -1 " strong-type bound for $1<p<\infty$ as the operators from [11]. In the current paper we have shown that the operators considered in [11] are CZ operators, and therefore are also weak-type $(1,1)$. It would be interesting to know if the same is true of the operators studied in [1].

## 3. A Method of Rotations for Lévy Multipliers

### 3.1 The proof of theorem 1.8.1

The main idea of the proof is to use a method of rotations to write $T_{m_{r}}$ as the weighted average of multipliers which can be studied using the Marcinkiewicz multiplier theorem.

Proof We first observe (see [26] Appendix D, p. 443) that

$$
\begin{equation*}
\int_{\mathbb{S}^{n-1}}|\xi \cdot \theta|^{r} d \sigma(\theta)=A_{n, r}|\xi|^{r} \tag{3.1}
\end{equation*}
$$

where $A_{n, r}=\frac{1}{2 \pi^{(n-1) / 2}} \frac{\Gamma\left(\frac{1+r}{2}\right)}{\Gamma\left(\frac{n+r}{2}\right)}$. Therefore,

$$
\begin{equation*}
m_{r}(\xi)=A_{n, r}^{-1} \int_{\mathbb{S}^{n-1}} \frac{|\xi \cdot \theta|^{r}}{|\xi|^{r}} \varphi(\theta) d \sigma(\theta) . \tag{3.2}
\end{equation*}
$$

Now for $\theta \in \mathbb{S}^{n-1}$, we let $m_{\theta}(\xi)=\frac{|\xi \cdot \theta|^{r}}{\left.|\xi|\right|^{r}}$. Using (3.2), we may write $T_{m_{r}}$ as a weighted average of the $T_{m_{\theta}}$ 's. More precisely, we shall prove the following lemma.

Lemma 3 For all $f \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$,

$$
T_{m_{r}} f(x)=A_{n, r}^{-1} \int_{\mathbb{S}^{n-1}} T_{m_{\theta}} f(x) \varphi(\theta) d \sigma(\theta),
$$

for almost every $x$.

Proof Let $f$ and $g \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$. Then by Plancherel's theorem, Fubini's theorem, and the Cauchy-Schwarz inequality,

$$
\begin{aligned}
& A_{n, r}^{-1} \int_{\mathbb{R}^{n}} \int_{\mathbb{S}^{n-1}} T_{m_{\theta}} f(x) \varphi(\theta) d \sigma(\theta) g(x) d x \\
= & A_{n, r}^{-1} \int_{\mathbb{S}^{n-1}} \varphi(\theta) \int_{\mathbb{R}^{n}} T_{m_{\theta}} f(x) g(x) d x d \sigma(\theta) \\
= & A_{n, r}^{-1} \int_{\mathbb{S}^{n-1}} \varphi(\theta) \int_{\mathbb{R}^{n}} m_{\theta}(\xi) \widehat{f}(\xi) \overline{\hat{g}}(\xi) d \xi d \sigma(\theta) \\
= & A_{n, r}^{-1} \int_{\mathbb{R}^{n}} \int_{\mathbb{S}^{n-1}} m_{\theta}(\xi) \varphi(\theta) d \sigma(\theta) \widehat{f}(\xi) \bar{g}(\xi) d \xi \\
= & \int_{\mathbb{R}^{n}} \widehat{T_{m_{r}} f}(\xi) \overline{\bar{g}}(\xi) d \xi \\
= & \int_{\mathbb{R}^{n}} T_{m_{r}} f(x) g(x) d x .
\end{aligned}
$$

We will also need to estimate the $L^{p}$ boundedness of the operators $T_{m_{\theta}}$. This is accomplished by the following lemma.

Lemma 4 There exists $0<C_{n}<\infty$ such that

$$
\left\|T_{m_{\theta}} f\right\|_{p} \leq C_{n}\left(p^{*}-1\right)^{6 n}\|f\|_{p},
$$

for all $f \in L^{p}\left(\mathbb{R}^{n}\right) . C_{n}$ depends only on $n$ and, in particular, does not depend on $r$ or $\theta$.

Before proving lemma 4, we will first show how it is used to give a simple proof of Theorem 1.8.1. By Minkowski's integral inequality,

$$
\begin{aligned}
\left\|T_{m_{r}} f\right\|_{p} & =A_{n, r}^{-1}\left(\int_{\mathbb{R}^{n}}\left|\int_{\mathbb{S}^{n-1}} \varphi(\theta) T_{m_{\theta}} f(x) d \sigma(\theta)\right|^{p} d x\right)^{1 / p} \\
& \leq A_{n, r}^{-1} \int_{\mathbb{S}^{n}-1}\left(\int_{\mathbb{R}^{n}}|\varphi(\theta)|^{p}\left|T_{m_{\theta}} f(x)\right|^{p} d x\right)^{1 / p} d \sigma(\theta) \\
& =A_{n, r}^{-1} \int_{\mathbb{S}^{n-1}}|\varphi(\theta)|\left(\int_{\mathbb{R}^{n}}\left|T_{m_{\theta}} f(x)\right|^{p} d x\right)^{1 / p} d \sigma(\theta) \\
& \leq A_{n, r}^{-1} \int_{\mathbb{S}^{n-1}}\left\|T_{m_{\theta}} f\right\|_{p} d \sigma(\theta) \\
& \leq A_{n, r}^{-1} C_{n}\left(p^{*}-1\right)^{6 n} \omega_{n-1}\|f\|_{p},
\end{aligned}
$$

where $\omega_{n-1}$ is the surface area of $\mathbb{S}^{n-1}$. Therefore, theorem 1.8 .1 is proved.

We shall now prove lemma 4
Proof For $\theta$ in $\mathbb{S}^{n-1}$, let $R$ be a rotation such that $R \theta=e_{1}$ and for $f \in L^{p}$ let $g(x)=$ $f\left(R^{-1} x\right)$. Then a simple change for variables shows that $T_{m_{\theta}} f(x)=T_{m_{e_{1}}} g(R x)$. Therefore, it suffices to show that

$$
\left\|T_{m_{e_{1}}} f\right\|_{p} \leq C_{n}\left(p^{*}-1\right)^{6 n}\|f\|_{p} \quad \text { for all } f \in L^{p}\left(\mathbb{R}^{n}\right)
$$

To prove this, we will show that $m_{e_{1}}$ satisfies the assumptions of theorem 1.3.1 and that we can take $K$ to be independent of $r$ in (1.8). Note that it follows from [39, p. 110] that for each fixed $r, T_{m_{e_{1}}}$ is a Marcinkiewicz multiplier, but it takes considerably more work to show that $K$ can be taken to be independent of $r$ in (1.8). $m_{e_{1}}(\xi)$ is even in each $\xi_{i}$ so it suffices to restrict attention to the region where all $\xi_{i}$ are positive. Noting that for all $A_{1}, \ldots, A_{k}>0$

$$
\int_{A_{1}}^{2 A_{1}} \cdots \int_{A_{k}}^{2 A_{k}} \frac{1}{\xi_{i_{1}} \xi_{i_{2}} \ldots \xi_{i_{k}}} d \xi_{i_{k}} \ldots d \xi_{i_{1}}=\log (2)^{k}
$$

we see that it suffices to prove there exists $C$ independent of $r$ such that

$$
\left|\partial_{i_{1}} \ldots \partial_{i_{k}} m_{e_{1}}(\xi)\right| \leq \frac{C}{\xi_{i_{1}} \xi_{i_{2}} \ldots \xi_{i_{k}}}
$$

or equivalently that

$$
\begin{equation*}
\xi_{i_{1}} \xi_{i_{2}} \ldots \xi_{i_{k}}\left|\partial_{i_{1}} \ldots \partial_{i_{k}} m_{e_{1}}(\xi)\right| \leq C . \tag{3.3}
\end{equation*}
$$

The left hand side of (3.3) is homogeneous of order zero, so it suffices to bound this quantity on the portion of the unit sphere where all $\xi_{i} \geq 0$. To do this, we will make use of two elementary lemma's which involve the use of Lagrange multipliers to bound polynomials on ellipses.

Lemma 5 Let $a, b, c, d>0$. The maximum value of

$$
f(x, y)=x^{a} y^{b}
$$

subject to the constraints $c x^{2}+d y^{2}=1, x, y \geq 0$, is given by

$$
\frac{\left(\frac{a}{c}\right)^{a / 2}\left(\frac{b}{d}\right)^{b / 2}}{(a+b)^{(a+b) / 2}}
$$

Proof It is easy to check using the method of Lagrange multipliers to show that $f$ is maximized when

$$
x^{2}=\frac{a}{c(a+b)} \quad \text { and } \quad y^{2}=\frac{b}{d(a+b)} .
$$

The result follows immediately.

Lemma 6 Let $1<k \leq n$, then the maximum value of $f(x, y, z)=(k-1) x^{2 k} y^{r}+$ $(n-k) x^{2 k-2} y^{r} z^{2}$ subject to the constraint that $g(x, y, z)=(k-1) x^{2}+y^{2}+(n-k) z^{2}=$ $1, x, y, z \geq 0$ is

$$
\frac{(2 k)^{k}}{(k-1)^{k-1}}\left(\frac{r}{2 k+r}\right)^{r / 2} \frac{1}{(2 k+r)^{k}} .
$$

Proof If $k=n$ then,

$$
f(x, y, z)=f(x, y)=(n-1) x^{2 n} y^{r} \quad \text { and } \quad g(x, y, z)=g(x, y)=(n-1) x^{2}+y^{2},
$$

so the result follows from lemma 5. If $1<k<n$, the method of Lagrange multipliers can be used to show that at any point at which $f$ achieves a local maximum, $z=0$. Therefore, the result again follows from lemma 5 .

Now, in order to verify that $m_{r}$ satisfies (3.3), we consider three cases.

Case $11 \notin\left\{i_{1}, \ldots, i_{k}\right\}$ :

By direct computation,

$$
\left|\partial_{i_{1}} \ldots \partial_{i_{k}} m_{e_{1}}(\xi)\right|=r(r+2) \ldots(r+2 k-2) \frac{\xi_{1}^{r} \xi_{i_{1}} \ldots \xi_{i_{k}}}{|\xi|^{r+2 k}}
$$

Therefore, we need to bound

$$
r(r+2) \ldots(r+2 k-2) \xi_{1}^{r} \xi_{i_{1}}^{2} \ldots \xi_{i_{k}}^{2}
$$

on the portion of the unit sphere where all coordinates are non-negative. By symmetry, it is clear that this last term is maximized when $\xi_{i_{1}}=\xi_{i_{2}}=\ldots=\xi_{i_{k}}$ and $\xi_{i}=0$, whenever $i \notin\left\{i_{1}, \ldots, i_{k}, 1\right\}$. Therefore, we are lead to the two-dimensional optimization problem of maximizing

$$
f(x, y)=x^{2 k} y^{r},
$$

subject to the constraint that $g(x, y)=k x^{2}+y^{2}=1$. By lemma 5 , the maximal value of $f$ subject to this constraint is less than

$$
C_{k}\left(\frac{1}{2 k+r}\right)^{k}
$$

Therefore, on the unit sphere

$$
r(r+2) \ldots(r+2 k-2) \xi_{j}^{r} \xi_{i_{1}}^{2} \ldots \xi_{i_{k}}^{2} \leq C_{k} \frac{r(r+2) \ldots(r+2 k-2)}{(2 k+r)^{k}} \leq C_{k}
$$

Case $2 k=1, i_{1}=1$ :

Differentiating, we see

$$
\left|\xi_{1} \partial_{1} m_{e_{1}}(\xi)\right|=r \frac{\xi_{1}^{r}}{|\xi|^{r+2}}\left(\xi_{2}^{2}+\ldots+\xi_{n}^{2}\right)
$$

and (3.3) can be verified by repeating the arguments of case 1 .

Case $3 k>1$ and $1 \in\left\{i_{1}, \ldots, i_{k}\right\}:$

Without loss of generality, we may assume $i_{k}=1$. Carrying out the computations, we see

$$
\begin{aligned}
& \left|\partial_{i_{1}} \ldots \partial_{i_{k-1}} \partial_{1} m(\xi)\right| \\
= & \left|\frac{r(r+2) \ldots(r+2 k-4) r \xi_{i_{1}} \ldots \xi_{i_{k-1}} \xi_{1}^{r-1}}{|\xi| r^{r+2 k-2}}-\frac{r(r+2) \ldots(r+2 k-2) \xi_{i_{1}} \ldots \xi_{i_{k-1}} \xi_{1}^{r+1}}{|\xi|^{r+2 k}}\right| \\
= & \frac{r(r+2) \ldots(r+2 k-4) \xi_{i_{1}} \ldots \xi_{i_{k-1}} \xi_{1}^{r-1}}{|\xi|^{r+2 k}}\left|r\left(\xi_{2}^{2}+\xi_{3}^{2}+\ldots+\xi_{n}^{2}\right)-(2 k-2) \xi_{1}^{2}\right| .
\end{aligned}
$$

Therefore, it suffices to show that there exists $C_{k}$ such that

$$
r(r+2) \ldots(r+2 k-4) \xi_{i_{1}}^{2} \ldots \xi_{i_{k-1}}^{2} \xi_{1}^{r+2}<C_{k}
$$

and

$$
r(r+2) \ldots(r+2 k-4) r \xi_{i_{1}}^{2} \ldots \xi_{i_{k-1}}^{2} \xi_{1}^{r}\left(\xi_{2}^{2}+\ldots+\xi_{n}^{2}\right)<C_{k},
$$

whenever $|\xi|=1$ and all $\xi_{i} \geq 0$. This can be done by using lemmas 5 and 6 in a manner similar to cases 1 and 2 .

Remark 4 In the case that $r=2 k$ is an even integer, we have that $T_{e_{1}}=R_{1}^{2 k}$, the $2 k$-th order Riesz transform in direction 1. Dimension free estimates for this operator were obtained by Iwaniec and Martin in [30] using a method that compared polynomials of the Riesz transforms to polynomials of the complex Riesz transforms and then in turn estimated the complex Riesz transforms by comparing them to the iterated Beurling-Ahlfors transform.

Identifying $\mathbb{C}^{n}$ with $\mathbb{R}^{2 n}$ the complex Riesz transforms are defined by

$$
C_{j}=R_{j}+i R_{n+j}
$$

for $1 \leq j \leq n$. For a polynomial $p(x)=\sum_{|\beta| \leq m} c_{\beta} x^{\beta}, p(\mathbf{R})$ and $p(\mathbf{C})$ are defined by

$$
p(\mathbf{R})=\sum_{|\beta| \leq m} c_{\beta} \mathbf{R}^{\beta} \quad \text { and } \quad p(\mathbf{C})=\sum_{|\beta| \leq m} c_{\beta} \mathbf{C}^{\beta},
$$

where $\mathbf{R}^{\beta}=R_{1}^{\beta_{1}} \circ \ldots \circ R_{n}^{\beta_{n}}$ and $\mathbf{C}^{\beta}=C_{1}^{\beta_{1}} \circ \ldots \circ C_{n}^{\beta_{n}}$. Iwaniec and Martin show that if $p_{2 k}$ is a homogeneous polynomial of degree $2 k$ we have that
$\left\|p_{2 k}(\mathbf{R})\right\|_{L^{p}\left(\mathbb{R}^{n}\right) \rightarrow L^{p}\left(\mathbb{R}^{n}\right)} \leq\left\|p_{2 k}(\mathbf{C})\right\|_{L^{p}\left(\mathbb{C}^{n}\right) \rightarrow L^{p}\left(\mathbb{C}^{n}\right)} \leq \frac{2 \Gamma(n+k)\left\|B^{k}\right\|_{p}}{k \pi^{n} \Gamma(k)} \int_{\mathbb{S}^{2 n-1}}\left|p_{2 k}(z)\right| d \sigma(z)$, where $\left\|B^{k}\right\|_{p}$ is the norm of the $k$-th iterated Beurling-Ahlfors transform on $L^{p}(\mathbb{C})$. Picking $p(x)=x_{1}^{2 k}$ and computing the integral on the right-hand side using the formulas in Appendix $D$ of [26], we see

$$
\left\|R_{1}^{2 k}\right\|_{L^{p}\left(\mathbb{R}^{n}\right) \rightarrow L^{p}\left(\mathbb{R}^{n}\right)} \leq\left\|B^{k}\right\|_{p} .
$$

The $L^{p}$ boundedness of $B^{k}$ was studied by Dragicevic, Petermichl, and Volberg in [22] where they showed that

$$
C_{1} k^{1-2 / p^{*}} p^{*} \leq\left\|B^{k}\right\|_{p} \leq C_{2} k^{1-2 / p^{*}} p^{*} .
$$

Combining this with (4) gives

$$
\left\|R_{1}^{2 k}\right\|_{L^{p}\left(\mathbb{R}^{n}\right) \rightarrow L^{p}\left(\mathbb{R}^{n}\right)} \leq C_{2} k^{1-2 / p^{*}} p^{*} .
$$

Therefore,

$$
\left\|T_{m_{r}} f\right\|_{p} \leq C_{n} \frac{\Gamma\left(\frac{n+r}{2}\right)}{\Gamma\left(\frac{n+1}{2}\right)}\left(\frac{r}{2}\right)^{1-2 / p^{*}} p^{*}\|f\|_{p}
$$

Like the bound obtained in theorem 1.8.2, this bound is linear in p. Futhermore, with $p$ fixed it has order $r^{(n+1) / 2-2 / p^{*}}$ as $r \rightarrow \infty$, which is slightly better than the bound obtained in theorem 1.8.2. However, this bound has the disadvantage of only being valid when $r$ is an even integer whereas the bound obtained in theorem 1.8.2 is valid for all sufficiently large $r$.

### 3.2 The proof of theorem 1.8.2

Proof It is clear that $\left\|m_{r}\right\|_{\infty} \leq 1$, so by (1.9) it suffices to show that

$$
\left(\sup _{R>0} R^{-n+2|\beta|} \int_{R<|\xi|<2 R}\left|\partial^{\beta} m_{r}(\xi)\right|^{2} d \xi\right)^{1 / 2} \leq C_{n} r^{|\beta|}
$$

for all multi-indexes with $|\beta| \leq n_{0}$. But since $m_{r}$ is homogeneous of order zero, we can make a change of variables and then use polar coordinates to see that

$$
\begin{aligned}
\sup _{R>0} R^{-n+2|\beta|} \int_{R<|\xi|<2 R}\left|\partial^{\beta} m_{r}(\xi)\right|^{2} d \xi & =\int_{1<|\xi|<2}\left|\partial^{\beta} m_{r}(\xi)\right|^{2} d \xi \\
& =\int_{1}^{2} t^{n-1} \int_{\mathbb{S}^{n-1}}\left|\partial^{\beta} m_{r}\left(t \xi^{\prime}\right)\right|^{2} d \sigma(\xi) d t \\
& =\int_{1}^{2} t^{n-1-2|\beta|} d t \int_{\mathbb{S}^{n-1}}\left|\partial^{\beta} m_{r}\left(\xi^{\prime}\right)\right|^{2} d \sigma(\xi) \\
& \leq C_{n} \int_{\mathbb{S}^{n-1}}\left|\partial^{\beta} m_{r}(\xi)\right|^{2} d \sigma(\xi),
\end{aligned}
$$

where $\xi^{\prime}=\frac{\xi}{|\xi|}$. Therefore, it suffices to show that for all multi-indexes $\beta$ with $|\beta| \leq n_{0}$,

$$
\begin{equation*}
\left(\int_{\mathbb{S}^{n-1}}\left|\partial^{\beta} m_{r}(\xi)\right|^{2} d \sigma(\xi)\right)^{1 / 2} \leq C_{n} r^{|\beta|} \tag{3.4}
\end{equation*}
$$

As in (3.1), we see that

$$
m_{r}(\xi)=C_{n} \frac{\left.\Gamma \frac{r+n}{2}\right)}{\Gamma\left(\frac{r+1}{2}\right)} n_{r}(\xi),
$$

where

$$
n_{r}(\xi)=\int_{\mathbb{S}^{n-1}} \frac{|\xi \cdot \theta|^{r}}{|\xi|^{r}} \varphi(\theta) d \sigma(\theta) .
$$

We will show that

$$
\left(\int_{\mathbb{S}^{n}-1}\left|\partial^{\beta} n_{r}(\xi)\right|^{2} d \sigma(\xi)\right)^{1 / 2} \leq C_{n} r^{|\beta|} \frac{\Gamma\left(\frac{r-n_{0}+1}{2}\right)}{\Gamma\left(\frac{r-n_{0}+n}{2}\right)},
$$

and so (3.4) will follow by observing that Sterling's formula implies that there exists $C_{n}$ such that for all $r \geq n_{0}$

$$
\frac{\Gamma\left(\frac{r+n}{2}\right)}{\Gamma\left(\frac{r+1}{2}\right)} \frac{\Gamma\left(\frac{r-n_{0}+1}{2}\right)}{\Gamma\left(\frac{r-n_{0}+n}{2}\right)} \leq C_{n} .
$$

For all $\theta \in \mathbb{S}^{n-1}$, let $m_{\theta}(\xi)=\frac{|\xi \cdot \theta|^{r}}{|\xi|^{r}}$ so that

$$
\partial^{\beta} n_{r}(\xi)=\int_{\mathbb{S}^{n-1}} \partial^{\beta} m_{\theta}(\xi) \varphi(\theta) d \sigma(\theta)
$$

We note that it suffices to show that for all $|\beta| \leq n_{0}$,

$$
\begin{equation*}
\left|\partial^{\beta} m_{\theta}(\xi)\right| \leq C_{n} r^{|\beta|}|\xi \cdot \theta|^{r-n_{0}} . \tag{3.5}
\end{equation*}
$$

For then we see that

$$
\begin{aligned}
\left(\int_{\mathbb{S}^{n-1}}\left|\partial^{\beta} n_{r}(\xi)\right|^{2} d \sigma(\xi)\right)^{1 / 2} & \leq C_{n} r^{|\beta|}\left(\int_{\mathbb{S}^{n-1}}\left(\int_{\mathbb{S}^{n-1}}\left|\partial^{\beta} m_{\theta}(\xi)\right| d \sigma(\theta)\right)^{2} d \sigma(\xi)\right)^{1 / 2} \\
& \leq C_{n} r^{|\beta|}\left(\int_{\mathbb{S}^{n-1}}\left(\int_{\mathbb{S}^{n-1}}|\xi \cdot \theta|^{r-n_{0}} d \sigma(\theta)\right)^{2} d \sigma(\xi)\right)^{1 / 2} \\
& =C_{n} r^{|\beta|} \frac{\Gamma\left(\frac{r-n_{0}+1}{2}\right)}{\Gamma\left(\frac{r-n_{0}+r}{2}\right)}
\end{aligned}
$$

Let $g_{\theta}(\xi)=|\xi \cdot \theta|^{r}$ and $h(\xi)=|\xi|^{-r}$ so that $m_{\theta}(\xi)=g_{\theta}(\xi) h(\xi)$. By Leibniz's rule

$$
\begin{aligned}
\left|\partial^{\beta} m_{\theta}(\xi)\right| & =\left|\sum_{\gamma \leq \beta}\binom{\beta}{\gamma} \partial^{\gamma} g_{\theta}(\xi) \partial^{\delta} h(\xi)\right| \\
& \leq C_{n} \sum_{\gamma \leq \beta}\left|\partial^{\gamma} g_{\theta}(\xi)\right|\left|\partial^{\delta} h(\xi)\right|
\end{aligned}
$$

where $\delta=\beta-\gamma$.

Letting $\gamma=\left(\gamma_{1}, \ldots, \gamma_{i}\right)$ and $\delta=\left(\delta_{1}, \ldots, \delta_{j}\right)$, we see that when $|\theta|=|\xi|=1$

$$
\begin{align*}
\left|\partial^{\gamma} g_{\theta}(\xi)\right| & =r(r-1) \ldots(r-i+1)|\xi \cdot \theta|^{(r-i)}\left|\theta_{\gamma_{1}} \ldots \theta_{\gamma_{i}}\right| \\
& \leq r^{i}|\xi \cdot \theta|^{r-n_{0}} \tag{3.6}
\end{align*}
$$

and

$$
\begin{equation*}
\left|\partial_{\delta} h(\xi)\right|=r(r+1) \ldots(r+j-1)|\xi|^{-r-2 j}\left|\xi_{\delta_{1}} \ldots \xi_{\delta_{j}}\right| \leq C_{n} r^{j} \tag{3.7}
\end{equation*}
$$

(3.5) follows immediately which completes the proof.

Remark 5 If we inspect the proof of theorem 1.8.2, we will see that if $r>n+1$, it follows from (3.6) and (3.7), that $m_{r}$ is multiplier which satisfies the estimate

$$
\begin{equation*}
|\xi|^{|\beta|}\left|\partial^{\beta} m_{r}(\xi)\right| \leq C_{r} \tag{3.8}
\end{equation*}
$$

for all multi-indexes with $|\beta| \leq n+1$. Therefore, by a result of McConnell [34], $m_{r}$ may be obtained using martingale transforms with respect to a Cauchy process.

### 3.3 The Method of Rotation for other Lévy Multipliers

We have seen that the Lévy multipliers which arise from martingale transforms with respect to $\alpha$-stable processes can be studied analytically using the method of rotations. This approach has the disadvantage that it does not allow us to obtain as good of constants as those that are obtained through probabilistic methods. However, it has the advantage of allowing us to remove the restriction that $\alpha<2$ and thereby obtain a larger class of operators which are bounded on $L^{p}\left(\mathbb{R}^{n}\right)$ for $1<p<\infty$. It is natural to wonder if this method can be applied to study the multipliers which arise from other Lévy processes and if so will it again let us remove restrictions on any relevant parameters.

Let $\left(X_{t}\right)_{t \geq 0}$ be a Lévy process whose Lévy measure $\nu$ is rotationally-symmetric and absolutely continuous with respect to the Lebesgue measure. Write $\nu$ in polar coordinates as $d \nu=v(r) d r d \sigma(\theta)$ for some function $v(r)$. Let $\varphi$ be a bounded function on $\mathbb{R}^{n}$ that is homogeneous of order zero, and consider the multiplier given by

$$
m_{\nu}(\xi)=\frac{\int_{\mathbb{R}^{n}}(\cos (\xi \cdot z)-1) \varphi(z) d \nu(z)}{\int_{\mathbb{R}^{n}}(\cos (\xi \cdot z)-1) d \nu(z)}
$$

Let $\rho(\xi)$ be the Lévy exponent corresponding to the Lévy triple $(0,0, \nu)$. Since the $\nu$ is symmetric, $\rho(\xi)$ is real, and therefore

$$
\begin{equation*}
\int_{\mathbb{R}^{n}}(\cos (\xi \cdot z)-1) d \nu(z)=\rho(\xi) \tag{3.9}
\end{equation*}
$$

To examine the numerator define $L: \mathbb{R} \rightarrow \mathbb{R}$ by

$$
\begin{equation*}
L(x)=\int_{0}^{\infty}(\cos (r x)-1) v(r) d r \tag{3.10}
\end{equation*}
$$

Then, we have that

$$
\begin{equation*}
\int_{\mathbb{R}^{n}}(\cos (\xi \cdot z)-1) \varphi(z) d \nu(z)=\int_{\mathbb{S}^{n-1}} L(\xi \cdot \theta) \varphi(\theta) d \sigma(\theta) \tag{3.11}
\end{equation*}
$$

Therefore, combining (3.9) and (3.11) we see that the multiplier which arises as the projection of martingale transforms with respect to $X_{t}$ is given by

$$
m_{\nu}(\xi)=\int_{\mathbb{S}^{n-1}} \frac{L(\xi \cdot \theta)}{\rho(\xi)} \varphi(\theta) d \sigma(\theta)
$$

Similarly to section 3.1, we set $m_{\theta}(\xi)=\frac{L(\xi \cdot \theta)}{\rho(\xi)}$ so that

$$
m_{\nu}(\xi)=\int_{\mathbb{S}^{n-1}} m_{\theta}(\xi) \varphi(\theta) d \sigma(\theta)
$$

Then repeating the arguments of section 3.1, we see that if $T_{m_{e_{1}}}$ is bounded on $L^{p}\left(\mathbb{R}^{n}\right)$, then $T_{m_{r}}$ is bound on $L^{p}\left(\mathbb{R}^{n}\right)$.

More generally, we have the following corollary.
Corollary 3 For any function $L: \mathbb{R} \rightarrow \mathbb{R}$, let $A L(\xi)=\int_{\mathbb{S}^{n-1}} L(\xi \cdot \theta) d \sigma(\theta)$. If $m_{e_{1}}(\xi)=\frac{L\left(\xi_{1}\right)}{A L(\xi)}$ is an $L^{p}$ multiplier for some $1<p<\infty$, then for all $\varphi \in L^{\infty}\left(\mathbb{S}^{n-1}\right)$.

$$
m_{L}(\xi)=\frac{\int_{\mathbb{S}^{n-1}} L(\xi \cdot \theta) \varphi(\theta) d \sigma(\theta)}{\int_{\mathbb{S}^{n-1}} L(\xi \cdot \theta) d \sigma(\theta)}
$$

is also an $L^{p}$ multiplier. In particular, if for some $C_{n, p}>0$,

$$
\left\|T_{m_{e_{1}}} f\right\|_{p} \leq C_{n, p}\|f\|_{p} \quad \text { for all } f \in L^{p}
$$

then

$$
\left\|T_{m_{L}} f\right\|_{p} \leq \omega_{n-1} C_{n, p}\|f\|_{p} \quad \text { for all } f \in L^{p}
$$

Consider now, for $0<\beta<\alpha<2$, the so-called "mixed-stable" process defined by, $Z_{t}=X_{t}+a Y_{t}$ where $X_{t}$ is a rotationally-symmetric $\alpha$-stable process, $Y_{t}$ is an independent rotationally symmetric $\beta$-stable process, and $a>0 . Z_{t}$ is a Lévy process with exponent $\rho(\xi)=-\left(|\xi|^{\alpha}+a^{\beta}|\xi|^{\beta}\right)$ and Lévy measure

$$
d \nu(z)=\left(C_{n, \alpha} r^{-1-\alpha}+C_{n, \beta} a^{\beta} r^{-1-\beta}\right) d r d \sigma(\theta) .
$$

In this case, by an argument similar to the $\alpha$-stable case, the corresponding multiplier is given by

$$
m_{\alpha, \beta}(\xi)=\frac{\int_{\mathbb{S}^{n-1}}\left(C_{n, \alpha}|\xi \cdot \theta|^{\alpha}+C_{n, \beta, a}|\xi \cdot \theta|^{\beta}\right) \varphi(\theta) d \sigma(\theta)}{\int_{\mathbb{S}^{n-1}}\left(C_{n, \alpha}|\xi \cdot \theta|^{\alpha}+C_{n, \beta, a}|\xi \cdot \theta|^{\beta}\right) d \sigma(\theta)} .
$$

It is already known that $m_{\alpha, \beta}$ is an $L^{p}$ multiplier for $1<p<\infty$ by the results of [9] and [5]. However, the method of rotations allows us to to remove the restriction that $0<\beta<\alpha<2$. More precisely, we can prove the following.

Corollary 4 Let $0<r<s<\infty$, let $C_{r}, C_{s}>0$, and let $\varphi \in L^{\infty}\left(\mathbb{R}^{n}\right)$. Then $m_{r, s}$ defined by

$$
m_{r, s}(\xi)=\frac{\int_{\mathbb{S}^{n}-1}\left(C_{r}|\xi \cdot \theta|^{r}+C_{s}|\xi \cdot \theta|^{s}\right) \varphi(\theta) d \sigma(\theta)}{\int_{\mathbb{S}^{n-1}}\left(C_{r}|\xi \cdot \theta|^{r}+C_{s}|\xi \cdot \theta|^{s}\right) d \sigma(\theta)} .
$$

is an $L^{p}$ multiplier, for all $1<p<\infty$ and

$$
\left\|T_{m_{r, s}} f\right\|_{p} \leq C_{n, r, s}\left(p^{*}-1\right)^{6 n}\|f\|_{p} \quad \text { for all } f \in L^{p}\left(\mathbb{R}^{n}\right)
$$

Proof As in the proof of theorem 1.8.1, the integral in the denominator can be computed directly and

$$
\int_{\mathbb{S}^{n}-1}\left(C_{r}|\xi \cdot \theta|^{r}+C_{s}|\xi \cdot \theta|^{s}\right) d \sigma(\theta)=C_{r}^{\prime}|\xi|^{r}+C_{s}^{\prime}|\xi|^{s}
$$

Therefore, in light of corollary 3 it suffices to show that

$$
m_{e_{1}}(\xi)=\frac{C_{r}\left|\xi_{1}\right|^{r}+C_{s}\left|\xi_{1}\right|^{s}}{C_{r}^{\prime}|\xi|^{r}+C_{s}^{\prime}|\xi|^{s}}
$$

is an Marcinkiewicz multiplier. As in the proof of lemma 4, we restrict attention to the region where all $\xi_{i}$ are non-negative, and check that $m_{e_{1}}$ satisfies (3.3). We already know that $\frac{\left|\xi_{1}\right|^{r}}{|\xi|^{r}}$ satisfies (3.3) so it suffices to show that

$$
n(\xi)=\frac{1+a\left|\xi_{1}\right|^{t}}{b+c|\xi|^{t}}
$$

satisfies (3.3) for all $a, b, c, t>0$ since it is easy to check using Leibniz's rule that the product of two multipliers which satisfy (3.3) is again a multiplier satisfying (3.3).

Applying Faá di Bruno's formula to the function $g(h(\xi))$, where $h(\xi)=|\xi|^{2}$ and $g(x)=\frac{1}{b+c x^{t / 2}}$, we see that $\partial_{i_{1}} \ldots \partial_{i_{k}} \frac{1}{b+\left.c|\xi|\right|^{\text {a }}}$ is a finite linear combination of terms of the form

$$
\left(\frac{|\xi|^{t}}{b+c|\xi|^{t}}\right)^{i} \frac{\xi_{i_{1}} \ldots \xi_{i_{k}}}{|\xi|^{2 k}} \frac{1}{b+c|\xi|^{t}}, \quad 0 \leq i \leq k .
$$

(3.3) then follows easily which completes the proof.

Consider again the relativistic $\alpha$-stable process which was introduced in section 2.2. Here we will show that the multipliers which arise from taking the projections of martingale transforms with respect to this process can be studied using the method of rotations. Unfortunately, unlike in the case of the mixed stable processes, the fact that $0<\alpha<2$ will play a crucial role in the proof. Therefore, we will not be able to remove that restriction and obtain a larger class of operators.

Corollary 5 Let $0<\alpha<2, M>0$, and $\varphi \in L^{\infty}\left(\mathbb{R}^{n}\right)$ homogeneous of order zero. Let $d \nu(z)=r^{-1-\alpha} \phi(r) d r d \theta$ be the Lévy measure corresponding to the relativistic $\alpha$-stable process with mass $M$ and let $L$ be defined as in (3.10). Then $\frac{L\left(\xi_{1}\right)}{\rho(\xi)}$ is a Marcinkiewicz multiplier and therefore, by corollary 3

$$
m_{\nu}=\frac{\int_{\mathbb{R}^{n}}(1-\cos (\xi \cdot \theta)) \varphi(\theta) d \nu(z)}{\int_{\mathbb{R}^{n}}(1-\cos (\xi \cdot \theta)) d \nu(z)}
$$

is an $L^{p}$ multiplier and

$$
\left\|T_{m_{\nu}} f\right\|_{p} \leq C_{n, \alpha}\left(p^{*}-1\right)^{6 n}\|f\|_{p}
$$

This is of course a weaker version of results already proven in [9] and [5], but nevertheless, it is interesting to observe that this result can also be obtained analytically.

Proof In [20], it is shown that the Lévy measure corresponding to $X_{t}$ can be written in polar coordinates by

$$
d \nu(z)=r^{-1-\alpha} \phi(r) d r d \sigma(\theta)
$$

where $\phi(r)$ is a bounded positive function that that satisfies

$$
\begin{equation*}
\phi(r) \leq C e^{-r} r^{(n+\alpha-1) / 2} \tag{3.12}
\end{equation*}
$$

when $r \geq 1$.
Now, by Faá di Bruno's formula, $\partial_{i_{1}} \ldots \partial_{i_{k}} \frac{1}{\rho(\xi)}$ is a finite linear combination of terms with the form

$$
\begin{equation*}
\frac{\xi_{i_{1}} \ldots \xi_{i_{k}}\left(|\xi|^{2}+M^{2 / \alpha}\right)^{\frac{\alpha}{2} j-k}}{\left(\left(|\xi|^{2}+M^{2 / \alpha}\right)^{\frac{\alpha}{2}}-M\right)^{j+1}}, \quad 0 \leq j \leq k . \tag{3.13}
\end{equation*}
$$

Therefore, we see that $\frac{1}{\rho(\xi)}$ is infinitely differentiable on $\mathbb{R}^{n} \backslash\{0\}$ and

$$
\left|\partial_{i_{1}} \ldots \partial_{i_{k}} \frac{1}{\rho(\xi)}\right| \leq O\left(\frac{1}{|\xi|^{\alpha+k}}\right) \quad \text { as }|\xi| \rightarrow \infty .
$$

Near 0 , each term in (3.13) is bounded above by

$$
\begin{aligned}
& C_{M, n, \alpha} \frac{1}{\left.\left(|\xi|^{2}+M^{2 / \alpha}\right)^{\frac{\alpha}{2}}-M\right)^{j+1}} \\
& \leq C_{M, n, \alpha} \frac{1}{\left.\left(|\xi|^{2}+M^{2 / \alpha}\right)^{\frac{\alpha}{2}}-M\right)} \leq O\left(\frac{1}{|\xi|^{2}}\right) \quad \text { as }|\xi| \rightarrow 0 .
\end{aligned}
$$

It is easy to check using the dominated convergence theorem, the mean value theorem and the fact that $r^{k-\alpha} \phi(r)$ is integrable on $(0, \infty)$ for all $k \geq 1$, that $L$ is infinitely differentiable on $(0, \infty)$. Therefore, in order to show that $\frac{L\left(\xi_{1}\right)}{\rho(\xi)}$ is a Marcinkiewicz multiplier it suffices to show that

$$
\begin{equation*}
|L(\xi)| \leq C_{\alpha} \min \left\{|\xi|^{\alpha},|\xi|^{2}\right\} \tag{3.14}
\end{equation*}
$$

and

$$
\begin{equation*}
\left|L^{\prime}(\xi)\right| \leq C_{\alpha} \min \left\{|\xi|^{\alpha-1},|\xi|\right\} . \tag{3.15}
\end{equation*}
$$

For then it will follow that $\frac{L\left(\xi_{1}\right)}{\rho(\xi)}$ satisfies (3.3) since

$$
\left|\xi_{i_{1}} \ldots \xi_{i_{k}} \partial_{i_{1}} \ldots \partial_{i_{k}} \frac{L\left(\xi_{1}\right)}{\rho(\xi)}\right|
$$

is a continuous function on $\mathbb{R}^{n} \backslash\{0\}$ which is bounded near the origin and as $|\xi| \rightarrow \infty$.
Making a change of variables, we see that

$$
\begin{aligned}
|L(x)| & =\left|\int_{0}^{\infty}(\cos (r x)-1) r^{-1-\alpha} \phi(r) d r\right| \\
& =|x|^{\alpha}\left|\int_{0}^{\infty}(\cos (s)-1) s^{-1-\alpha} \phi\left(\frac{s}{|x|}\right) d s\right| \leq C_{\alpha}|x|^{\alpha},
\end{aligned}
$$

where the last inequality uses the boundedness of $\phi$. On the other hand we can use the inequality $|\cos (x)-1| \leq x^{2}$, along with (3.12) and the boundedness of $\phi$ to see that

$$
\begin{aligned}
|L(x)| & \left.=\mid \int_{0}^{\infty}(\cos (r x)-1) r^{-1-\alpha}\right) \phi(r) d r \mid \\
& \leq|x|^{2}\left|\int_{0}^{\infty} r^{1-\alpha} \phi(r) d r\right| \leq C_{\alpha}|x|^{2}
\end{aligned}
$$

This proves (3.14). Note that the fact that $0<\alpha<2$ is needed in order for this integral to converge.

To prove (3.15) observe that

$$
L^{\prime}(x)=\int_{0}^{\infty} \sin (r x) r^{-\alpha} \phi(r) d r
$$

Using the fact that $|\sin (x)| \leq|x|$, it follows that $\left|L^{\prime}(x)\right| \leq C_{\alpha}|x|$ by mimicing the above arguments. To obtain the other part of (3.15) we make a change of variables, and use the fact that $\varphi$ is decreasing to see

$$
\begin{aligned}
\left|L^{\prime}(x)\right| & =|x|^{\alpha-1}\left|\int_{0}^{\infty} \sin (t) t^{-\alpha} \varphi\left(\frac{t}{x}\right) d t\right| \\
& =|x|^{\alpha-1} \sum_{n=0}^{\infty}(-1)^{n} \int_{n \pi}^{(n+1) \pi}\left|\sin (t) t^{-\alpha} \varphi\left(\frac{t}{x}\right)\right| d t \\
& \leq|x|^{\alpha-1}\left|\int_{0}^{\pi} \sin (t) t^{-\alpha} \varphi\left(\frac{t}{x}\right) d t\right| \\
& \leq C_{\alpha}|x|^{\alpha-1}
\end{aligned}
$$

This completes the proof of corollary (5).

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## VITA

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