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# Dimension free bounds for the Hardy-Littlewood maximal operator associated to convex sets ${ }^{(*)}$ 

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#### Abstract

This survey is based on a series of lectures given by the authors at the working seminar "Convexité et Probabilités" at UPMC Jussieu, Paris, during the spring 2013. It is devoted to maximal functions associated to symmetric convex sets in high dimensional linear spaces, a topic mainly developed between 1982 and 1990 but recently renewed by further advances.

The series focused on proving these maximal function inequalities in $L^{p}\left(\mathbb{R}^{n}\right)$, with bounds independent of the dimension $n$ and for all $p \in$ $(1,+\infty]$ in the best cases. This program was initiated in 1982 by Elias Stein, who obtained the first theorem of this kind for the family of Euclidean balls in arbitrary dimension. We present several results along this line, proved by Bourgain, Carbery and Müller during the period 19861990, and a new one due to Bourgain (2014) for the family of cubes in arbitrary dimension. We complete the cube case with a negative answer to the possible dimensionless behavior of the weak type $(1,1)$ constant, due to Aldaz, Aubrun and Iakovlev-Strömberg between 2009 and 2013.


Résumé. - Ces Notes reprennent et complètent une série d'exposés donnés par les auteurs au groupe de travail «Convexité et Probabilités» à l'UPMC Jussieu, Paris, au cours du printemps 2013. Elles sont consacrées à l'étude des fonctions maximales de type Hardy-Littlewood associées aux corps convexes symétriques dans $\mathbb{R}^{n}$. On s'intéresse tout particulièrement au comportement des constantes intervenant dans les estimations lorsque

[^0]la dimension $n$ tend vers l'infini. Ce sujet a été développé principalement entre 1982 et 1990, mais a été relancé par des avancées récentes.

Le but de la série d'exposés était de prouver des inégalités maximales dans $L^{p}\left(\mathbb{R}^{n}\right)$ avec des bornes indépendantes de la dimension $n$, pour certaines familles de corps convexes. Dans les meilleurs cas, on a pu obtenir de tels résultats pour toutes les valeurs de $p$ dans $(1,+\infty]$. Ce thème de recherche a été initié en 1982 par Elias Stein [75], qui a démontré le premier théorème de ce genre pour la famille des boules euclidiennes en dimension arbitraire, obtenant pour tout $p \in(1,+\infty]$ une borne dans $L^{p}\left(\mathbb{R}^{n}\right)$ indépendante de $n$. Nous présentons ce théorème de Stein ainsi que plusieurs autres résultats dans cette direction, démontrés par Bourgain, par Carbery et par Müller dans la période 1986-1990. En 1986, Bourgain [9] obtient une borne indépendante de $n$ valable dans $L^{2}\left(\mathbb{R}^{n}\right)$ pour tout corps convexe symétrique dans $\mathbb{R}^{n}$, puis Bourgain [10] et Carbery [21] étendent le résultat $L^{p}\left(\mathbb{R}^{n}\right)$ de Stein aux corps convexes symétriques quelconques, mais sous la condition que $p>3 / 2$. Müller [59] obtient un résultat valable pour tout $p>1$ quand un certain paramètre géométrique, lié aux volumes des projections du corps convexe sur les hyperplans, reste borné. Ce paramètre ne reste pas borné pour tous les convexes, en particulier, il tend vers l'infini pour les cubes de grande dimension. Nous donnons un théorème récent (2014) dû à Bourgain [13] qui obtient pour tout $p>1$ une borne dans $L^{p}\left(\mathbb{R}^{n}\right)$ indépendante de $n$ pour la famille des fonctions maximales associées aux cubes en dimension arbitraire. Nous complétons l'étude du cas du cube par des résultats pour la constante de type faible $(1,1)$, dus à Aldaz [1], à Aubrun [3] et à Iakovlev-Strömberg [46] entre 2009 et 2013 . À l'inverse du cas $L^{p}\left(\mathbb{R}^{n}\right), 1<p \leqslant+\infty$, cette constante de type faible ne reste pas bornée quand la dimension tend vers l'infini.

## Introduction

First defined by Hardy and Littlewood [44] in the one-dimensional setting, the Hardy-Littlewood maximal operator was generalized in arbitrary dimension by Wiener [83]. It turned out to be a powerful tool, for instance in harmonic or Fourier analysis, in differentiation theory or in singular integrals theory. It was extended to various situations, including not only homogeneous settings, as in the book of Coifman and Weiss [23], but also non-homogeneous, like noncompact symmetric spaces in works by Clerc and Stein [22] or Strömberg [78]. Also studied in vector-valued settings with the Fefferman-Stein type inequalities [33], it gave rise to several kinds of maximal operators which are now important in real analysis.

We shall denote by M the classical centered Hardy-Littlewood maximal operator, defined on the class of locally integrable functions $f$ on $\mathbb{R}^{n}$ by

$$
\begin{equation*}
(\mathrm{M} f)(x)=\sup _{r>0} \frac{1}{\left|B_{r}\right|} \int_{B_{r}}|f(x-y)| \mathrm{d} y, \quad x \in \mathbb{R}^{n} \tag{0.1}
\end{equation*}
$$

where $B_{r}$ is the Euclidean ball of radius $r$ and center 0 in $\mathbb{R}^{n}$, and $|S|$ denotes here the $n$-dimensional Lebesgue volume of a Borel subset $S$ of $\mathbb{R}^{n}$. It is well known that this nonlinear operator M is of strong type $(p, p)$ when $1<p \leqslant+\infty$ and of weak type $(1,1)$, as stated in the following famous theorem. We write $L^{p}\left(\mathbb{R}^{n}\right)$ for the $L^{p}$-space corresponding to the Lebesgue measure on $\mathbb{R}^{n}$.

Theorem 0.1 (Hardy-Littlewood maximal theorem). - Let $n$ be an integer $\geqslant 1$.
(1) For every function $f \in L^{1}\left(\mathbb{R}^{n}\right)$ and $\lambda>0$, the weak type inequality

$$
\begin{equation*}
\left|\left\{x \in \mathbb{R}^{n}:(\mathrm{M} f)(x)>\lambda\right\}\right| \leqslant \frac{C(n)}{\lambda}\|f\|_{L^{1}\left(\mathbb{R}^{n}\right)} \tag{WT}
\end{equation*}
$$

holds true, with a constant $C(n)$ depending only on the dimension $n$.
(2) Let $1<p \leqslant+\infty$. There exists a constant $C(n, p)$ such that for every function $f$ in $L^{p}\left(\mathbb{R}^{n}\right)$, one has

$$
\begin{equation*}
\|\mathrm{M} f\|_{L^{p}\left(\mathbb{R}^{n}\right)} \leqslant C(n, p)\|f\|_{L^{p}\left(\mathbb{R}^{n}\right)} . \tag{ST}
\end{equation*}
$$

The weak type inequality is optimal in the sense that $\mathrm{M} f$ is never in $L^{1}\left(\mathbb{R}^{n}\right)$, unless $f=0$ almost everywhere. Zygmund introduced the so-called " $L \log L$ class" to give a sufficient condition for the local integrability of the Hardy-Littlewood maximal function, a condition that is actually necessary, as proved by Stein [72]. The proof of Theorem 0.1 by Hardy and Littlewood was combinatorial and used decreasing rearrangements. The authors said: "The problem is most easily grasped when stated in the language of cricket, or any other game in which a player compiles a series of scores of which an average is recorded". Passing through the Vitali covering lemma, which is recalled below, has become later a standard approach.

A natural question that can be raised is the following. Could we compute the best constant in both inequalities (WT) and (ST)? This question seems to be out of reach in full generality. There is a very remarkable exception to this statement, the one-dimensional case where Melas has shown in [57] by a mixture of combinatorial, geometric and analytic arguments, that the best constant in $(\mathrm{WT})$ is $(11+\sqrt{61}) / 12$. The case $p>1$ is still open, even in the one-dimensional case, despite of substantial progress by Grafakos, Montgomery-Smith and Motrunich [41], who obtained by variational methods the best constant in (ST) for the class of positive functions on the line that are convex except at one point. The uncentered maximal operator $f \mapsto f^{*}$ is better understood [40], the uncentered maximal function $f^{*}$ being defined for every $x \in \mathbb{R}^{n}$ by

$$
\begin{equation*}
f^{*}(x)=\sup _{B \in \mathcal{B}(x)} \frac{1}{|B|} \int_{B}|f(u)| \mathrm{d} u \tag{0.2}
\end{equation*}
$$

where $\mathcal{B}(x)$ denotes the family of Euclidean balls $B$ containing $x$, with arbitrary center $y$ and radius $>d(y, x)$. It is clear that $f^{*} \geqslant \mathrm{M} f$, and the maximal theorem also holds for $f^{*}$ since any "uncentered" ball $B \in \mathcal{B}(x)$ of radius $r$ is contained in $B(x, 2 r)$, yielding the far from sharp pointwise inequality $f^{*} \leqslant 2^{n} \mathrm{M} f$.

Lacking for exact values, one may address the question of the asymptotic behavior of the constants when the dimension $n$ tends to infinity. This program was initiated at the beginning of the 80 s by Stein. In the usual proof of the Hardy-Littlewood maximal theorem based on the Vitali covering lemma, the dependence on the dimension $n$ in the weak type result is exponential, of the form $C(n)=C^{n}$ for some $C>1$. Then, by interpolation of Marcinkiewicz-type between the weak- $L^{1}$ case and the trivial $L^{\infty}$ case, one can get for the strong type in $L^{p}\left(\mathbb{R}^{n}\right)$ a constant of the form $C(n, p)=p C^{n / p} /(p-1)$, when $1<p \leqslant+\infty$ (see [39, Exercises, 1.3.3 (a)]). In [75], Stein has improved this asymptotic behavior in a spectacular fashion. Indeed, by using a spherical maximal operator together with a lifting method, he showed that for every $p>1$, one can replace the bound $C(n, p)$ in (ST) by a bound $C(p)$ independent of $n$. The detailed proof appeared in the paper [77] by Stein and Strömberg.

The use of an appropriate spherical maximal operator is now a decisive approach for bounding the $L^{p}$ norm of Hardy-Littlewood-type maximal operators independently of the dimension $n$, when $p>1$. This is the case, for instance, for the Heisenberg group [84] or for hyperbolic spaces [54]. Moreover, Stein and Strömberg proved that the weak type $(1,1)$ constant grows at most like $\mathrm{O}(n)$, and it is still unknown whether or not this constant may be bounded independently of the dimension. The proof in [77] draws on the Hopf-Dunford-Schwartz ergodic theorem, about which Stein says in [73] that it is "one of the most powerful results in abstract analysis". The strategy, which exploits the relationship between averages on balls and either the heat semi-group or the Poisson semi-group, is well explained in [24], and has been applied in several different settings $[27,52,53,55]$.

In a large part of these Notes, we shall replace Euclidean balls in the definition (0.1) of the maximal operator by other centrally symmetric convex bodies in $\mathbb{R}^{n}$ (in what follows, we shall omit "centrally" and abbreviate it as symmetric convex body). For example, replacing averages over Euclidean balls $B_{r}$ of radius $r$ by averages over $n$-dimensional cubes $Q_{r}$ with side $2 r$ gives an operator $\mathrm{M}_{Q}$ which satisfies both the weak type and strong type maximal inequalities. Indeed, since $B_{r} \subset Q_{r} \subset \sqrt{n} B_{r}$, it is obvious that $\mathrm{M}_{Q}$ is bounded in $L^{p}\left(\mathbb{R}^{n}\right)$ with $C(n, p)$ replaced by $n^{n / 2} C(n, p)$, but this painless route badly spoils the constants. Several results specific to the cube case have been obtained, as we shall indicate below.

More generally, as in Stein and Strömberg [77], one can give a symmetric convex body $C$ in $\mathbb{R}^{n}$ and introduce the maximal operator $\mathrm{M}_{C}$ associated to the convex set $C$ as follows: for every $f \in L_{\mathrm{loc}}^{1}\left(\mathbb{R}^{n}\right)$ one defines the function $\mathrm{M}_{C} f$ on $\mathbb{R}^{n}$ by

$$
\begin{align*}
\left(\mathrm{M}_{C} f\right)(x) & =\sup _{t>0} \frac{1}{|t C|} \int_{x+t C}|f(y)| \mathrm{d} y  \tag{0.3.M}\\
& =\sup _{t>0} \frac{1}{|C|} \int_{C}|f(x+t v)| \mathrm{d} v, \quad x \in \mathbb{R}^{n}
\end{align*}
$$

where $x+t C:=\{x+t c: c \in C\}$. One may also consider $\mathrm{M}_{C}$ when $C$ is not symmetric but has its centroid at 0, see Fradelizi [34, Section 1.5]. The maximal operator $\mathrm{M}_{C}$ satisfies, again, a maximal theorem of HardyLittlewood type.

Let $C$ be a symmetric convex body in $\mathbb{R}^{n}$. The weak type $(1,1)$ property for $\mathrm{M}_{C}$ can be deduced from the Vitali covering lemma: given a finite family of translated-dilated sets $x_{i}+r_{i} C, i \in I, x_{i} \in \mathbb{R}^{n}, r_{i}>0$, one can extract a disjoint subfamily $\left(x_{j}+r_{j} C\right)_{j \in J}, J \subset I$, such that each set $x_{i}+r_{i} C$, $i \in I$, of the original family is contained in the dilate $x_{j}+3 r_{j} C$ of some member $x_{j}+r_{j} C, j \in J$, of the extracted disjoint family. One may explain the constant 3 by the use of the triangle inequality for the norm on $\mathbb{R}^{n}$ whose unit ball is $C$. Passing to the Lebesgue measure in $\mathbb{R}^{n}$, this statement naturally introduces a factor $3^{n}$ corresponding to the dilation factor 3 . If $f_{C}^{*}$ denotes the corresponding uncentered maximal function of $f$ associated to $C$, then for every $\lambda>0$, one has that

$$
\begin{equation*}
\left|\left\{x \in \mathbb{R}^{n}: f_{C}^{*}(x)>\lambda\right\}\right| \leqslant \frac{3^{n}}{\lambda} \int_{\left\{f_{C}^{*}>\lambda\right\}}|f(x)| \mathrm{d} x \tag{0.4}
\end{equation*}
$$

We briefly sketch a proof, very similar to that of Doob's maximal inequality presented in Section 1.1. It is convenient here to consider that $C$ is an open subset of $\mathbb{R}^{n}$. Given an arbitrary compact subset $K$ of the open set $U_{\lambda}=$ $\left\{f_{C}^{*}>\lambda\right\}$, one applies the Vitali lemma to a finite covering of $K$ by open sets $S_{i}=x_{i}+r_{i} C$ such that $\int_{S_{i}}|f|>\lambda\left|S_{i}\right|$. A simple feature of $f_{C}^{*}$ is that each such $S_{i}$ is actually contained in $U_{\lambda}$. If $J \subset I$ corresponds to the disjoint family given by Vitali, then

$$
|K| \leqslant \sum_{j \in J}\left|x_{j}+3 r_{j} C\right|=3^{n} \sum_{j \in J}\left|x_{j}+r_{j} C\right| \leqslant \frac{3^{n}}{\lambda} \int_{U_{\lambda}}|f(x)| \mathrm{d} x
$$

implying (0.4). Next, a direct argument involving only Fubini and Hölder can give an $L^{p}$ bound, exactly as in the proof of Doob's Theorem 1.1 below, but giving a factor $3^{n}$ instead of $3^{n / p}$ obtained by interpolation. This Vitali method does not depend upon the symmetric body $C$, does not distinguish
the centered and uncentered operators, and introduces a quite unsatisfactory exponential constant.

Stein and Strömberg have greatly improved this exponential dependence in [77]. By a clever covering argument with less overlap than in Vitali's lemma, they proved that the weak type constant admits a bound of the form $\mathrm{O}(n \log n)$, and by using the Calderón-Zygmund method of rotations, they obtained for the strong type property a constant which behaves as $n p /(p-1)$. Concerning the weak type constant, Naor and Tao [60] have established the same $n \log n$ behavior for the large class of $n$-strong micro-doubling metric measure spaces (see also [25]). Several powerful results about the strong type constant for maximal functions associated to convex sets, beyond the one of Stein-Strömberg, have been established between 1986 and 1990. First of all, Bourgain proved a dimensionless theorem for general symmetric convex bodies in the $L^{2}$ case [9], applying geometrical arguments and methods from Fourier analysis. This result has been generalized to $L^{p}\left(\mathbb{R}^{n}\right)$, for all $p>3 / 2$, by Bourgain [10] and Carbery [21] in two independent papers. They both bring into play an auxiliary dyadic maximal operator, but Bourgain uses it together with square function techniques while Carbery uses multipliers associated to fractional derivatives. Detlef Müller extended in [59] the $L^{p}$ bound to every $p>1$, but under an additional geometrical condition on the family of convex sets $C$ under study. Müller also proved that for every fixed $q \in[1,+\infty)$, his condition is fulfilled by the family $\mathcal{F}_{q}$ of $\ell_{n}^{q}$ balls, $n \in \mathbb{N}^{*}$.

After Müller's article, activity in this area slowed down. Nevertheless, Bourgain recently proved in [13] that for all $p>1$, the strong type constant can be bounded independently of the dimension when we average over cubes. In order to attack this problem, Bourgain applies an arsenal of techniques, including a holomorphic semi-group theorem due to Pisier [62] and ideas inspired by martingale theory. The cube case is rather well understood since Aldaz [1] has proved that the weak type $(1,1)$ constant $\kappa_{Q, n}$ for cubes must tend to infinity with the dimension $n$. The best lower bound known at the time of our writing is due to Iakovlev-Strömberg [46] who obtained $\kappa_{Q, n} \geqslant$ $\kappa n^{1 / 4}$, improving a previous estimate $\kappa_{Q, n} \geqslant \kappa_{\varepsilon}(\log n)^{1-\varepsilon}$ for every $\varepsilon>0$, which was obtained by Aubrun [3] following the Aldaz result.

In the present survey, except for Section 9 on the Aldaz "negative" result, we shall restrict ourselves to $p>1$ and examine the strong type $(p, p)$ behavior of maximal functions associated to symmetric convex bodies in $\mathbb{R}^{n}$. We shall present the dimensionless result of Stein for Euclidean balls, the works of Bourgain, Carbery and Müller during the 80s and the recent dimensionless theorem of Bourgain for cubes. As we shall see, the proofs require a lot
of methods and tools, including multipliers, square functions, LittlewoodPaley theory, complex interpolation, holomorphic semi-groups and geometrical arguments involving convexity. The study of weak type inequalities for Hardy-Littlewood-type operators needs powerful methods as well: not only the aforementioned Hopf-Dunford-Schwartz ergodic theorem, but also sharp estimates for heat or Poisson semi-group, Iwasawa decomposition, $K$ -bi-invariant convolution-type operators, expander-type estimates...

The first two sections contain general dimension free inequalities obtained respectively by probabilistic methods or by Fourier transform methods. The Poisson semi-group plays an important rôle in Stein's book [73], and appears also in Bourgain's articles [9, 10] and in Carbery [21]. We give a presentation of this semi-group, both on the probabilistic and Fourier analytic viewpoints. The third section is about some analytic tools that are employed later on, namely, estimates for the Gamma function in the complex plane, and the complex interpolation scheme for linear operators, as developed in Stein [70]. The Stein result for Euclidean balls in arbitrary dimension is our Theorem 4.1. Section 5 is about Bourgain's $L^{2}$-theorem in arbitrary dimension $n$, stating that there exists a constant $\kappa_{2}$ independent of $n$ such that for any symmetric convex body $C$ in $\mathbb{R}^{n}$, one has

$$
\left\|\mathrm{M}_{C} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)} \leqslant \kappa_{2}\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)}
$$

for every $f \in L^{2}\left(\mathbb{R}^{n}\right)$. The next section presents Carbery's proof of the generalization to $L^{p}$ of the latter bound, obtained by Bourgain [10] and Carbery [21]. In both papers, the $L^{p}$ result for general symmetric convex bodies is proved for $p>3 / 2$ only. A theorem due to Detlef Müller [59] is given in Section 7; for families of symmetric convex sets $C$ for which a certain parameter $q(C)$ remains bounded, it extends the dimensionless $L^{p}$ bound to every $p>1$. This parameter is related to the $(n-1)$-dimensional measure of hyperplane projections of a specific volume one linear image of $C$, the so-called isotropic position. Section 8 presents the result of Bourgain about cubes in arbitrary dimension. In this special case, an $L^{p}$ bound independent of the dimension is valid for all $p>1$, although the Müller condition is not satisfied. Bourgain's proof is highly dependent on the product structure of the cube. In Section 9, we prove the Aldaz result that the weak type $(1,1)$ constant for cubes is not bounded when the dimension $n$ tends to infinity. We mention the quantitative improvement by Aubrun [3], and give a proof for the lower bound $\kappa n^{1 / 4}$ due to Iakovlev-Strömberg [46].

We have put a notable emphasis on the notion of log-concavity. We shall see that with not much more effort, most maximal theorems for convex sets generalize to symmetric log-concave probability densities. This kind of extension from convex sets to log-concave functions has attracted a lot of attention in convex geometry in recent years, see [5, 42, 49, 50] among many
others. In fact, Bourgain's estimate (5.17.B), which is crucial to all results in Section 5 and after, is only based on properties of log-concave distributions.

We have chosen a very elementary expository style. We shall give fully detailed proofs, except in the first two introductory sections. Most readers will know the contents of these sections and may start by reading Section 4. Some may be happy though to see a gentle introduction to a few points they are less familiar with. Our choice of topics in these two first sections owes a lot to Stein's monograph Topics in harmonic analysis [73]. In the next sections, we have chosen to recall and usually follow the methods from the original papers. This leads sometimes to unnecessary complications, but we shall try to give hints to other possibilities.

We believe that most of our notation is standard. We write $\lfloor x\rfloor,\lceil x\rceil$ for the floor and ceiling of a real number $x$, integers satisfying $x-1<\lfloor x\rfloor \leqslant$ $x \leqslant\lceil x\rceil<x+1$. We pay a special attention to constants independent of the dimension, for instance those appearing in results about martingale inequalities, Riesz transforms, and try to keep specific letters for these constants throughout the paper, such as $c_{p}, \rho_{p}, \ldots$ We use the letter $\kappa$ to denote a "universal" constant that does not deserve to be remembered. Most often in our Notes, "we" is a two-letter abbreviation for "the author", namely, Stein, Bourgain, Carbery, Müller and several others... We include an index and a notation index.

## 1. General dimension free inequalities, first part

This first section is devoted to general facts obtained by probabilistic methods, or merely employing the probabilistic language. We begin by reviewing the basic definitions. The functions here are real or complex valued, or they take values in a finite dimensional real or complex linear space $F$ equipped with a norm denoted by $|x|$, for every vector $x \in F$. If $\Omega$ is a set, a $\sigma$-field $\mathcal{G}$ of subsets of $\Omega$ is a family of subsets that is closed under countable unions $\bigcup_{n \in \mathbb{N}} A_{n}$, closed under taking complement $A \mapsto A^{c}$, and such that $\emptyset \in \mathcal{G}$. If $\Omega$ is a set and $\mathcal{G}$ a $\sigma$-field of subsets of $\Omega$, one says that a function $g$ on $\Omega$ is $\mathcal{G}$-measurable when for every Borel subset $B$ of the range space, the inverse image $g^{-1}(B)$, also denoted by

$$
\{g \in B\}:=\{\omega \in \Omega: g(\omega) \in B\}
$$

belongs to the collection $\mathcal{G}$.
A probability space $(\Omega, \mathcal{F}, P)$ consists of a set $\Omega$, a $\sigma$-field $\mathcal{F}$ of subsets of $\Omega$ and a probability measure $P$ on $(\Omega, \mathcal{F})$, i.e., a nonnegative $\sigma$-additive measure on $(\Omega, \mathcal{F})$ such that $P(\Omega)=1$. If a function $f$ is $\mathcal{F}$-measurable (we
say then that $f$ is a random variable) and if $f$ is $P$-integrable, the expectation of $f$ is the integral of $f$ with respect to $P$, denoted by

$$
\mathrm{E} f:=\int_{\Omega} f(\omega) \mathrm{d} P(\omega)
$$

Random variables $\left(f_{i}\right)_{i \in I}$ on $(\Omega, \mathcal{F}, P)$ are independent if for any finite subset $J \subset I$, one has $\mathrm{E}\left(\prod_{j \in J} h_{j} \circ f_{j}\right)=\prod_{j \in J} \mathrm{E}\left(h_{j} \circ f_{j}\right)$ for all nonnegative Borel functions $\left(h_{j}\right)_{j \in J}$ on the range space. The distribution of the random variable $f$ with values in $Y=\mathbb{R}, \mathbb{C}$ or $F$ is the image probability measure $\mu=f_{\#} P$, defined on the Borel $\sigma$-field $\mathcal{B}_{Y}$ of $Y$ by letting $\mu(B)=P(\{f \in B\})$ for every $B \in \mathcal{B}_{Y}$. If $\mu$ is a distribution on the Euclidean space $F$, the marginals of $\mu$ on the linear subspaces $F_{0}$ of $F$ are the distributions $\mu_{F_{0}}$ obtained from $\mu$ as images by orthogonal projection, i.e., one sets $\mu_{F_{0}}=\left(\pi_{0}\right)_{\#} \mu$ where $\pi_{0}$ is the orthogonal projection from $F$ onto $F_{0}$. If $f$ is $F$-valued and if $\mu$ is the distribution of $f$, then $\mu_{F_{0}}$ is that of $\pi_{0} \circ f$.

If $\mathcal{G}$ is a sub- $\sigma$-field of $\mathcal{F}$, the conditional expectation on $\mathcal{G}$ of an integrable function $f$ is the unique element $\mathrm{E}(f \mid \mathcal{G})$ of $L^{1}(\Omega, \mathcal{F}, P)$ possessing a $\mathcal{G}$-measurable representative $g$ such that

$$
\mathrm{E}\left(\mathbf{1}_{A} f\right)=\mathrm{E}\left(\mathbf{1}_{A} g\right)=\mathrm{E}\left(\mathbf{1}_{A} \mathrm{E}(f \mid \mathcal{G})\right)
$$

for every set $A \in \mathcal{G}$, where $\mathbf{1}_{A}$ denotes the indicator function of $A$, equal to 1 on $A$ and 0 outside. It follows that

$$
\mathrm{E}(h f)=\mathrm{E}(h \mathrm{E}(f \mid \mathcal{G})), \text { and actually } \mathrm{E}(h f \mid \mathcal{G})=h \mathrm{E}(f \mid \mathcal{G})
$$

for every bounded $\mathcal{G}$-measurable scalar function $h$ on $\Omega$. When $f$ is scalar and belongs to $L^{2}(\Omega, \mathcal{F}, P)$, the conditional expectation of $f$ on $\mathcal{G}$ is the orthogonal projection of $f$ onto the closed linear subspace $L^{2}(\Omega, \mathcal{G}, P)$ of $L^{2}(\Omega, \mathcal{F}, P)$ formed by $\mathcal{G}$-measurable and square integrable functions. When $A$ is an atom of $\mathcal{G}$, i.e., a minimal non-empty element of $\mathcal{G}$, and if $P(A)>0$, the value of $\mathrm{E}(f \mid \mathcal{G})$ on the atom $A$ is the average of $f$ on $A$, hence

$$
\mathrm{E}(f \mid \mathcal{G})(\omega)=\frac{1}{P(A)} \int_{A} f\left(\omega^{\prime}\right) \mathrm{d} P\left(\omega^{\prime}\right), \quad \omega \in A
$$

The conditional expectation operator $\mathrm{E}(\cdot \mid \mathcal{G})$ is linear and positive, i.e., it sends nonnegative functions to nonnegative functions. It follows that we have the inequality $\varphi(\mathrm{E}(f \mid \mathcal{G})) \leqslant \mathrm{E}(\varphi(f) \mid \mathcal{G})$ when the real-valued function $\varphi$ is convex on the range space of $f$. In particular, one has that $|\mathrm{E}(f \mid \mathcal{G})| \leqslant$ $\mathrm{E}(|f| \mid \mathcal{G})$, and

$$
\|\mathrm{E}(f \mid \mathcal{G})\|_{L^{p}(\Omega, \mathcal{F}, P)} \leqslant\|f\|_{L^{p}(\Omega, \mathcal{F}, P)}, \quad 1 \leqslant p<+\infty
$$

The inequality is true also when $p=+\infty$, it is easy and treated separately.

### 1.1. Doob's maximal inequality

A (discrete time) martingale on a probability space $(\Omega, \mathcal{F}, P)$ consists of a filtration, i.e., an increasing sequence $\left(\mathcal{F}_{k}\right)_{k \in I}$ of sub- $\sigma$-fields of $\mathcal{F}$ indexed by a subset $I$ of $\mathbb{Z}$, and of a sequence $\left(M_{k}\right)_{k \in I}$ of integrable functions on $\Omega$ such that for all $k, \ell \in I$ with $k \leqslant \ell$, one has

$$
M_{k}=\mathrm{E}\left(M_{\ell} \mid \mathcal{F}_{k}\right)
$$

Notice that each $M_{k}, k \in I$, is $\mathcal{F}_{k}$-measurable. If $I$ has a maximal element $N$, the martingale is completely determined by its last element $M_{N}$, since we have then that $M_{k}=\mathrm{E}\left(M_{N} \mid \mathcal{F}_{k}\right)$ for every $k \in I$. In the case of a finite field $\mathcal{F}_{k}$, the martingale condition means that the value of $M_{k}$ on each atom of $\mathcal{F}_{k}$ is the average of the values of $M_{\ell}$ on that atom, for every $\ell \in I$ with $\ell>k$. Clearly, any subsequence $\left(M_{k}\right)_{k \in J}, J \subset I$, is a martingale with respect to the filtration $\left(\mathcal{F}_{k}\right)_{k \in J}$.

Let us consider a finite martingale $\left(M_{k}\right)_{k=0}^{N}$ on $(\Omega, \mathcal{F}, P)$, with respect to a filtration $\left(\mathcal{F}_{k}\right)_{k=0}^{N}$. This martingale can be real or complex valued, or may take values in a finite dimensional normed space $F$. We introduce the maximal process $\left(M_{k}^{*}\right)_{k=0}^{N}$, which is defined by $M_{k}^{*}=\max _{0 \leqslant j \leqslant k}\left|M_{j}\right|$ for $k=0, \ldots, N$. In the vector-valued case, $\left|M_{j}\right|$ is the function assigning to each $\omega \in \Omega$ the norm of the vector $M_{j}(\omega) \in F$. We also employ the lighter notation $\|M\|_{p}$ for the norm $\|M\|_{L^{p}}$ of a function $M$ in $L^{p}(\Omega, \mathcal{F}, P)$, when $1 \leqslant p \leqslant+\infty$.

THEOREM 1.1 (Doob's inequality). - Let $\left(M_{k}\right)_{k=0}^{N}$ be a martingale (real, complex or vector-valued). For every real number $c>0$, one has that

$$
c P\left(\left\{M_{N}^{*}>c\right\}\right) \leqslant \int_{\left\{M_{N}^{*}>c\right\}}\left|M_{N}\right| \mathrm{d} P
$$

Furthermore, for every $p \in(1,+\infty]$, one has when $M_{N} \in L^{p}(\Omega, \mathcal{F}, P)$ that

$$
\begin{equation*}
\left\|M_{N}^{*}\right\|_{p} \leqslant \frac{p}{p-1}\left\|M_{N}\right\|_{p} \tag{1.1}
\end{equation*}
$$

Proof. - We cut the set $\left\{M_{N}^{*}>c\right\}$ into disjoint events $A_{0}, \ldots, A_{N}$, corresponding to the first time $k$ when $\left|M_{k}\right|>c$. Let $A_{0}=\left\{\left|M_{0}\right|>c\right\}$ and for each integer $k$ between 1 and $N$, let $A_{k}$ denote the set of $\omega \in \Omega$ such that $\left|M_{k}(\omega)\right|>c$ and $M_{k-1}^{*}(\omega) \leqslant c$. On the set $A_{k}$, we have $\left|M_{k}\right|>c$, and $A_{k}$ belongs to the $\sigma$-field $\mathcal{F}_{k}$ since $\left|M_{k}\right|$ and $M_{k-1}^{*}$ are $\mathcal{F}_{k}$-measurable, hence

$$
\begin{aligned}
c P\left(A_{k}\right) & \leqslant \int_{A_{k}}\left|M_{k}\right| \mathrm{d} P=\int_{A_{k}}\left|\mathrm{E}\left(M_{N} \mid \mathcal{F}_{k}\right)\right| \mathrm{d} P \\
& \leqslant \int_{A_{k}} \mathrm{E}\left(\left|M_{N}\right| \mid \mathcal{F}_{k}\right) \mathrm{d} P=\int_{A_{k}}\left|M_{N}\right| \mathrm{d} P
\end{aligned}
$$

On the other hand, we see that $\left\{M_{N}^{*}>c\right\}=\bigcup_{k=0}^{N} A_{k}$, union of pairwise disjoint sets, therefore

$$
\begin{align*}
c P\left(\left\{M_{N}^{*}>c\right\}\right)=\sum_{k=0}^{N} c P & \left(A_{k}\right) \\
& \leqslant \sum_{k=0}^{N} \int_{A_{k}}\left|M_{N}\right| \mathrm{d} P=\int_{\left\{M_{N}^{*}>c\right\}}\left|M_{N}\right| \mathrm{d} P \tag{1.2}
\end{align*}
$$

The result for $L^{p}$ when $1<p<+\infty$ follows. For each value $t>0$, we apply (1.2) with $c=t$, we use Fubini's theorem and Hölder's inequality, obtaining thus

$$
\begin{aligned}
\mathrm{E}\left(\left(M_{N}^{*}\right)^{p}\right) & =\mathrm{E}\left(\int_{0}^{M_{N}^{*}} p t^{p-1} \mathrm{~d} t\right)=\int_{0}^{+\infty} p t^{p-1} P\left(\left\{M_{N}^{*}>t\right\}\right) \mathrm{d} t \\
& \leqslant \int_{0}^{+\infty} p t^{p-2} \mathrm{E}\left(\mathbf{1}_{\left\{M_{N}^{*}>t\right\}}\left|M_{N}\right|\right) \mathrm{d} t=\mathrm{E}\left(\frac{p}{p-1}\left(M_{N}^{*}\right)^{p-1}\left|M_{N}\right|\right) \\
& \leqslant \frac{p}{p-1}\left(\mathrm{E}\left(\left(M_{N}^{*}\right)^{p}\right)\right)^{1-1 / p}\left(\mathrm{E}\left(\left|M_{N}\right|^{p}\right)\right)^{1 / p}
\end{aligned}
$$

hence $\left\|M_{N}^{*}\right\|_{p} \leqslant p(p-1)^{-1}\left\|M_{N}\right\|_{p}$. The case $p=+\infty$ is straightforward.
Remark 1.2. - In some contexts, it is useful to observe that the notion of conditional expectation on a sub- $\sigma$-field $\mathcal{F}_{0}$ of $\mathcal{F}$ remains well defined if we have a possibly infinite measure $\mu$ on $(\Omega, \mathcal{F})$, but which is $\sigma$-finite on $\mathcal{F}_{0}$, in other words, if $\Omega$ can be split in countably many sets $A_{i}$ in $\mathcal{F}_{0}$ such that $\mu\left(A_{i}\right)<+\infty$ for each $i$. If this condition is fulfilled by $\mu$ and by the smallest sub- $\sigma$-field $\mathcal{F}_{0}$ of a filtration $\left(\mathcal{F}_{k}\right)_{k=0}^{N}$, we can also speak about martingales with respect to the infinite measure $\mu$, and Theorem 1.1 remains true with the same proof, simply replacing the words "probability of an event" by "measure of a set".

We can always consider the orthogonal projection $\pi_{0}$ from $L^{2}(\Omega, \mathcal{F}, \mu)$ onto $L^{2}\left(\Omega, \mathcal{F}_{0}, \mu\right)$, but $L^{2}\left(\Omega, \mathcal{F}_{0}, \mu\right)=\{0\}$ when $\mathcal{F}_{0}$ does not contain any set with finite positive measure. On the other hand, when $A \in \mathcal{F}_{0}$ has finite measure, the formula $\pi_{0}\left(\mathbf{1}_{A} f\right)=\mathbf{1}_{A} \pi_{0}(f)$ allows one to work on $A$ as in the case of a probability measure.

### 1.2. The Hopf maximal inequality

We are given a measure space $(X, \Sigma, \mu)$ and a linear operator $T$ from $L^{1}(X, \Sigma, \mu)$ to itself. We shall only consider $\sigma$-finite measures throughout these Notes, and we work in this section with the space $L^{1}(X, \Sigma, \mu)$ of realvalued functions. We assume that $T$ is positive and nonexpansive, which
means that for every nonnegative function $g \in L^{1}(X, \Sigma, \mu), T g$ is nonnegative, and that the norm of $T$ is $\leqslant 1$. We can sum up these two properties by saying that when $g \geqslant 0$, then $T g \geqslant 0$ and $\int_{X} T g \mathrm{~d} \mu \leqslant \int_{X} g \mathrm{~d} \mu$.

Let us consider a function $f$ in $L^{1}(X, \Sigma, \mu)$, and for every integer $k \geqslant 0$ let

$$
S_{k}(f)=f+T f+T^{2} f+\cdots+T^{k} f
$$

If $N$ is a nonnegative integer, we set $S_{N}^{*}(f)=\max \left\{S_{j}(f): 0 \leqslant j \leqslant N\right\}$.
Lemma 1.3 (Hopf). - With the preceding notation, we have for every function $f \in L^{1}(X, \Sigma, \mu)$ and $N \geqslant 0$ that

$$
\int_{\left\{S_{N}^{*}(f)>0\right\}} f \mathrm{~d} \mu \geqslant 0
$$

Proof, after Garsia [38]. - Let us simply write $S_{k}$ for $S_{k}(f)$ and $S^{*}$ for $S_{N}^{*}(f)$. By definition, we have $S_{k} \leqslant S^{*}$ for each integer $k \leqslant N$; since $T$ is positive and linear, we see that

$$
T S_{k} \leqslant T S^{*}, \quad \text { and } \quad S_{k+1}=f+T S_{k} \leqslant f+T S^{*}
$$

In order to get for $S_{0}=f$ an inequality similar to $S_{k+1} \leqslant f+T S^{*}$, we replace $S^{*}$ by its nonnegative part $S^{*+}=\max \left(S^{*}, 0\right) \geqslant S^{*}$. Using positivity, we can write

$$
S_{0}=f \leqslant f+T\left(S^{*+}\right), \quad S_{k+1} \leqslant f+T S^{*} \leqslant f+T\left(S^{*+}\right)
$$

Taking the supremum of $S_{k}$ s for $0 \leqslant k \leqslant N$, we obtain the crucial inequality

$$
\begin{equation*}
S^{*} \leqslant f+T\left(S^{*+}\right), \text { or } f \geqslant S^{*}-T\left(S^{*+}\right) \tag{1.3}
\end{equation*}
$$

Since $T$ is positive and nonexpansive on $L^{1}(X, \Sigma, \mu)$, we have

$$
\int_{\left\{S^{*}>0\right\}} S^{*} \mathrm{~d} \mu=\int_{X} S^{*+} \mathrm{d} \mu \geqslant \int_{X} T\left(S^{*+}\right) \mathrm{d} \mu \geqslant \int_{\left\{S^{*}>0\right\}} T\left(S^{*+}\right) \mathrm{d} \mu
$$

and the result follows by (1.3), because

$$
\int_{\left\{S^{*}>0\right\}} f \mathrm{~d} \mu \geqslant \int_{\left\{S^{*}>0\right\}}\left(S^{*}-T\left(S^{*+}\right)\right) \mathrm{d} \mu \geqslant 0
$$

We go on with the same linear operator $T$. For each integer $k \geqslant 0$, let us define the $k$ th average operator $a_{k}=a_{k, T}$ associated to $T$ by writing

$$
a_{k}(f)=\frac{f+T f+\cdots+T^{k} f}{k+1}=\frac{S_{k}(f)}{k+1}, \quad f \in L^{1}(X, \Sigma, \mu)
$$

For each integer $N \geqslant 0$, let $a_{N}^{*}(f)=\max \left\{a_{j}(f): 0 \leqslant j \leqslant N\right\}$. It is clear that the set $\left\{a_{N}^{*}(f)>0\right\}$ coincides with the set $\left\{S_{N}^{*}(f)>0\right\}$ which appears in Lemma 1.3.

We continue in a simplified setting where we also assume that $\mu$ is finite and that $T \mathbf{1}=\mathbf{1}$. It follows that $a_{k}(\mathbf{1})=\mathbf{1}$ for each $k \geqslant 0$ and $a_{k}(f-c)=$ $a_{k}(f)-c$ for every $c \in \mathbb{R}$, thus $a_{N}^{*}(f-c)=a_{N}^{*}(f)-c$. Lemma 1.3 yields

$$
\int_{\left\{a_{N}^{*}(f-c)>0\right\}}(f-c) \mathrm{d} \mu=\int_{\left\{S_{N}^{*}(f-c)>0\right\}}(f-c) \mathrm{d} \mu \geqslant 0 .
$$

Equivalently, for every $f \in L^{1}(X, \Sigma, \mu)$, we have

$$
\begin{equation*}
c \mu\left(\left\{a_{N}^{*}(f)>c\right\}\right) \leqslant \int_{\left\{a_{N}^{*}(f)>c\right\}} f \mathrm{~d} \mu, \quad N \geqslant 0, c \in \mathbb{R} \tag{1.4}
\end{equation*}
$$

This inequality makes sense also when $\mu$ is infinite. Note that if $c<0$ and if $\mu$ is infinite, then $\mu(\{f \leqslant c\}) \leqslant \mu(\{|f| \geqslant|c|\})<+\infty$, the measure of $\left\{a_{N}^{*}(f)>c\right\}$ is thus infinite and (1.4) is trivial. We can extend (1.4) to an infinite $\mu$ if there exists an increasing sequence $\left(C_{\ell}\right)_{\ell \geqslant 0}$ of subsets of $X$ with finite measure such that

$$
\begin{align*}
& T^{j} \mathbf{1}_{C_{\ell}} \leqslant \mathbf{1} \text { for all } j, \ell \geqslant 0, \\
& T^{j} \mathbf{1}_{C_{\ell}} \underset{\ell \rightarrow+\infty}{\longrightarrow} \mathbf{1} \text { pointwise for each } j \geqslant 0 . \tag{1.5}
\end{align*}
$$

Let $c, \varepsilon>0$ and abbreviate $\left\{a_{N}^{*}(f)>t\right\}$ as $D(t)$, for $t>0$. Choose $c^{\prime}>c$ such that $\int_{D(c) \backslash D\left(c^{\prime}\right)}(1+|f|) \mathrm{d} \mu<\varepsilon$. Let $E\left(c^{\prime}, \ell\right)=\left\{\min _{0 \leqslant j \leqslant N} T^{j} \mathbf{1}_{C_{\ell}} \leqslant\right.$ $\left.c / c^{\prime}\right\}$, choose a large $\ell$ such that $\mu\left(D\left(c^{\prime}\right) \backslash C_{\ell}\right)<\varepsilon$ and $\int_{E\left(c^{\prime}, \ell\right)}|f| \mathrm{d} \mu<\varepsilon$, then observe that

$$
D\left(c^{\prime}\right) \subset\left\{a_{N}^{*}\left(f-c^{\prime} \mathbf{1}_{C_{\ell}}\right)>0\right\} \subset D(c) \cup E\left(c^{\prime}, \ell\right)
$$

and apply Lemma 1.3 to $f-c^{\prime} \mathbf{1}_{C_{\ell}}$. The assumption (1.5) is fulfilled when $T$ is an operator of convolution with a probability measure on $\mathbb{R}^{n}$, acting on $L^{1}\left(\mathbb{R}^{n}\right)$.

For each function $f \in L^{1}(X, \Sigma, \mu)$, let us define

$$
a^{*}(f)=\sup _{k \geqslant 0} a_{k}(f)=\sup _{k \geqslant 0} \frac{f+T f+\cdots+T^{k} f}{k+1}=\lim _{N \rightarrow+\infty} a_{N}^{*}(f) .
$$

The set $\left\{a^{*}(f)>c\right\}$ is the increasing union of the sets $\left\{a_{N}^{*}(f)>c\right\}, N \geqslant 0$, so, passing to the limit by dominated convergence, we deduce from (1.4) that

$$
\begin{equation*}
c \mu\left(\left\{a^{*}(f)>c\right\}\right) \leqslant \int_{\left\{a^{*}(f)>c\right\}} f \mathrm{~d} \mu, \quad c \in \mathbb{R} \tag{1.6}
\end{equation*}
$$

Following [29, Lemma VIII.6.7], we now get a variant of (1.6). Assume $c>0$ in what follows. We define $f_{c}$ by $f_{c}(x)=f(x)$ when $f(x)>c$ and $f_{c}(x)=0$ otherwise, for $x \in X$. Note that $f \leqslant f_{c}+c$. If $a^{*}\left(f_{c}\right)(x) \leqslant c$, then $f_{c}(x)=$ $a_{0}\left(f_{c}\right)(x) \leqslant c$ thus $f_{c}(x)=0$ by construction. Hence $f_{c}$ vanishes outside
$\left\{a^{*}\left(f_{c}\right)>c\right\}$ and

$$
\int_{\{f>c\}} f \mathrm{~d} \mu=\int_{X} f_{c} \mathrm{~d} \mu=\int_{\left\{a^{*}\left(f_{c}\right)>c\right\}} f_{c} \mathrm{~d} \mu
$$

Using the positivity of $T$ and of $a_{k}$ for each $k \geqslant 0$, we infer from $f_{c} \geqslant f-c$ that $a^{*}\left(f_{c}\right) \geqslant a^{*}(f-c)=a^{*}(f)-c$. Then, by (1.6) for $f_{c}$ and since $c>0$, we get

$$
\int_{\{f>c\}} f \mathrm{~d} \mu \geqslant c \mu\left(\left\{a^{*}\left(f_{c}\right)>c\right\}\right) \geqslant c \mu\left(\left\{a^{*}(f)-c>c\right\}\right) .
$$

Finally, we have obtained

$$
\begin{equation*}
c \mu\left(\left\{a^{*}(f)>2 c\right\}\right) \leqslant \int_{\{f>c\}} f \mathrm{~d} \mu, \quad c>0 . \tag{1.7}
\end{equation*}
$$

Let us define $A^{*}(f)=\sup _{k \geqslant 0}\left|a_{k}(f)\right|=\max \left(a^{*}(f), a^{*}(-f)\right)$. Still assuming $c>0$, we decompose the set $\left\{A^{*}(f)>c\right\}=\left\{a^{*}(f)>c\right\} \cup$ $\left\{a^{*}(-f)>c\right\}$ into three disjoint pieces, $E_{0}=\left\{a^{*}(f)>c, a^{*}(-f) \leqslant c\right\}$, $E_{1}=\left\{a^{*}(f)>c, a^{*}(-f)>c\right\}$, and $E_{2}=\left\{a^{*}(f) \leqslant c, a^{*}(-f)>c\right\}$. According to (1.6) we have

$$
\begin{align*}
c \mu\left(\left\{A^{*}(f)>c\right\}\right) & \leqslant c \mu\left(\left\{a^{*}(f)>c\right\}\right)+c \mu\left(\left\{a^{*}(-f)>c\right\}\right) \\
& \leqslant \int_{\left\{a^{*}(f)>c\right\}} f \mathrm{~d} \mu+\int_{\left\{a^{*}(-f)>c\right\}}(-f) \mathrm{d} \mu \\
& =\int_{E_{0}} f \mathrm{~d} \mu+\int_{E_{2}}(-f) \mathrm{d} \mu \leqslant \int_{\left\{A^{*}(f)>c\right\}}|f| \mathrm{d} \mu, \tag{1.8}
\end{align*}
$$

noting that the integrals of $f$ and $-f$ on $E_{1}$ cancel each other. In the same way, we can get from (1.7) the variant form $c \mu\left(\left\{A^{*} f>2 c\right\}\right) \leqslant$ $\int_{\{|f|>c\}}|f| \mathrm{d} \mu$. Notice that the latter "variant form" will be inherited by any linear operator $S$ satisfying that $\left|S^{k} f\right| \leqslant T^{k}|f|$ for every $k \geqslant 0$, and see Remark 1.5.

When $1<p<+\infty$, we deduce from (1.8) the $L^{p}$ inequality

$$
\begin{equation*}
\left\|\sup _{k \geqslant 0} \frac{\left|f+T f+\cdots+T^{k} f\right|}{k+1}\right\|_{p} \leqslant \frac{p}{p-1}\|f\|_{p} \tag{1.9}
\end{equation*}
$$

as we have seen with Doob's inequality (1.1), while the variant form leads to a constant $2(p /(p-1))^{1 / p}$ which is larger than $p /(p-1)$ for every $p>1$.

Let now $\left(T_{t}\right)_{t \geqslant 0}$ be a semi-group of linear operators on $L^{1}(X, \Sigma, \mu)$, i.e., operators satisfying $T_{s+t}=T_{s} \circ T_{t}$ for all $s, t \geqslant 0$. We assume in addition that each $T_{t}$ is positive and nonexpansive on $L^{1}$. We also assume that $T_{t}$ is actually defined on $L^{1}(X, \Sigma, \mu)+L^{\infty}(X, \Sigma, \mu)$ and that $T_{t} \mathbf{1}=\mathbf{1}$ for every
$t \geqslant 0$. This implies that $T_{t}$ is continuous from $L^{\infty}$ to $L^{\infty}$, with norm 1 . By interpolation, we get that the norm $\left\|T_{t}\right\|_{p \rightarrow p}$ on $L^{p}$, for $p \in[1,+\infty]$, is $\leqslant 1$. Suppose that the semi-group is strongly continuous on $L^{1}$, which means that $\left\|f-T_{t} f\right\|_{1} \rightarrow 0$ as $t \rightarrow 0$, for each $f \in L^{1}$. Combined with our assumptions, this fact implies that $t \mapsto T_{t} f$ is continuous from $[0,+\infty)$ to $L^{p}$ for every function $f \in L^{p}$ and $1 \leqslant p<+\infty$. For $f \in L^{p}(X, \Sigma, \mu)$ let

$$
a^{*} f=\sup _{t>0} \frac{1}{t} \int_{0}^{t} T_{s} f \mathrm{~d} s, \quad A^{*} f=\sup _{t>0}\left|\frac{1}{t} \int_{0}^{t} T_{s} f \mathrm{~d} s\right|,
$$

where the supremum can be defined as an essential supremum, see the discussion in Section 3.3. Yet, for the main examples of semi-groups of interest in these Notes, namely, the Gaussian semi-group or the Poisson semi-group on $\mathbb{R}^{n}$, the function $t \mapsto\left(T_{t} f\right)(x)$ is continuous on $(0,+\infty)$ for each fixed $x \in \mathbb{R}^{n}$ and $f \in L^{1}\left(\mathbb{R}^{n}\right)$, so $a^{*} f$ and $A^{*} f$ have then a well defined pointwise value, possibly $+\infty$.

Suppose now that the measure $\mu$ is finite (or that a continuous analog of (1.5) is satisfied). When $1<p<+\infty$, we obtain from (1.9) the $L^{p}$ inequality

$$
\begin{equation*}
\left\|A^{*} f\right\|_{p} \leqslant \frac{p}{p-1}\|f\|_{p} \tag{1.10}
\end{equation*}
$$

If $T_{t}$ is positive and $T_{t} \mathbf{1}=\mathbf{1}$, the case $p=+\infty$ in (1.10) is clear.
Since $t \mapsto a(t, f):=t^{-1} \int_{0}^{t} T_{s} f \mathrm{~d} s$ is continuous from $(0,+\infty)$ to $L^{p}$, we can reach any $a(t, f), t>0$, as an almost everywhere limit of a sequence $\left(a\left(t_{j}, f\right)\right)_{j \geqslant 0}$, where each $t_{j}$ is rational and $>0$. It follows that $A^{*} f$ can be defined as the supremum of $|a(t, f)|$ for $t>0$ rational. For all integers $k \geqslant 0$ and $n \geqslant 1$, observe that

$$
\begin{aligned}
a\left(\frac{k+1}{n}, f\right) & =\frac{n}{k+1} \sum_{i=0}^{k} \int_{i / n}^{(i+1) / n} T_{s} f \mathrm{~d} s \\
& =\left(\frac{\sum_{i=0}^{k} T_{i / n}}{k+1}\right)\left(n \int_{0}^{1 / n} T_{s} f \mathrm{~d} s\right) .
\end{aligned}
$$

Letting $f_{n}=n \int_{0}^{1 / n} T_{s} f \mathrm{~d} s=a(1 / n, f)$ and $T=T_{1 / n}$ we see that

$$
a\left(\frac{k+1}{n}, f\right)=\frac{f_{n}+T f_{n}+\cdots+T^{k} f_{n}}{k+1}=a_{k, T}\left(f_{n}\right) .
$$

Let $Q_{n}$ be the set of positive multiples of $1 / n$. By (1.9) applied to $T_{1 / n}$ and $f_{n}$, and because $a(1 / n, \cdot)$ is an average of operators with norm $\leqslant 1$ on $L^{p}$, we get

$$
\left\|\sup _{t \in Q_{n}}|a(t, f)|\right\|_{p}=\left\|\sup _{j \geqslant 1}|a(j / n, f)|\right\|_{p} \leqslant \frac{p}{p-1}\|a(1 / n, f)\|_{p} \leqslant \frac{p}{p-1}\|f\|_{p} .
$$

We see that $Q_{m} \subset Q_{m n}$ for all $m, n \geqslant 1$. The sets $Q_{n}$ corresponding to $n=\ell$ ! for $\ell \geqslant 1$ are increasing with $\ell$, and they cover the set of positive rationals. We can conclude by noticing that $A^{*} f$ is the increasing limit of $\sup _{t \in Q_{\ell!}}|a(t, f)|$.

Applying (1.6) we can obtain a version of Hopf's maximal inequality as

$$
c \mu\left(\left\{a^{*} f>c\right\}\right) \leqslant \int_{\left\{a^{*} f>c\right\}} f \mathrm{~d} \mu, \quad c \in \mathbb{R}, \quad f \in L^{1}(X, \Sigma, \mu)
$$

and from (1.8), we have $c \mu\left(\left\{A^{*} f>c\right\}\right) \leqslant \int_{\left\{A^{*} f>c\right\}}|f| \mathrm{d} \mu$ when $c>0$.
By the preceding remark about the sets $Q_{\ell!}$ it is enough to prove the inequality with $a_{n}^{*}:=\sup _{t \in Q_{n}} a(t, f)=\sup _{k \geqslant 0} a_{k, T_{1 / n}}\left(f_{n}\right)$ replacing $a^{*} f$, with $n \geqslant 1$ arbitrary and with a vanishing error term. By (1.6) we have $c \mu\left(\left\{a_{n}^{*}>c\right\}\right) \leqslant \int_{\left\{a_{n}^{*}>c\right\}} f_{n}$. Since the semi-group $\left(T_{t}\right)_{t \geqslant 0}$ is strongly continuous, we know that $\left\|f_{n}-f\right\|_{1} \rightarrow 0$ and we can conclude because $\int_{\left\{a_{n}^{*} f>c\right\}} f_{n} \mathrm{~d} \mu-\int_{\left\{a_{n}^{*} f>c\right\}} f \mathrm{~d} \mu$ tends to zero.
We have made here assumptions more restrictive than those of the Hopf-Dunford-Schwartz statement [29, Chap. VIII] praised by Stein [73], which does not assume $T_{t}$ positive, nor $\mu$ finite and $T_{t} \mathbf{1}=\mathbf{1}$. Theorem 1.4 below contains Lemma VIII.7. 6 and Theorem VIII.7.7 from [29] in a slightly simplified form (the set $U$ there has only one element here). The semi-group $\left(T_{t}\right)_{t \geqslant 0}$ on $L^{1}(X, \Sigma, \mu)$ is said to be strongly measurable if, for each $f$ in $L^{1}(X, \Sigma, \mu)$, the mapping $t \mapsto T_{t} f \in L^{1}(X, \Sigma, \mu)$ is measurable with respect to the Lebesgue measure on $[0,+\infty)$.

ThEOREM 1.4 ([29]). - Let $\left(T_{t}\right)_{t \geqslant 0}$ be a strongly measurable semi-group on the space $L^{1}(X, \Sigma, \mu)$, with $\left\|T_{t}\right\|_{1 \rightarrow 1} \leqslant 1$ and $\left\|T_{t}\right\|_{\infty \rightarrow \infty} \leqslant 1$ for all $t \geqslant 0$. For every function $f \in L^{1}(X, \Sigma, \mu)$ and every $c>0$ one has

$$
c \mu\left(\left\{A^{*} f>2 c\right\}\right) \leqslant \int_{\{|f|>c\}}|f| \mathrm{d} \mu
$$

If $1<p<+\infty$ and $f \in L^{p}(X, \Sigma, \mu)$, the function $A^{*} f$ is in $L^{p}(X, \Sigma, \mu)$ and

$$
\left\|A^{*} f\right\|_{p} \leqslant 2\left(\frac{p}{p-1}\right)^{1 / p}\|f\|_{p}
$$

Remark 1.5. - In [29, Section VIII.6], the authors consider first a linear operator $T$ acting from $L^{1}$ to $L^{1}$ with norm $\leqslant 1$ and also acting from $L^{\infty}$ to $L^{\infty}$ with norm $\leqslant 1$; in this discrete parameter case, they study

$$
A_{T}^{*} f=\sup _{n \geqslant 1} \frac{1}{n}\left|\sum_{k=0}^{n-1} T^{k} f\right|
$$

before going to the continuous setting of a semi-group $\left(T_{t}\right)_{t \geqslant 0}$. One of the steps in their proof consists in introducing a positive operator $P$ which acts
from $L^{1}$ to $L^{1}$ and from $L^{\infty}$ to $L^{\infty}$, with norm $\leqslant 1$ in both cases, and such that

$$
\forall n \geqslant 0, \quad\left|T^{n} f\right| \leqslant P^{n}(|f|), \quad f \in L^{\infty} \cap L^{1}
$$

This step is easy when the measure is the uniform measure on a finite set. The assumptions imply that $T$ is given by a matrix $\left(t_{i, j}\right)$ such that the sum of absolute values in each row and in each column is $\leqslant 1$. It is then enough to take $P$ to be the matrix with entries $p_{i, j}$ equal to the absolute values $\left|t_{i, j}\right|$ of the entries of $T$.

### 1.3. From martingales to semi-groups, via an argument of Rota

The arguments in this section, due to Rota [67], are presented in a more sophisticated manner in Stein's book [73, Chap. 4, §4]. We consider a Markov chain $X_{0}, \ldots, X_{N}$ with transition matrix $P$, assumed to be symmetric. We suppose for simplicity that the state space $\mathcal{E}$ is finite, with cardinality $Z$. For every $e_{0} \in \mathcal{E}$, we have

$$
\sum_{e \in \mathcal{E}} P\left(e_{0}, e\right)=1
$$

For each integer $k$ such that $0 \leqslant k<N$ and for all $e_{0}, e_{1} \in \mathcal{E}$, the probability that $X_{k+1}=e_{1}$ knowing that $X_{k}=e_{0}$ is given by the entry $P\left(e_{0}, e_{1}\right)$ of the matrix $P$. This statement introduces implicitly the Markov property, which loosely speaking, prescribes that what happens after time $k$ depends only on what we know at the instant $k$, regardless of the past positions at times $j<k$. For each integer $j \geqslant 2$, the power $P^{j}$ of the matrix $P$ controls the moves in $j$ successive steps, the entry $P^{j}\left(e_{0}, e\right)$ giving the probability of moving from $e_{0}$ to $e$ in exactly $j$ steps. If $Q$ is a transition matrix and $f$ a scalar function on $\mathcal{E}$, we introduce the notation

$$
(Q f)(x)=\sum_{y \in \mathcal{E}} Q(x, y) f(y), \quad x \in \mathcal{E}
$$

When applied to a power $P^{j}$, the notation $P^{j} f$ corresponds to the semigroup notation $P_{t} f$, with $j \in \mathbb{N}$ replacing $t \geqslant 0$. If the transition matrix $Q$ is symmetric, hence bistochastic, and if $1 \leqslant p \leqslant+\infty$, convexity implies that $\|Q f\|_{p} \leqslant\|f\|_{p}$ with respect to the uniform measure on $\mathcal{E}$. Let $f$ be a function on $\mathcal{E}$ and let $j, k$ be two nonnegative integers with $j+k \leqslant N$. If we fix $x_{0} \in \mathcal{E}$, the mean of the values $f(y)$, when the chain makes $j$ steps from the position $x_{0}$ at time $k$ to the position $y$ at time $k+j$, is equal to $\left(P^{j} f\right)\left(x_{0}\right)$.

A simple but important symmetric example is that of the Bernoulli random walk on $\mathbb{Z}$, where for all $x, y \in \mathbb{Z}$ we have $P(x, y)=1 / 2$ when $|x-y|=1$, and $P(x, y)=0$ otherwise. This is not a finite example, but it can be "approximated" by considering on the finite set $\mathcal{E}_{N}=\{-N, \ldots, N\}$,
for $N$ large, the modified matrix $P_{N}$ which still has $P_{N}(x, y)=1 / 2$ when $|x-y|=1$, for $x, y \in \mathcal{E}_{N}$, but where $P_{N}(N, N)=P_{N}(-N,-N)=1 / 2$. One can also consider the Bernoulli random walk on $\mathbb{Z}^{n}$, for which $P(x, y)=2^{-n}$ when $\left|x_{i}-y_{i}\right|=1$ for all coordinates $x_{i}, y_{i}, i=1, \ldots, n$, of the points $x=\left(x_{1}, \ldots, x_{n}\right)$ and $y=\left(y_{1}, \ldots, y_{n}\right)$ in $\mathbb{Z}^{n}$.

Assume that the distribution of the initial position $X_{0}$ is uniform, that is to say, that $P\left(X_{0}=e_{0}\right)=1 / Z$ for every $e_{0} \in \mathcal{E}$. Then for each $e_{1} \in \mathcal{E}$, we have

$$
\begin{aligned}
P\left(X_{1}=e_{1}\right) & =\sum_{e \in \mathcal{E}} P\left(X_{0}=e, X_{1}=e_{1}\right) \\
& =\sum_{e \in \mathcal{E}} \frac{1}{Z} P\left(e, e_{1}\right)=\frac{1}{Z} \sum_{e \in \mathcal{E}} P\left(e_{1}, e\right)=\frac{1}{Z},
\end{aligned}
$$

since the matrix $P$ is symmetric. The distribution of the position $X_{1}$ of the chain at time $i=1$ remains the uniform distribution, as well as that of $X_{2}, \ldots, X_{N}$. The uniform distribution is invariant under the action of $P$. Recalling the meaning of the transition matrix in terms of conditional probability, using Markov's property and letting $A_{N-1}=\left\{X_{0}=e_{0}, X_{1}=\right.$ $\left.e_{1}, \ldots, X_{N-1}=e_{N-1}\right\}$, we have that

$$
\begin{aligned}
E & :=P\left(X_{0}=e_{0}, X_{1}=e_{1}, \ldots, X_{N}=e_{N}\right) \\
& =P\left(A_{N-1}, X_{N}=e_{N}\right)=P\left(A_{N-1}\right) P\left(X_{N}=e_{N} \mid A_{N-1}\right) \\
& =P\left(A_{N-1}\right) P\left(X_{N}=e_{N} \mid X_{N-1}=e_{N-1}\right)=P\left(A_{N-1}\right) P\left(e_{N-1}, e_{N}\right)
\end{aligned}
$$

We may go on, and by the symmetry property of the matrix we get

$$
\begin{aligned}
E & =\cdots=\frac{1}{Z} P\left(e_{0}, e_{1}\right) P\left(e_{1}, e_{2}\right) \ldots P\left(e_{N-2}, e_{N-1}\right) P\left(e_{N-1}, e_{N}\right) \\
& =\frac{1}{Z} P\left(e_{N}, e_{N-1}\right) P\left(e_{N-1}, e_{N-2}\right) \ldots P\left(e_{2}, e_{1}\right) P\left(e_{1}, e_{0}\right) \\
& =P\left(X_{N}=e_{0}, X_{N-1}=e_{1}, \ldots, X_{1}=e_{N-1}, X_{0}=e_{N}\right)
\end{aligned}
$$

We see that the "reversed" chain has the same behavior as that of the original chain. Since the matrix is symmetric, we certainly have, whatever the distribution of $X_{0}$ can be, that the probability to arrive at a fixed $y_{0}$ at time $N$, starting from an arbitrary point $x$ at time $k=N-j$, is given by $P^{j}\left(x, y_{0}\right)=P^{j}\left(y_{0}, x\right)$, the probability of moving from $y_{0}$ at time 0 to $x$ at time $j$. But under the invariant distribution, we can say more: if $g$ is a function on $\mathcal{E}$, the mean of the values $g(x)$ on all trajectories starting from $x$ at time $k$ and arriving at $y_{0}$ at time $N$ is equal to $\left(P^{j} g\right)\left(y_{0}\right)$. Clearly, this statement is not true in general, since this mean value depends on the distribution of $X_{k}$, hence on that of $X_{0}$. Under the uniform distribution, we see by reversing the chain that the preceding mean is equal to the mean of
$g(x)$, when starting from $y_{0}$ at time 0 and arriving at $x$ at time $j$, namely, this mean is equal to $\left(P^{j} g\right)\left(y_{0}\right)$.

Let us describe the situation more formally. Let $\Omega=\mathcal{E}^{N+1}$ denote the space of all possible trajectories $\left(e_{0}, e_{1}, \ldots, e_{N}\right) \in \mathcal{E}^{N+1}$ for the chain. On this model space $\Omega$ and for $k=0, \ldots, N$, we set

$$
X_{k}(\omega)=\omega_{k} \in \mathcal{E}, \quad \omega=\left(\omega_{0}, \ldots, \omega_{N}\right) \in \mathcal{E}^{N+1}
$$

It is easy to determine the probability measure $\mathbf{P}$ on $\Omega$ that corresponds to the behavior of our Markov chain under the invariant distribution. For each singleton $\{\omega\}=\left\{\left(\omega_{0}, \ldots, \omega_{N}\right)\right\}$ in $\mathcal{P}(\Omega)$, we must have that

$$
\mathbf{P}\left(\left\{\left(\omega_{0}, \ldots, \omega_{N}\right)\right\}\right)=\frac{1}{Z} P\left(\omega_{0}, \omega_{1}\right) P\left(\omega_{1}, \omega_{2}\right) \ldots P\left(\omega_{N-1}, \omega_{N}\right)
$$

For $k=0, \ldots, N$, let $\mathcal{F}_{k}$ denote the finite field of subsets of $\Omega$ whose atoms $A$ are of the following form: to any $e_{0}, \ldots, e_{k}$ fixed in $\mathcal{E}$ we associate $A_{\mathbf{e}} \in \mathcal{F}_{k}$ defined by

$$
A=A_{\mathbf{e}}=\left\{\omega=\left(\omega_{0}, \ldots, \omega_{N}\right): \omega_{j}=e_{j}, 0 \leqslant j \leqslant k\right\} \in \mathcal{F}_{k}, \quad \mathbf{e}=\left(e_{0}, \ldots, e_{k}\right)
$$

This $\mathcal{F}_{k}$ is the "field of past events" at time $k$, it increases with $k$. Let $\mathcal{G}_{k}$ denote the field of events occurring precisely at time $k$, whose atoms $B$ are of the form

$$
B=\left\{\omega=\left(\omega_{0}, \ldots, \omega_{N}\right): \omega_{k}=e_{k}\right\} \in \mathcal{G}_{k}
$$

Clearly, we have $\mathcal{G}_{k} \subset \mathcal{F}_{k}$. A function on $\Omega$ which is $\mathcal{G}_{k}$-measurable depends only on the coordinate $\omega_{k}$, and is thus of the form $g\left(X_{k}\right)$ with $g$ a function on $\mathcal{E}$. If $f$ is a function on $\mathcal{E}$, the Markov property yields

$$
\mathrm{E}\left(f\left(X_{N}\right) \mid \mathcal{F}_{k}\right)=\mathrm{E}\left(f\left(X_{N}\right) \mid \mathcal{G}_{k}\right)=g\left(X_{k}\right)
$$

where $g(x)=\left(P^{N-k} f\right)(x)$ for every $x \in \mathcal{E}$. The preliminary discussion shows that

$$
\begin{equation*}
\left(P^{N-k} f\right)\left(X_{k}\right)=\mathrm{E}\left(f\left(X_{N}\right) \mid \mathcal{F}_{k}\right), \quad\left(P^{N-k} g\right)\left(X_{N}\right)=\mathrm{E}\left(g\left(X_{k}\right) \mid \mathcal{G}_{N}\right) \tag{1.11}
\end{equation*}
$$

We introduce the "canonical" martingale associated to a function $f$ on $\mathcal{E}$, by letting

$$
\begin{equation*}
M_{i}=\left(P^{N-i} f\right)\left(X_{i}\right)=\mathrm{E}\left(f\left(X_{N}\right) \mid \mathcal{F}_{i}\right), \quad 0 \leqslant i \leqslant N \tag{1.12}
\end{equation*}
$$

We see that in (1.11), one occurrence of $P^{N-k}$ relates to the expectation at time $k<N$ of future positions $f\left(X_{N}\right)$, while the other is about expectation at time $N$ of past positions $g\left(X_{k}\right)$. Combining the two equalities in (1.11) in a "back and forth" move, by taking $g=P^{j} f$ and $j=N-k$, we conclude that

$$
\begin{equation*}
\left(P^{2 j} f\right)\left(X_{N}\right)=\mathrm{E}\left(M_{N-j} \mid \mathcal{G}_{N}\right) \tag{1.13}
\end{equation*}
$$

Since the conditional expectation operator on $\mathcal{G}_{N}$ is positive, we see that for every $j=N-k=0, \ldots, N$, we have the inequality

$$
\max _{0 \leqslant j \leqslant N}\left|\left(P^{2 j} f\right)\left(X_{N}\right)\right|=\max _{0 \leqslant j \leqslant N}\left|\mathrm{E}\left(M_{N-j} \mid \mathcal{G}_{N}\right)\right| \leqslant \mathrm{E}\left(\max _{0 \leqslant i \leqslant N}\left|M_{i}\right| \mid \mathcal{G}_{N}\right)
$$

It implies when $1<p \leqslant+\infty$, according to Doob's inequality (1.1) and to the non-expansivity on $L^{p}$ of conditional expectations, the chain of inequalities

$$
\begin{align*}
&\left\|\max _{0 \leqslant j \leqslant N}\left|\left(P^{2 j} f\right)\left(X_{N}\right)\right|\right\|_{p} \leqslant\left\|\mathrm{E}\left(\max _{0 \leqslant i \leqslant N}\left|M_{i}\right| \mid \mathcal{G}_{N}\right)\right\|_{p} \\
& \leqslant\left\|\max _{0 \leqslant i \leqslant N}\left|M_{i}\right|\right\|_{p} \leqslant \frac{p}{p-1}\left\|M_{N}\right\|_{p}=\frac{p}{p-1}\left\|f\left(X_{N}\right)\right\|_{p} \tag{1.14}
\end{align*}
$$

We could recover the odd indices $2 j+1$ by applying the latter inequality to $P f$ instead of $f$ and using $\|P f\|_{p} \leqslant\|f\|_{p}$, to the cost of an extra factor 2.

Estimating the maximal function of semi-groups is a central theme in [73]. The discrete case of (1.14) was obtained by Stein in the short article [71], independently of Rota [67], by methods preluding those of [73]. Theorem 1 in [71] applies to self-adjoint operators $P$ on $L^{2}(X, \Sigma, \mu)$ satisfying also $\|P\|_{1 \rightarrow 1} \leqslant 1$ and $\|P\|_{\infty \rightarrow \infty} \leqslant 1$.
One can play the same game with convex functions other than the supremum function on $\mathbb{R}^{N+1}$. For example, let us begin with the convexity inequality

$$
\left(\sum_{0 \leqslant i \leqslant N}\left|\mathrm{E}\left(f_{i} \mid \mathcal{G}\right)\right|^{2}\right)^{1 / 2} \leqslant \mathrm{E}\left(\left(\sum_{0 \leqslant i \leqslant N}\left|f_{i}\right|^{2}\right)^{1 / 2} \mid \mathcal{G}\right)
$$

and make use of the Burkholder-Gundy inequalities of Theorem 1.6, in order to obtain, when $0 \leqslant j_{0}<j_{1}<\ldots<j_{r} \leqslant N, 1<p<+\infty$, and with respect to the invariant measure $\mu$, the inequality

$$
\begin{equation*}
\left\|\left(\sum_{k=1}^{r}\left|\left(P^{2 j_{k}} f-P^{2 j_{k-1}} f\right)\right|^{2}\right)^{1 / 2}\right\|_{L^{p}(\mu)} \leqslant c_{p}\|f\|_{L^{p}(\mu)} \tag{1.15}
\end{equation*}
$$

Indeed, we have seen in (1.13) that $\left(P^{2 j_{k}} f\right)\left(X_{N}\right)$ is the projection on $\mathcal{G}_{N}$ of the member $M_{N-j_{k}}=\mathrm{E}\left(f\left(X_{N}\right) \mid \mathcal{F}_{N-j_{k}}\right)$ of the martingale $\left(M_{j}\right)_{j=0}^{N}$ in (1.12). Then $L_{i}=M_{N-j_{r-i}}, i=0, \ldots, r$ is another martingale, and

$$
\left(P^{2 j_{k-1}} f\right)\left(X_{N}\right)-\left(P^{2 j_{k}} f\right)\left(X_{N}\right)=\mathrm{E}\left(M_{N-j_{k-1}}-M_{N-j_{k}} \mid \mathcal{G}_{N}\right)
$$

appears as projection on $\mathcal{G}_{N}$ of the martingale difference $d_{r-k+1}=L_{r-k+1}-$ $L_{r-k}$ (see Section 1.4.2) when $1 \leqslant k \leqslant r$. This principle can be applied for bounding diverse convex functions of a semi-group, by considering them as projections of corresponding functions of a martingale, for which we may have an " $L^{p}$ inequality".

Let us come back to (1.14). Since the distribution of $X_{N}$ is uniform, we can restate (1.14) when $1<p \leqslant+\infty$ as

$$
\left(\frac{1}{Z} \sum_{x \in \mathcal{E}} \max _{0 \leqslant j \leqslant N}\left|\left(P^{2 j} f\right)(x)\right|^{p}\right)^{1 / p} \leqslant \frac{p}{p-1}\left(\frac{1}{Z} \sum_{x \in \mathcal{E}}|f(x)|^{p}\right)^{1 / p}
$$

or else, changing the normalization and letting $N$ tend to infinity, we obtain

$$
\begin{equation*}
\left(\sum_{x \in \mathcal{E}} \sup _{j \geqslant 0}\left|\left(P^{2 j} f\right)(x)\right|^{p}\right)^{1 / p} \leqslant \frac{p}{p-1}\left(\sum_{x \in \mathcal{E}}|f(x)|^{p}\right)^{1 / p} \tag{1.16}
\end{equation*}
$$

We can also write

$$
\left(\sum_{x \in \mathcal{E}} \sup _{j \geqslant 0}\left|\left(P^{j} f\right)(x)\right|^{p}\right)^{1 / p} \leqslant \frac{2 p}{p-1}\left(\sum_{x \in \mathcal{E}}|f(x)|^{p}\right)^{1 / p}
$$

If we want to deal with a countably infinite state space $\mathcal{E}$ such as $\mathcal{E}=\mathbb{Z}^{n}$, we may accept (as Stein [73] does) to work with an infinite invariant measure, uniform on $\mathcal{E}$, that gives measure 1 to each singleton $\{e\}, e \in \mathcal{E}$. We then obtain the same maximal inequality (1.16), applying Remark 1.2. If we do not accept an "infinite probability", we may, for example with the Bernoulli random walk, work with "boxes" finite but large enough: if $f$ is finitely supported in $\mathbb{Z}^{n}$ and if $N$ is fixed, we can find a finite box $B$ in $\mathcal{E}$, so big that $P^{j} f$ vanishes outside $B$ for every $j \leqslant 2 N$. Changing the Bernoulli transition matrix $P(x, y)$ at the boundary of $B$, in order to force the Markov chain to remain inside, we are back to the finite case.

### 1.4. Brownian motion, and more on martingales

### 1.4.1. Gaussian distributions and Brownian motion

Let $|x|$ denote here the Euclidean norm of a vector $x$ in $\mathbb{R}^{n}$. For every probability measure $\mu$ on $\mathbb{R}^{n}$ having a finite first order moment $\int_{\mathbb{R}^{n}}|x| \mathrm{d} \mu(x)$, one defines the barycenter of $\mu$ as

$$
\operatorname{bar} \mu=\int_{\mathbb{R}^{n}} x \mathrm{~d} \mu(x) \in \mathbb{R}^{n}
$$

To a probability measure $\mu$ on $\mathbb{R}^{n}$ with finite second order moment $\int_{\mathbb{R}^{n}}|x|^{2} \mathrm{~d} \mu(x)$, one associates the quadratic form

$$
Q_{\mu}: \xi \mapsto \int_{\mathbb{R}^{n}}((x-\operatorname{bar} \mu) \cdot \xi)^{2} \mathrm{~d} \mu(x), \quad \xi \in \mathbb{R}^{n}
$$

The matrix Q of $Q_{\mu}$ with respect to the canonical basis of $\mathbb{R}^{n}$ is the covariance matrix of $\mu$. The quadratic form $Q_{\mu}$ is positive definite when $\mu$ is
not supported on any affine hyperplane, for example when $\mu$ is the uniform probability measure on a bounded convex set $C$ with non empty interior, i.e., a convex body $C$. We say that $\mu$ is centered when bar $\mu=0$, and in this case the expression of $Q_{\mu}$ simplifies to $Q_{\mu}(\xi)=\int_{\mathbb{R}^{n}}(x \cdot \xi)^{2} \mathrm{~d} \mu(x)$ for every $\xi \in \mathbb{R}^{n}$.

When $f$ is a probability density on $\mathbb{R}$ with finite second order moment, the variance $\sigma^{2}$ of $f(x) \mathrm{d} x$ is defined by

$$
\sigma^{2}=\int_{\mathbb{R}}\left(x-\int_{\mathbb{R}} y f(y) \mathrm{d} y\right)^{2} f(x) \mathrm{d} x
$$

When $f$ is centered, one has that $\sigma^{2}=\int_{\mathbb{R}} x^{2} f(x) \mathrm{d} x$.
A Gaussian random variable with distribution $N\left(0, \mathrm{I}_{n}\right)$ takes values in $\mathbb{R}^{n}$, its distribution $\gamma_{n}$ is symmetric, thus centered, defined on $\mathbb{R}^{n}$ by

$$
\begin{equation*}
\mathrm{d} \gamma_{n}(x)=(2 \pi)^{-n / 2} \mathrm{e}^{-|x|^{2} / 2} \mathrm{~d} x \tag{1.17}
\end{equation*}
$$

and $\gamma_{n}$ admits the identity matrix $\mathrm{I}_{n}$ as covariance matrix. If $F$ is an $n$ dimensional Euclidean space, we denote by $\gamma_{F}$ the image of $\gamma_{n}$ under an (any) isometry from $\mathbb{R}^{n}$ onto $F$. If $X$ is a $N\left(0, \mathrm{I}_{n}\right)$ Gaussian random variable and $\sigma>0$, then the multiple $\sigma X$ admits the distribution $\mathrm{d} \gamma_{n, \sigma}(x)=$ $(2 \pi)^{-n / 2} \mathrm{e}^{-|x / \sigma|^{2} / 2} \mathrm{~d}(x / \sigma)$, called the $N\left(0, \sigma^{2} \mathrm{I}_{n}\right)$ distribution, with $\sigma^{2} \mathrm{I}_{n}$ as covariance matrix. One can consider that the Dirac probability measure $\delta_{0}$ at the origin of $\mathbb{R}^{n}$ corresponds to $N\left(0,0_{n}\right)$.

The (absolute) moments of the one-dimensional distribution $\gamma_{1}$ can be computed in terms of values of the Gamma function. For every $p>-1$, one has that

$$
\int_{\mathbb{R}}|x|^{p} \mathrm{~d} \gamma_{1}(x)=(2 \pi)^{-1 / 2} \int_{\mathbb{R}}|x|^{p} \mathrm{e}^{-x^{2} / 2} \mathrm{~d} x=2^{p / 2} \pi^{-1 / 2} \Gamma((p+1) / 2)
$$

As $p$ tends to $+\infty$, it follows from Stirling's formula that

$$
\begin{equation*}
g_{p}:=\left(\int_{\mathbb{R}}|x|^{p} \mathrm{~d} \gamma_{1}(x)\right)^{1 / p} \simeq \sqrt{p / \mathrm{e}} \tag{1.18}
\end{equation*}
$$

An $n$-dimensional Brownian motion $\left(B_{t}\right)_{t \geqslant 0}$ starting at $x_{0} \in \mathbb{R}^{n}$ is an $\mathbb{R}^{n}$ valued random process defined on some probability space $(\Omega, \mathcal{F}, P)$, such that $B_{0}=x_{0}$, such that $B_{t}-B_{s}$ is a Gaussian random variable with distribution $N\left(0,(t-s) \mathrm{I}_{n}\right)$ whenever $0 \leqslant s \leqslant t$, and with independent increments: for every integer $k \geqslant 1$, when $0 \leqslant t_{0}<\ldots<t_{k}$ are given, then

$$
B_{t_{0}}, B_{t_{1}}-B_{t_{0}}, B_{t_{2}}-B_{t_{1}}, \ldots, B_{t_{k}}-B_{t_{k-1}}
$$

are independent. The coordinates $\left(B_{t, i}\right)_{t \geqslant 0}, i=1, \ldots, n$, are independent one-dimensional Brownian motions. It is possible to choose everywhere defined measurable functions $\left(B_{t}\right)_{t \geqslant 0}$ satisfying the above properties in such a
way that the trajectories $0 \leqslant t \mapsto B_{t}(\omega)$, or random paths, are continuous for (almost) every $\omega \in \Omega$. The Brownian motion is a martingale with continuous time parameter $t \geqslant 0$, with respect to a continuous time filtration $\left(\mathcal{F}_{t}\right)_{t \geqslant 0}$ where $\mathcal{F}_{t}$ is generated by the variables $B_{s}, 0 \leqslant s \leqslant t$. See for example Durrett [31] for a detailed account.

It is well known that the Brownian motion on $\mathbb{R}^{n}$ is the limit of Markov chains with symmetric transition matrix, namely, a limit of suitably scaled Bernoulli random walks. Indeed, if $\delta>0$ is given and if we consider a Bernoulli walk on the real line moving at each time $k \delta, k \in \mathbb{N}^{*}$, by a step $\pm \sqrt{\delta}$, so that

$$
X_{t}^{(\delta)}=\sqrt{\delta} \sum_{k=1}^{\lfloor t / \delta\rfloor} \varepsilon_{k}, \quad t \geqslant 0, \quad \varepsilon_{k}= \pm 1
$$

then the distribution of $\left(X_{t}^{(\delta)}\right)_{t \geqslant 0}$ tends when $\delta \rightarrow 0$ to that of a onedimensional Brownian motion. Here, $\left(\varepsilon_{k}\right)_{k=1}^{\infty}$ is a sequence of independent Bernoulli random variables, taking values $\pm 1$ with probability $1 / 2$. If $\left(B_{t}\right)_{t \geqslant 0}$ is the Brownian motion in $\mathbb{R}^{n}$, starting at 0 , and if we consider the associated Gaussian semi-group $\left(G_{s}\right)_{s \geqslant 0}$ defined for $f \in L^{1}\left(\mathbb{R}^{n}\right)$ and $s>0$ by

$$
\begin{align*}
\left(G_{s} f\right)(x) & =\mathrm{E} f\left(x+B_{s}\right) \\
& =(2 \pi s)^{-n / 2} \int_{\mathbb{R}^{n}} f(x+y) \mathrm{e}^{-|y|^{2} /(2 s)} \mathrm{d} y, \quad x \in \mathbb{R}^{n}, \tag{1.19}
\end{align*}
$$

we can show an inequality analogous to (1.16). For every $p$ in $(1,+\infty]$ and for every function $f \in L^{p}\left(\mathbb{R}^{n}\right)$, we have a maximal inequality for the Gaussian semi-group with a bound independent of the dimension $n$, stating that

$$
\begin{equation*}
\left(\int_{\mathbb{R}^{n}} \sup _{s>0}\left|\left(G_{s} f\right)(x)\right|^{p} \mathrm{~d} x\right)^{1 / p} \leqslant \frac{p}{p-1}\left(\int_{\mathbb{R}^{n}}|f(x)|^{p} \mathrm{~d} x\right)^{1 / p} \tag{*}
\end{equation*}
$$

If we just need a maximal inequality possibly dimension dependent, there is an easy proof relating the Gaussian maximal function to the classical maximal function $\mathrm{M} f$, because the Gaussian kernel is radial and radially decreasing, see (4.6). Once Stein's Theorem 4.1 giving dimensionless estimates for $\mathrm{M} f$ is established, this easy bound of $G_{s} f$ by $\mathrm{M} f$ implies a dimensionless estimate for the Gaussian semi-group, or for the Poisson semi-group as well. With Bourgain, Carbery and Müller, we shall follow the opposite route, from the semi-group estimates to $\mathrm{M} f$ or $\mathrm{M}_{C} f$. We sketch an argument for obtaining (1.20. $G^{*}$ ) from the Bernoulli case.

Let us give some more details in dimension $n=1$. Let $\left(\varepsilon_{k}\right)_{k=1}^{\infty}$ be a sequence of independent Bernoulli random variables, taking values $\pm 1$ with probability $1 / 2$. The associated semi-group $\left(P_{j}\right)$, indexed by $j \in \mathbb{N}$, is
defined by

$$
\left(P_{j} g\right)(i)=\mathrm{E} g\left(i+\sum_{1 \leqslant k \leqslant j} \varepsilon_{k}\right), \quad j \geqslant 0, \quad i \in \mathbb{Z}
$$

and it satisfies (1.16). As a consequence of the de Moivre-Laplace theorem and by classical tail estimates, we know that

$$
\left(1+x^{2}\right) P\left(\left\{N^{-1 / 2} \sum_{1 \leqslant k \leqslant N} \varepsilon_{k}<x\right\}\right)
$$

tends to $\left(1+x^{2}\right) \gamma_{1}((-\infty, x))$ when $N \rightarrow \infty$, uniformly in $x$ real. It follows that $\operatorname{Ef}\left(N^{-1 / 2} \sum_{1 \leqslant k \leqslant N} \varepsilon_{k}\right)$ tends to $\int_{\mathbb{R}} f(y) d \gamma_{1}(y)$, uniformly on Lipschitz functions having a Lipschitz constant bounded by some fixed $C$. If $f$ is Lipschitz on $\mathbb{R}$, then

$$
\operatorname{E} f\left(x+N^{-1 / 2} \sum_{1 \leqslant k<s N} \varepsilon_{k}\right) \underset{N}{\longrightarrow}(2 \pi s)^{-1 / 2} \int_{\mathbb{R}} f(x+y) \mathrm{e}^{-y^{2} /(2 s)} \mathrm{d} y,
$$

uniformly in $x \in \mathbb{R}$ and $s \in\left[t_{0}, t_{1}\right]$, with $0<t_{0} \leqslant t_{1}$ fixed. This implies that for any given $\varepsilon>0$ and $N$ large enough, letting $g_{N}(i)=f(i / \sqrt{N})$ for $i \in \mathbb{Z}$ and assuming $s N-1 \leqslant j_{N}<s N$, we have that

$$
\left|P_{j_{N}} g_{N}(i)-\left(G_{s} f\right)(i / \sqrt{N})\right|<\varepsilon, \quad i \in \mathbb{Z}
$$

for every $s \in\left[t_{0}, t_{1}\right]$. Applying (1.16) to $g_{N}$, we obtain when $s_{0}, s_{1}, \ldots, s_{k}$ and $a>0$ are given that

$$
\int_{-a}^{a} \max _{0 \leqslant j \leqslant k}\left|\left(G_{s_{j}} f\right)(x)\right|^{p} \mathrm{~d} x \leqslant \eta^{p}(\varepsilon)+\left(\frac{p}{p-1}\right)^{p} \frac{1}{\sqrt{N}} \sum_{i \in \mathbb{Z}}\left|f\left(\frac{i}{\sqrt{N}}\right)\right|^{p}
$$

where $\eta(\varepsilon)$ tends to 0 with $\varepsilon$, implying (1.20.G*) when $\varepsilon \rightarrow 0, N \rightarrow+\infty$, $a \rightarrow+\infty$ and if the sequence $\left\{s_{j}\right\}_{j \geqslant 0}$ is dense in $(0,+\infty)$. The same argument works in $\mathbb{R}^{n}$, thanks to the product structure of the Bernoulli and Gaussian measures and to the fact that the linear space generated by products $f\left(x_{1}, \ldots, x_{n}\right)=\prod_{j=1}^{n} f_{j}\left(x_{j}\right)$ is uniformly dense in the space of compactly supported Lipschitz functions on $\mathbb{R}^{n}$.

These considerations generalize to semi-groups of convolution with symmetric probability measures $\left(\mu_{t}\right)_{t \geqslant 0}$ on $\mathbb{R}^{n}$, that is to say, when $\mu_{s} * \mu_{t}=$ $\mu_{s+t}, s, t \geqslant 0$, and $\mu_{t}(A)=\mu_{t}(-A)$ for every Borel subset $A \subset \mathbb{R}^{n}$. Given $k>1$, one can find a finitely supported symmetric probability measure $\nu_{1 / k}$ on $\mathbb{R}^{n}$ which is an approximation of $\mu_{1 / k}$, in the sense that the integrals of a given finite family of functions $f$ on $\mathbb{R}^{n}$ are nearly the same for $\mu_{j / k}$ and for $\nu_{1 / k}^{* j}$ whenever $j \leqslant k^{2}$. We may assume that $\nu_{1 / k}$ is supported in $\varepsilon \mathbb{Z}, \varepsilon>0$. The symmetric Markov chain $\left(X_{j}\right)_{j \leqslant k^{2}}$ on $\mathcal{E}=\varepsilon \mathbb{Z}$ with transition governed by $\nu_{1 / k}$ permits us to approximate the maximal function $\sup _{t}\left|\mu_{t} * f\right|$ of the semi-group, replacing it with $\max _{j \leqslant k^{2}}\left|\nu_{1 / k}^{* j} * f\right|$.

It follows that some convex functions of the convolution semi-group can be estimated in $L^{p}$ by projecting functions of a martingale. For example, the sum of squares of differences, already mentioned in the Gaussian case, can be studied also in the Poisson case by relating it to the square function of a martingale and applying the Burkholder-Gundy inequalities presented in the next section.

### 1.4.2. The Burkholder-Gundy inequalities

When $\left(M_{k}\right)_{k=0}^{N}$ is a martingale with respect to a filtration $\left(\mathcal{F}_{k}\right)_{k=0}^{N}$, one introduces the difference sequence $\left(d_{k}\right)_{k=0}^{N}$, which is defined by $d_{0}=M_{0}$ and $d_{k}=M_{k}-M_{k-1}$ if $0<k \leqslant N$. Observe that $d_{k}$ is $\mathcal{F}_{k}$-measurable for $0 \leqslant k \leqslant N$ and that $\mathrm{E}\left(d_{k} \mid \mathcal{F}_{k-1}\right)=0$ for $k>0$. Conversely, given a sequence $\left(d_{k}\right)_{k=0}^{N}$ with these two properties, we obtain a martingale by setting $M_{k}=\sum_{j=0}^{k} d_{j}$, for $0 \leqslant k \leqslant N$. For a scalar martingale $\left(M_{k}\right)_{k=0}^{N}$, we define the square function process $\left(S_{k}\right)_{k=0}^{N}$ of the martingale by

$$
S_{k}=\left(\sum_{j=0}^{k}\left|d_{j}\right|^{2}\right)^{1 / 2}, \quad k=0, \ldots, N
$$

For a real or complex martingale in $L^{2}$, the differences $d_{k}$ and $d_{\ell}$ are orthogonal when $k \neq \ell$. If $k<\ell$ for example, then $d_{k}$ and its complex conjugate $\overline{d_{k}}$ are $\mathcal{F}_{\ell-1}$-measurable, thus $\mathrm{E}\left(\overline{d_{k}} d_{\ell}\right)=\mathrm{E}\left(\overline{d_{k}} \mathrm{E}\left(d_{\ell} \mid \mathcal{F}_{\ell-1}\right)\right)=0$. It follows that

$$
\begin{equation*}
\mathrm{E}\left|M_{N}\right|^{2}=\sum_{k=0}^{N} \mathrm{E}\left|d_{k}\right|^{2}=\mathrm{E}\left|S_{N}\right|^{2} \tag{1.21}
\end{equation*}
$$

This equality $\left\|M_{N}\right\|_{2}=\left\|S_{N}\right\|_{2}$ appears as an evident case of the following result.

Theorem 1.6 (Burkholder-Gundy [17]). - For every $p$ in $(1,+\infty)$, there exists a constant $c_{p} \geqslant 1$ such that for every integer $N \geqslant 1$, for every real or complex martingale $\left(M_{k}\right)_{k=0}^{N}$, one has

$$
c_{p}^{-1}\left\|M_{N}\right\|_{p} \leqslant\left\|S_{N}\right\|_{p} \leqslant c_{p}\left\|M_{N}\right\|_{p}
$$

The Khinchin inequalities (see for example Zygmund [85, vol. I, V.8, Th. 8.4]) are a very particular instance of the preceding theorem. Let $\left(\varepsilon_{k}\right)_{k=1}^{N}$ be a sequence of independent Bernoulli random variables defined on a probability space $(\Omega, \mathcal{F}, P)$, taking the values $\pm 1$ with probability $1 / 2$. For every $p$ in $(0,+\infty)$, there exist constants $A_{p}, B_{p}>0$ such that for every $N \geqslant 1$
and all scalars $\left(a_{k}\right)_{k=0}^{N}$, one has

$$
\begin{gather*}
A_{p}\left(\sum_{k=1}^{N}\left|a_{k}\right|^{2}\right)^{1 / 2} \leqslant\left(\mathrm{E}\left|\sum_{k=1}^{N} a_{k} \varepsilon_{k}\right|^{p}\right)^{1 / p} \leqslant\left(\sum_{k=1}^{N}\left|a_{k}\right|^{2}\right)^{1 / 2}, 0<p \leqslant 2  \tag{1.22.K}\\
\left(\sum_{k=1}^{N}\left|a_{k}\right|^{2}\right)^{1 / 2} \leqslant\left(\mathrm{E}\left|\sum_{k=1}^{N} a_{k} \varepsilon_{k}\right|^{p}\right)^{1 / p} \leqslant B_{p}\left(\sum_{k=1}^{N}\left|a_{k}\right|^{2}\right)^{1 / 2}, 2 \leqslant p
\end{gather*}
$$

The exact values of the constants $A_{p}, B_{p}$ are known ([79, 43]). In order to relate these inequalities to Theorem 1.6 when $1<p<+\infty$, we consider a special filtration on $(\Omega, \mathcal{F}, P)$, generated by the sequence $\left(\varepsilon_{k}\right)_{k=1}^{N}$. Let $\mathcal{F}_{0}$ be the trivial field consisting of $\Omega$ and $\emptyset$, and for $k>0$, let $\mathcal{F}_{k}$ be the finite field generated by $\varepsilon_{1}, \ldots, \varepsilon_{k}$. This field $\mathcal{F}_{k}$ has $2^{k}$ atoms of the form

$$
\begin{equation*}
A=A_{\mathbf{u}}=\left\{\omega \in \Omega: \varepsilon_{j}(\omega)=u_{j}, j=1, \ldots, k\right\}, \quad \mathbf{u}=\left(u_{1}, \ldots, u_{k}\right) \tag{1.23}
\end{equation*}
$$

where $u_{j}= \pm 1$. We shall call this particular sequence $\left(\mathcal{F}_{k}\right)_{k=0}^{N}$ of finite fields a dyadic filtration. In this framework, for $1 \leqslant k \leqslant N$, any scalar multiple $a_{k} \varepsilon_{k}$ of $\varepsilon_{k}$ is a martingale difference $d_{k}$. For the associated martingale with $M_{N}=\sum_{k=1}^{N} a_{k} \varepsilon_{k}$, the square function $S_{N}$ is the constant function equal to $\left(\sum_{k=1}^{N}\left|a_{k}\right|^{2}\right)^{1 / 2}$ and the Khinchin inequalities appear indeed as a simple example of application of Theorem 1.6. Of course, the latter sentence is historically totally inaccurate.

We shall prove only special cases of Theorem 1.6. We say that a sequence of random variables $\left(m_{k}\right)_{k=0}^{N}$ is predictable when
$m_{0}$ is $\mathcal{F}_{0}$-measurable, and $m_{k}$ is $\mathcal{F}_{k-1}$-measurable for $0<k \leqslant N$.
If $\left(m_{k}\right)_{k=0}^{N}$ is scalar valued and predictable, and if $\left(d_{k}\right)_{k=0}^{N}$ is a martingale difference sequence, then $\left(m_{k} d_{k}\right)_{k=0}^{N}$ is again a martingale difference sequence since one has that $\mathrm{E}\left(m_{k} d_{k} \mid \mathcal{F}_{k-1}\right)=m_{k} \mathrm{E}\left(d_{k} \mid \mathcal{F}_{k-1}\right)=0$. The new martingale $\left(L_{k}\right)_{k=0}^{N}$ defined by $L_{k}=\sum_{j=0}^{k} m_{j} d_{j}$ is said to be obtained as martingale transform, see [15, 16].

Consider a dyadic filtration $\left(\mathcal{F}_{k}\right)_{k=0}^{N}$ as defined above. Notice that each atom $A$ of $\mathcal{F}_{k}$ as in (1.23) has probability $2^{-k}$, and is split into two atoms $A_{ \pm}$of $\mathcal{F}_{k+1}, A_{ \pm}:=A \cap\left\{\varepsilon_{k+1}= \pm 1\right\}$, according to the value of $\varepsilon_{k+1}$. Let $d_{k+1}$ be a martingale difference with respect to these dyadic fields. The function $d_{k+1}$ should have mean 0 on the atom $A$ of $\mathcal{F}_{k}$, and be constant on each of the two atoms $A_{ \pm}$of $\mathcal{F}_{k+1}$ contained in $A$, which have equal measure $P(A) / 2$. It follows that $d_{k+1}$ must take on $A$ two opposite values $\pm v$. Consequently, the modulus (or the norm) of $d_{k+1}$ is constant on $A$, thus $\left|d_{k+1}\right|$ is $\mathcal{F}_{k}$-measurable, so that $\left(\left|d_{k}\right|\right)_{k=0}^{N}$ is predictable, as defined in (1.24). We shall call Bernoulli martingale any martingale $\left(M_{k}\right)_{k=0}^{N}$ with
respect to this dyadic filtration $\left(\mathcal{F}_{k}\right)_{k=0}^{N}$. A Bernoulli martingale with values in a vector space can be pictured as a tree $\left(v_{\varepsilon_{1}, \ldots, \varepsilon_{k}}\right)$ of vectors, $0 \leqslant k \leqslant N$ and $\varepsilon_{j}= \pm 1$, such that each vector $v_{\varepsilon_{1}, \ldots, \varepsilon_{k}}$ in the tree is the midpoint of his two successors $v_{\varepsilon_{1}, \ldots, \varepsilon_{k}, 1}$ and $v_{\varepsilon_{1}, \ldots, \varepsilon_{k},-1}$. The vectors $v_{\varepsilon_{1}, \ldots, \varepsilon_{k}}$ are the values of the $k$ th random variable $M_{k}$ of the martingale, which can be defined by $M_{k}\left(\varepsilon_{1}, \ldots, \varepsilon_{k}\right)=v_{\varepsilon_{1}, \ldots, \varepsilon_{k}}$.

The next Lemma contains an easier case of a result due to Burgess Davis [26], namely, the left-hand inequality when $p=1$. The rest of the statement presents a mixture of Doob's and Burkholder-Gundy's inequalities.

Lemma 1.7. - For every $p$ with $1 \leqslant p \leqslant 2$ and for every real or complex Bernoulli martingale $\left(M_{k}\right)_{k=0}^{N}$, one has that

$$
6^{-1}\left\|M_{N}^{*}\right\|_{p} \leqslant\left\|S_{N}\right\|_{p} \leqslant 6\left\|M_{N}^{*}\right\|_{p}
$$

Partial proof, after [56]. - We consider the case $p=1$. The general strategy is to bring the problem to $L^{2}$, where $\left\|S_{N}\right\|_{2}=\left\|M_{N}\right\|_{2}$ by (1.21), and this is essentially done by dividing $f=M_{N} \in L^{1}$ by a "parent" of $\sqrt{|f|}$, in order to get an element in $L^{2}$ "similar" to $\sqrt{|f|}$. One then applies known facts in $L^{2}$, and finally come back to $L^{1}$ by multiplication with a suitable $L^{2}$ function. We begin with the proof of the left-hand inequality in Lemma 1.7.

Let $\left(M_{k}\right)_{k=0}^{N}$ be a Bernoulli martingale. We know that $\left(\left|d_{k}\right|\right)_{k=0}^{N}$ is predictable, as well as $\left(S_{k}\right)_{k=0}^{N}$. Consider the martingale transform $L_{k}=$ $\sum_{j=0}^{k} S_{j}^{-1 / 2} d_{j}$. In $L^{2}$ we know that $\mathrm{E}\left|L_{N}\right|^{2}=\sum_{j=0}^{N} \mathrm{E}\left(S_{j}^{-1}\left|d_{j}\right|^{2}\right)$. We see that $S_{0}^{-1}\left|d_{0}\right|^{2}=S_{0}$, and $S_{j}^{-1}\left|d_{j}\right|^{2} \leqslant 2\left(S_{j}-S_{j-1}\right)$ for $j \geqslant 1$ because, letting $t=S_{j-1}^{2}$ and $h=\left|d_{j}\right|^{2}$, we have

$$
2(\sqrt{t+h}-\sqrt{t})=\int_{t}^{t+h} u^{-1 / 2} \mathrm{~d} u \geqslant h(t+h)^{-1 / 2}
$$

It follows that

$$
\begin{equation*}
\mathrm{E}\left|L_{N}\right|^{2} \leqslant 2 \mathrm{E} S_{N} \tag{1.25}
\end{equation*}
$$

Notice that $\left|\sum_{j=0}^{s} S_{j}^{-1 / 2} d_{j}\right|=\left|L_{s}\right| \leqslant L_{N}^{*}$ and $\left|\sum_{j=r+1}^{s} S_{j}^{-1 / 2} d_{j}\right|=\mid L_{s}-$ $L_{r} \mid \leqslant 2 L_{N}^{*}$ when $0 \leqslant r<s \leqslant N$. Multiplying termwise the sequence $\left(S_{k}^{-1 / 2} d_{k}\right)_{k=0}^{N}$ by the non-decreasing sequence $\left(S_{k}^{1 / 2}\right)_{k=0}^{N}$, we obtain for every $s \leqslant N$ by Abel's summation method that

$$
\left|M_{s}\right|=\left|\sum_{j=0}^{s} d_{j}\right| \leqslant S_{s}^{1 / 2} \sup _{0 \leqslant r \leqslant s}\left|\sum_{j=r}^{s} S_{j}^{-1 / 2} d_{j}\right| \leqslant 2 S_{N}^{1 / 2} L_{N}^{*},
$$

thus $M_{N}^{*} \leqslant 2 S_{N}^{1 / 2} L_{N}^{*}$. By Cauchy-Schwarz, Doob's inequality (1.1) with $p=2$, and by (1.25), we get the conclusion

$$
\mathrm{E} M_{N}^{*} \leqslant 2\left(\mathrm{E} S_{N}\right)^{1 / 2}\left\|L_{N}^{*}\right\|_{2} \leqslant 2^{2}\left(\mathrm{E} S_{N}\right)^{1 / 2}\left\|L_{N}\right\|_{2} \leqslant 2^{5 / 2} \mathrm{E} S_{N} \leqslant 6 \mathrm{E} S_{N}
$$

We leave the rewriting of this proof when $1<p<2$ as an easy exercise for the reader, and we pass to the right-hand side inequality using the same method, with the help of the non-decreasing predictable sequence $\left(A_{k}\right)_{k=0}^{N}$ defined by

$$
A_{0}=\left|d_{0}\right|=\left|M_{0}\right|, \quad A_{k}=\max \left(A_{k-1}, M_{k-1}^{*}+\left|d_{k}\right|\right) \geqslant\left|M_{k}\right|, \quad k=1, \ldots, N
$$

and of the martingale transform $L_{k}=\sum_{j=0}^{k} A_{j}^{-1 / 2} d_{j}, k=0, \ldots, N$. Observe that $\left|d_{k}\right| \leqslant\left|M_{k}\right|+\left|M_{k-1}\right| \leqslant 2 M_{N}^{*}$, thus $A_{N} \leqslant 3 M_{N}^{*}$. By Abel, writing $d_{k}=M_{k}-M_{k-1}$ for $k \geqslant 1$, we see that

$$
\begin{aligned}
\left|L_{N}\right| & =\left|A_{N}^{-1 / 2} M_{N}+\sum_{k=0}^{N-1} M_{k}\left(A_{k}^{-1 / 2}-A_{k+1}^{-1 / 2}\right)\right| \\
& \leqslant A_{N}^{1 / 2}+\sum_{k=0}^{N-1} A_{k}\left(A_{k}^{-1 / 2}-A_{k+1}^{-1 / 2}\right) \\
& \leqslant A_{N}^{1 / 2}+\sum_{k=0}^{N-1}\left(\sqrt{A_{k+1}}-\sqrt{A_{k}}\right) \leqslant 2 A_{N}^{1 / 2}
\end{aligned}
$$

where we make use of $u^{2}\left(u^{-1}-v^{-1}\right) \leqslant v-u$ when $0<u \leqslant v$. In $L^{2}$ we know that $\mathrm{E}\left(\sum_{k=0}^{N} A_{k}^{-1}\left|d_{k}\right|^{2}\right)=\mathrm{E}\left|L_{N}\right|^{2} \leqslant 4 \mathrm{E} A_{N}$, and we go back to $L^{1}$ with Cauchy-Schwarz and the obvious inequality

$$
\sum_{k=0}^{N}\left|d_{k}\right|^{2} \leqslant A_{N} \sum_{k=0}^{N} A_{k}^{-1}\left|d_{k}\right|^{2}
$$

We obtain

$$
\mathrm{E} S_{N}=\mathrm{E}\left(\sum_{k=0}^{N}\left|d_{k}\right|^{2}\right)^{1 / 2} \leqslant\left(\mathrm{E} A_{N}\right)^{1 / 2}\left\|L_{N}\right\|_{2} \leqslant 2 \mathrm{E} A_{N} \leqslant 6\left\|M_{N}^{*}\right\|_{1}
$$

Remark. - The Brownian martingales can be approximated by Bernoulli martingales, and we can obtain the analogous result for them. Actually, the preceding proof is even simpler to write in this case. Brownian martingales are defined by means of (Itō's) stochastic integrals

$$
M_{t}(\omega)=\int_{0}^{t} m_{s}(\omega) \mathrm{d} B_{s}(\omega), \quad t \geqslant 0
$$

where $\left(m_{s}\right)_{s \geqslant 0}$ is an adapted process, meaning essentially that each $m_{s}$, $s \geqslant 0$, is $\mathcal{F}_{s}$-measurable. The square function is then defined by $S_{t}^{2}(\omega)=$
$\int_{0}^{t}\left|m_{s}(\omega)\right|^{2} \mathrm{~d} s$ for every $t \geqslant 0$, and one can replace in the proof of Lemma 1.7 the Abel summation method by the more pleasant integration by parts.

Remark 1.8. - Together with Doob's inequality, Lemma 1.7 implies Theorem 1.6 for Bernoulli martingales when $1<p \leqslant 2$. The BurkholderGundy inequalities are equivalent to saying that martingale difference sequences are unconditional in $L^{p}$ when $1<p<+\infty$, that is to say, that there exists a constant $\kappa_{u, p}$ such that for each integer $N \geqslant 0$, all scalars $\left(a_{k}\right)_{k=0}^{N}$ with $\left|a_{k}\right| \leqslant 1$ and all martingale differences $\left(d_{k}\right)_{k=0}^{N}$, we have

$$
\begin{equation*}
\left\|\sum_{k=0}^{N} a_{k} d_{k}\right\|_{p} \leqslant \kappa_{u, p}\left\|\sum_{k=0}^{N} d_{k}\right\|_{p} . \tag{1.26}
\end{equation*}
$$

Going from Theorem 1.6 to unconditionality is simple, since the square function of the martingale at the left-hand side of (1.26) is less than that on the right-hand side, and we can take $\kappa_{u, p}=c_{p}^{2}$. The other direction follows from Khinchin, by averaging over signs $a_{k}= \pm 1$. Indeed, one obtains from (1.22.K) for $\left(f_{k}\right)_{k=1}^{N}$ in $L^{p}(X, \Sigma, \mu), 1 \leqslant p<+\infty$, that

$$
\begin{align*}
A_{p}^{p}\left\|\left(\sum_{k=1}^{N}\left|f_{k}\right|^{2}\right)^{1 / 2}\right\|_{L^{p}(\mu)}^{p} & \leqslant \mathrm{E} \int_{X}\left|\sum_{k=1}^{N} \varepsilon_{k} f_{k}\right|^{p} \mathrm{~d} \mu \\
& \leqslant B_{p}^{p}\left\|\left(\sum_{k=1}^{N}\left|f_{k}\right|^{2}\right)^{1 / 2}\right\|_{L^{p}(\mu)}^{p} . \tag{1.27}
\end{align*}
$$

It is possible (see Pisier [64, Section 5.8]) to obtain the general case of unconditionality of martingale differences by approximating general martingale difference sequences by blocks of Bernoulli martingale differences. Also, one can see that (1.26) is self-dual and obtain by duality the Burkholder-Gundy inequalities for $2 \leqslant p<+\infty$.

The proof of Lemma 1.7 is valid with almost no change when the martingale takes values in a Hilbert space $H$, because $L^{2}(\Omega, \mathcal{F}, P, H)$ is a Hilbert space where the $H$-valued martingale differences are orthogonal. For values in a Banach space, two difficulties arise. First, the relevant "square function" has to be defined, and second, the Banach space-valued martingale differences are not unconditional in general. The Banach spaces where martingale differences are unconditional form a nice class of spaces, see Pisier [64, Chap. 5, The UMD property for Banach spaces].

Remark 1.9. - Let $f=\sum_{k=0}^{N} d_{k}$ be the sum of a Bernoulli martingale and let $g=\sum_{k=0}^{N} a_{k} d_{k}$ be obtained from $f$ by a martingale transform
operation, with $\left|a_{k}\right| \leqslant 1$ for $k=0, \ldots, N$. By Lemma 1.7 and Doob's inequality (1.1), we have

$$
\|g\|_{p} \leqslant\left\|g^{*}\right\|_{p} \leqslant 6\|S(g)\|_{p} \leqslant 6\|S(f)\|_{p} \leqslant 36\left\|f^{*}\right\|_{p} \leqslant \frac{36 p}{p-1}\|f\|_{p}, \quad 1<p \leqslant 2
$$

which shows that the constant $\kappa_{u, p}$ in (1.26) is of order $1 /(p-1)$ in this case. Actually, Burkholder has found the exact value of the unconditional constant for general martingale transforms and for every $p \in(1,+\infty)$. It is given by

$$
\kappa_{u, p}=p^{*}-1, \text { where } p^{*}:=\max (p, p /(p-1))
$$

One can consult [16] and the references given there to several other articles by Burkholder. One can also find in [16, Section 5.4] a bound $c_{p} \leqslant p^{*}-1$ for the constant $c_{p}$ in Theorem 1.6.

### 1.4.3. A consequence of the "reflection principle"

Consider a Brownian motion $\left(B_{s}\right)_{s \geqslant 0}$ on $\mathbb{R}$, defined on a probability space $(\Omega, \mathcal{F}, P)$ and with respect to a filtration $\left(\mathcal{F}_{s}\right)_{s \geqslant 0}$. We assume that $B_{0}=0$, we fix a real number $v>0$, and we let $S_{v}(\omega)$ denote the first time when the trajectory $s \mapsto B_{s}(\omega), s \geqslant 0$, which is continuous for almost every $\omega \in \Omega$, reaches the point $v$. It is clear that if $s_{0}>0$ is given, one has $\left\{B_{s_{0}} \geqslant v\right\} \subset\left\{S_{v} \leqslant s_{0}\right\}$, thus
$P\left(\left\{S_{v} \leqslant s_{0}\right\}\right) \geqslant P\left(\left\{B_{s_{0}} \geqslant v\right\}\right)=P\left(\left\{B_{1} \geqslant v / \sqrt{s_{0}}\right\}\right)=\int_{v / \sqrt{s_{0}}}^{+\infty} \mathrm{e}^{-y^{2} / 2} \frac{\mathrm{~d} y}{\sqrt{2 \pi}}$.
From now on, we write $P\left(S_{v} \leqslant s_{0}\right)$ for $P\left(\left\{S_{v} \leqslant s_{0}\right\}\right)$. We will show that actually

$$
P\left(S_{v} \leqslant s_{0}\right)=2 P\left(B_{s_{0}} \geqslant v\right)=2 \int_{v / \sqrt{s_{0}}}^{+\infty} \mathrm{e}^{-y^{2} / 2} \frac{\mathrm{~d} y}{\sqrt{2 \pi}}
$$

which proves in passing that $S_{v}$ is finite almost surely, since we have then

$$
P\left(S_{v}<+\infty\right)=2 \int_{0}^{+\infty} \mathrm{e}^{-y^{2} / 2} \frac{\mathrm{~d} y}{\sqrt{2 \pi}}=1
$$

The reasoning makes use of the reflection of the Brownian motion after a stopping time $\tau$. A stopping time is a random variable $\tau$ with values in $[0,+\infty]$, such that for every $t \geqslant 0$, the event $\{\tau \leqslant t\}$ belongs to the $\sigma$-field $\mathcal{F}_{t}$ of the past of time $t$. Intuitively, a stopping time corresponds to a decision to quit at time $\tau(\omega)$ that an observer, embarked on a path $t \mapsto X_{t}(\omega)$ of the random process $\left(X_{t}\right)_{t \geqslant 0}$ since the time $t=0$, can take from his only knowledge of what happened on his way between 0 and the present time.

The random time $S_{v}$ is an excellent example of stopping time, with a quite simple rule: I stop when I reach the point $v>0$.

The Brownian reflected after the random time $\tau$ changes its direction, its trajectory becomes the symmetric of the original trajectory with respect to the point $\left(B_{\tau}\right)(\omega):=B_{\tau(\omega)}(\omega)$ that was reached at time $\tau(\omega)$. Let us denote by $\left(B_{s}^{\tau}\right)_{s \geqslant 0}$ the reflected Brownian, given by

$$
\begin{aligned}
B_{s}^{\tau}(\omega) & =B_{s}(\omega) & & \text { if } 0 \leqslant s \leqslant \tau(\omega) \\
\frac{B_{s}^{\tau}(\omega)+B_{s}(\omega)}{2} & =B_{\tau(\omega)}(\omega) & & \text { if } s \geqslant \tau(\omega) .
\end{aligned}
$$

The reflected Brownian $B^{\tau}$ is still a Brownian motion. Consider first the simplest stopping time and reflection. Choosing a set $A_{1}$ in the $\sigma$-field $\mathcal{F}_{s_{1}}$ at time $s_{1}>0$, we define a stopping time $\tau_{1}$ equal to $s_{1}$ on $A_{1}$ and to $+\infty$ outside. The corresponding reflection $\left(B_{s}^{\tau_{1}}\right)_{s \geqslant 0}$ is given by

$$
\begin{aligned}
& B_{s}^{\tau_{1}}(\omega)=B_{s}(\omega) \quad \text { if } \quad 0 \leqslant s \leqslant s_{1} \quad \text { or } \quad \omega \notin A_{1}, \\
& \frac{B_{s}^{\tau_{1}}(\omega)+B_{s}(\omega)}{2}=B_{s_{1}}(\omega) \quad \text { if } s \geqslant s_{1} \quad \text { and } \omega \in A_{1} .
\end{aligned}
$$

One shows easily that $\left(B_{s}^{\tau_{1}}\right)_{s \geqslant 0}$ is a Brownian motion. Iterating this operation, one can reach discrete stopping times, and pass to the limit for dealing with general stopping times. Indeed, a stopping time $\tau$ can be approximated by the first time $\tau_{k}>\tau$ such that $2^{k} \tau_{k}$ is an integer, i.e., $\tau_{k}=2^{-k}\left(\left\lfloor 2^{k} \tau\right\rfloor+1\right)$, for every $k \in \mathbb{N}$.

Another important property that can be checked following the same route is the following: if $\tau$ is an almost surely finite stopping time, the process "starting afresh at time $\tau$ ", defined by $X_{s}=B_{\tau+s}-B_{\tau}$, i.e., $X_{s}(\omega)=$ $B_{\tau(\omega)+s}(\omega)-B_{\tau(\omega)}(\omega)$, is also a Brownian motion.

Consider the Brownian reflected after the stopping time $S_{v}$, with $v>0$. Since the Brownian paths are continuous and $B_{0}=0$, we have $B_{S_{v}(\omega)}(\omega)=v$ and for every $s_{0}>0$, the event $\left\{B_{s_{0}}>v\right\}$ is contained in $\left\{S_{v}<s_{0}\right\}$. Clearly, the event $\left\{B_{s_{0}}^{S_{v}}>v\right\}$ is also contained in $\left\{S_{v}<s_{0}\right\}$ and disjoint from $\left\{B_{s_{0}}>v\right\}$. Actually, since on the set $\left\{S_{v}<s_{0}\right\}$ one has $B_{s_{0}}^{S_{v}}+B_{s_{0}}=2 v$, one sees that

$$
\left\{S_{v}<s_{0}\right\} \backslash\left\{B_{s_{0}} \geqslant v\right\}=\left\{B_{s_{0}}^{S_{v}}>v\right\}
$$

The event $\left\{B_{s_{0}}^{S_{v}}>v\right\}$ has the same probability as $\left\{B_{s_{0}}>v\right\}$, since $\left(B_{s}^{S_{v}}\right)_{s \geqslant 0}$ is another Brownian, and $P\left(S_{v}=s_{0}\right) \leqslant P\left(B_{s_{0}}=v\right)=0$. We have therefore that

$$
P\left(S_{v} \leqslant s_{0}\right)=P\left(S_{v}<s_{0}\right)=2 P\left(B_{s_{0}}>v\right)=2 \int_{v}^{+\infty} \mathrm{e}^{-u^{2} /\left(2 s_{0}\right)} \frac{\mathrm{d} u}{\sqrt{2 \pi s_{0}}}
$$

Consequently, for every $s>0$, we obtain

$$
P\left(S_{v} \leqslant s\right)=P\left(\sup _{0 \leqslant u \leqslant s} B_{u} \geqslant v\right)=2 \int_{v / \sqrt{s}}^{+\infty} \mathrm{e}^{-y^{2} / 2} \frac{\mathrm{~d} y}{\sqrt{2 \pi}} .
$$

This allows us to find the density $h_{v}$ of the distribution of $S_{v}$, which is given by

$$
\begin{equation*}
h_{v}(s)=\mathbf{1}_{s \geqslant 0} \frac{v s^{-3 / 2}}{\sqrt{2 \pi}} \mathrm{e}^{-v^{2} /(2 s)}, \quad s \in \mathbb{R} . \tag{1.28}
\end{equation*}
$$

Remark. - A variant of the preceding reasoning applies to the exit time $S$ from an open convex subset $D$ of $\mathbb{R}^{n}$ containing the starting point $x_{0}$ of an $n$-dimensional Brownian motion. Suppose that this Brownian motion touches the boundary of $D$, for the first time, at the point $x=x(\omega)$ and at time $S(\omega)$. Let $E_{x}$ be an affine half-space tangent to $D$ at $x$, and exterior to $D$ (this $E_{x}$ is not unique in general). Starting again from $x$ at time $S(\omega)$, there is a probability $1 / 2$ to end in $E_{x}$ at time $s_{0}>S(\omega)$, so there is at least one chance out of two to end up outside $D$ at time $s_{0}$. The set $\left\{B_{s_{0}} \notin D\right\}$ is a subset of $\left\{S<s_{0}\right\}$ that occupies thus at least one half of it. We have therefore

$$
P\left(S<s_{0}\right) \leqslant 2 P\left(B_{s_{0}} \notin D\right)
$$

This inequality says that the probability to be outside $D$ at a time between 0 and $s_{0}$ is bounded by twice the probability to be outside $D$ at time $s_{0}$. This can be readily interpreted in terms of maximal function. If $\|\cdot\|_{C}$ denotes the norm on $\mathbb{R}^{n}$ associated to a symmetric convex body $C$ in $\mathbb{R}^{n}$, we deduce maximal inequalities in $L^{p}\left(\mathbb{R}^{n}\right)$ for the $\|\cdot\|_{C}$ norm of the martingale $\left(B_{s}\right)_{s \geqslant 0}$ that are better than Doob's inequality. Namely, for every $p>0$ we have

$$
\begin{aligned}
\mathrm{E} \max _{0 \leqslant s \leqslant s_{0}}\left\|B_{s}\right\|_{C}^{p} & =p \int_{0}^{+\infty} t^{p-1} P\left(\max _{0 \leqslant s \leqslant s_{0}}\left\|B_{s}\right\|_{C}>t\right) \mathrm{d} t \\
& \leqslant 2 p \int_{0}^{+\infty} t^{p-1} P\left(\left\|B_{s_{0}}\right\|_{C}>t\right) \mathrm{d} t=2 \mathrm{E}\left\|B_{s_{0}}\right\|_{C}^{p}
\end{aligned}
$$

For $p \leqslant 1$, there is no Doob's inequality in $L^{p}$, and when $p>1$, one has always that $2^{1 / p}<p /(p-1)$, because $(1-x) 2^{x}<(1-x) \mathrm{e}^{x} \leqslant 1$ for $0<x<1$.

One could get a similar estimate when the set $D$ is no longer convex, but has the property that for every boundary point $x$ of $D$, there is a cone $E_{x}$ based at $x$, disjoint from $D$ and with a solid angle bounded below by $\delta>0$ independent of $x$. If we measure the angle as the proportion of the unit sphere $S^{n-1}$ of $\mathbb{R}^{n}$ intersected by the cone $E_{x}-x$ based at 0 , then the constant 2 above has to be replaced by $\delta^{-1}$.

### 1.5. The Poisson semi-group

Let us recall that the $S$ chwartz class $\mathcal{S}\left(\mathbb{R}^{n}\right)$ consists of all $C^{\infty}$ functions $\varphi$ such that $\left(1+|x|^{k}\right) \varphi^{(\ell)}(x)$ is bounded on $\mathbb{R}^{n}$ for all integers $k, \ell \geqslant 0$. We shall denote by $\left(P_{t}\right)_{t \geqslant 0}$ the Poisson semi-group on $\mathbb{R}^{n}$, which can be defined, for $f$ in the Schwartz class $\mathcal{S}\left(\mathbb{R}^{n}\right)$, by

$$
\begin{equation*}
\left(P_{t} f\right)(x)=u(x, t), \quad x \in \mathbb{R}^{n}, t \geqslant 0 \tag{1.29}
\end{equation*}
$$

where $u(x, t)$ is the (bounded) harmonic extension of $f$ to the upper halfspace $H^{+}$of $\mathbb{R}^{n+1}$ formed by all $(x, t)$ with $x \in \mathbb{R}^{n}$ and $t \geqslant 0$. For $x \in \mathbb{R}^{n}$ one has $u(x, 0)=f(x), \Delta u(x, t)=0$ when $t>0$, and $u$ is continuous on $H^{+}$. The semi-group property $P_{t+s}=P_{t} P_{s}$ amounts to saying that the harmonic extension of the function $f_{s}$ defined on $\mathbb{R}^{n}$ by $f_{s}(x)=u(x, s)$ is given by $v(x, t)=u(x, t+s)$.

The Poisson semi-group is intimately related to the Brownian motion $\left(B_{s}\right)_{s \geqslant 0}$ in $\mathbb{R}^{n+1}$. If the Brownian $\left(B_{s}\right)_{s \geqslant 0}$ starts at time $s=0$ from the point $\left(x_{0}, t_{0}\right)$, where $x_{0} \in \mathbb{R}^{n}$ and $t_{0}>0$, we know that almost every path $s \mapsto B_{s}(\omega)$ will hit the hyperplane $H_{0}=\{t=0\}$ at some time $\tau_{t_{0}}(\omega)<+\infty$. If we decompose $B_{s}$ into $\left(x_{0}+X_{s}, t_{0}+T_{s}\right)$, then $T_{s}$ is a one-dimensional Brownian motion, starting from 0 at time 0 , and $X_{s}$ is a $n$-dimensional Brownian motion, starting from the point 0 in $\mathbb{R}^{n}$ and independent of $T_{s}$. The stopping time $\tau_{t_{0}}$ is the first time $s>0$ when $T_{s}=-t_{0}$. If $f$ is reasonable, for example continuous and bounded on $\mathbb{R}^{n}$, one sees that the (bounded) harmonic extension $u$ of $f$ to the upper half-space is given by

$$
u\left(x_{0}, t_{0}\right)=\operatorname{E} F\left(B_{\tau_{t_{0}}}\right)=\operatorname{E} f\left(x_{0}+X_{\tau_{t_{0}}}\right)=\int_{\Omega} f\left(x_{0}+X_{\tau_{t_{0}}(\omega)}(\omega)\right) \mathrm{d} P(\omega)
$$

where $F$ is defined on the hyperplane $H_{0}$ of $\mathbb{R}^{n+1}$ by $F(x, 0)=f(x)$ for every $x \in \mathbb{R}^{n}$. The Poisson probability measure $P_{t_{0}}(x) \mathrm{d} x$ on $\mathbb{R}^{n}$ is the distribution of $X_{\tau_{t_{0}}}$, distribution of the Brownian motion $\left(X_{s}\right)$ starting from $0 \in \mathbb{R}^{n}$ and stopped at time $\tau_{t_{0}}$, when $B_{s}$ reaches $H_{0}$. We shall employ the same notation $P_{t}$ for the semi-group, for the Poisson distribution on $\mathbb{R}^{n}$, and for its density $P_{t}(x)$. The operator $P_{t}$ is the convolution with the corresponding probability measure, it acts thus on $L^{p}\left(\mathbb{R}^{n}\right)$ for $1 \leqslant p \leqslant+\infty$. We shall say that $t$ is the parameter of $P_{t}$.

The distribution of the stopping time $\tau_{t_{0}}$ is clearly the same as the distribution of the first time $S_{t_{0}}$ when the one-dimensional Brownian motion starting from 0 reaches $t_{0}>0$, and we know by (1.28) the density $h_{t}$ of the distribution of $S_{t}$. The Poisson distribution $P_{t}$ on $\mathbb{R}^{n}$ is obtained by mixing Gaussian distributions on $\mathbb{R}^{n}$, distributions of $X_{s}$ at various times $s$, the mixing being done according to the distribution of $S_{t}$. In the portion of the space $\Omega$ where $s_{0} \leqslant \tau_{t} \leqslant s_{0}+\delta s$, the coordinate $x$ of the Brownian point
$B_{s}=\left(X_{s}, t+T_{s}\right)$ at time $\tau_{t}$ is approximately $X_{s_{0}}$, with probability of order $h_{t}\left(s_{0}\right) \delta s$, and $\left(X_{s}\right)_{s \geqslant 0}$ is independent of $\tau_{t}$. The point $(x, 0)=\left(X_{s_{0}}, 0\right)$ is the point where the Brownian $B_{s}$ touches the hyperplane $H_{0}$, knowing that $\tau_{t}=s_{0}$. This is the reason behind the subordination principle of the Poisson semi-group to the Gaussian semi-group, which implies in particular that the maximal function of the Poisson semi-group is bounded by that of the Gaussian semi-group $\left(G_{s}\right)_{s \geqslant 0}$ on $\mathbb{R}^{n}$. Indeed, we have by (1.28) that $P_{t}$ is "in the (closed) convex hull" of the Gaussian semi-group, since

$$
\begin{equation*}
P_{t}=\int_{0}^{+\infty} G_{s} \frac{t s^{-3 / 2}}{\sqrt{2 \pi}} \mathrm{e}^{-t^{2} /(2 s)} \mathrm{d} s \tag{1.30}
\end{equation*}
$$

It follows that

$$
\left|P_{t} * f\right| \leqslant \int_{0}^{+\infty}\left|G_{s} * f\right| \frac{t s^{-3 / 2}}{\sqrt{2 \pi}} \mathrm{e}^{-t^{2} /(2 s)} \mathrm{d} s \leqslant \sup _{u \geqslant 0}\left|G_{u} * f\right|
$$

We get a dimensionless estimate for the maximal function of the Poisson semi-group, consequence of the one in (1.20. $G^{*}$ ) for the Gaussian case. We have

$$
\begin{equation*}
\left\|\sup _{t>0}\left|P_{t} f\right|\right\|_{L^{p}\left(\mathbb{R}^{n}\right)} \leqslant \frac{p}{p-1}\left(\int_{\mathbb{R}^{n}}|f(x)|^{p} \mathrm{~d} x\right)^{1 / p} \tag{*}
\end{equation*}
$$

The remarks about comparing to $\mathrm{M} f$ are still in order here. Stein [73, Lemma 1, p. 48] proves (1.31. $P^{*}$ ) with different constants and in a different way, capable of easier generalizations to non Euclidean settings. He does not deal with the Gaussian maximal function, but applies the Hopf maximal inequality (1.10) to the Gaussian semi-group together with the subordination principle. Using subordination, Stein shows that the Poisson maximal function $P^{*} f=\sup _{t>0}\left|P_{t} f\right|$ is bounded by an average of expressions $t^{-1} \int_{0}^{t}\left(G_{s} f\right) \mathrm{d} s$ that are controlled by Hopf.

The formula (1.30) proves that the marginals of $P_{t}$ are other Poisson distributions: indeed, the mixing distribution, which has density $h_{t}$, does not depend on the dimension $n$, and the projections on $\mathbb{R}^{\ell}, 1 \leqslant \ell<n$, of Gaussian distributions $N\left(0, \sigma^{2} \mathrm{I}_{n}\right)$ on $\mathbb{R}^{n}$ are $N\left(0, \sigma^{2} \mathrm{I}_{\ell}\right)$ Gaussian distributions. We can also deduce the density of the distribution $P_{t}$ for each $t>0$, writing

$$
\begin{aligned}
P_{t}(x) & =\int_{0}^{+\infty} \mathrm{e}^{-|x|^{2} /(2 s)}(2 \pi s)^{-n / 2} \frac{t s^{-3 / 2}}{\sqrt{2 \pi}} \mathrm{e}^{-t^{2} /(2 s)} \mathrm{d} s \\
& =t \int_{0}^{+\infty}(2 \pi s)^{-n / 2-1 / 2} \mathrm{e}^{-\left(t^{2}+|x|^{2}\right) /(2 s)} \frac{\mathrm{d} s}{s}, \quad x \in \mathbb{R}^{n}
\end{aligned}
$$

Setting $u=s /\left(t^{2}+|x|^{2}\right)$, then $v=1 /(2 u)$, we get

$$
P_{t}(x)=t\left(\pi\left(t^{2}+|x|^{2}\right)\right)^{-(n+1) / 2} \int_{0}^{+\infty} \mathrm{e}^{-v} v^{(n+1) / 2} \frac{\mathrm{~d} v}{v}
$$

The Poisson kernel $P_{t}$ on $\mathbb{R}^{n}$ is thus given by the formula

$$
\begin{equation*}
P_{t}(x)=P_{t}^{(n)}(x)=\frac{\Gamma[(n+1) / 2]}{\pi^{(n+1) / 2}} \frac{t}{\left(t^{2}+|x|^{2}\right)^{(n+1) / 2}}, \quad x \in \mathbb{R}^{n}, t>0 \tag{1.32}
\end{equation*}
$$

In dimension $n=1$, the Poisson kernel is the Cauchy kernel, equal to

$$
\begin{equation*}
P_{t}(x)=P_{t}^{(1)}(x)=\frac{t}{\pi\left(t^{2}+x^{2}\right)}, \quad x \in \mathbb{R}, t>0 \tag{1.33.C}
\end{equation*}
$$

The coefficient that comes into the $n$-dimensional formula (1.32) satisfies the asymptotic estimate

$$
\frac{\Gamma[(n+1) / 2]}{\pi^{(n+1) / 2}} \simeq \sqrt{\frac{2}{\pi n}} \frac{1}{\omega_{n}}=\sqrt{\frac{2 n}{\pi}} \frac{1}{s_{n-1}}
$$

where $\omega_{n}$ is the volume of the unit ball in $\mathbb{R}^{n}$ and $s_{n-1}$ the ( $n-1$ )-dimensional measure of the unit sphere $S^{n-1}$ in $\mathbb{R}^{n}$, given by

$$
\begin{equation*}
\omega_{n}=\frac{\pi^{n / 2}}{(n / 2)!}:=\frac{\pi^{n / 2}}{\Gamma((n / 2)+1)}, \quad s_{n-1}=n \omega_{n} \tag{1.34}
\end{equation*}
$$

From this, we obtain estimates on the measure of Euclidean balls for the probability measure $P_{1}(x) \mathrm{d} x$ on $\mathbb{R}^{n}$. Writing $P_{1}(x)=F(|x|)$, we get an exact asymptotic estimate when the dimension $n$ tends to infinity: for $\nu>0$ fixed, we have

$$
\begin{aligned}
& \int_{\{|x|>\sqrt{n} / \nu\}} P_{1}(x) \mathrm{d} x \\
& =s_{n-1} \int_{\sqrt{n} / \nu}^{+\infty} r^{n-1} F(r) \mathrm{d} r \simeq \sqrt{\frac{2}{\pi}} \int_{\sqrt{n} / \nu}^{+\infty} \frac{\sqrt{n}}{r}\left(\frac{r^{2}}{1+r^{2}}\right)^{(n+1) / 2} \frac{\mathrm{~d} r}{r} \\
& =\sqrt{\frac{2}{\pi}} \int_{1 / \nu}^{+\infty}\left(1+\frac{1}{n u^{2}}\right)^{-(n+1) / 2} \frac{\mathrm{~d} u}{u^{2}}=\sqrt{\frac{2}{\pi}} \int_{0}^{\nu}\left(1+\frac{y^{2}}{n}\right)^{-(n+1) / 2} \mathrm{~d} y .
\end{aligned}
$$

Therefore, when $n$ tends to infinity, we see that

$$
\begin{equation*}
\int_{\{\nu|x|>\sqrt{n}\}} P_{1}(x) \mathrm{d} x \longrightarrow 2 \int_{0}^{\nu} \mathrm{e}^{-y^{2} / 2} \frac{\mathrm{~d} y}{\sqrt{2 \pi}} . \tag{1.35}
\end{equation*}
$$

## 2. General dimension free inequalities, second part

In this section, we gather results that depend on the Fourier transform. In order that the Fourier transform be isometric on $L^{2}\left(\mathbb{R}^{n}\right)$, we set

$$
\forall \xi \in \mathbb{R}^{n}, \quad \widehat{f}(\xi)=\int_{\mathbb{R}^{n}} f(x) \mathrm{e}^{-2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} x, \quad \widehat{\mu}(\xi)=\int_{\mathbb{R}^{n}} \mathrm{e}^{-2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} \mu(x)
$$

when $f$ is in $L^{1}\left(\mathbb{R}^{n}\right) \cap L^{2}\left(\mathbb{R}^{n}\right)$ or when $\mu$ is a bounded measure on $\mathbb{R}^{n}$. By the Plancherel theorem (some say Parseval's theorem), we know that
this defines a mapping from $L^{1}\left(\mathbb{R}^{n}\right) \cap L^{2}\left(\mathbb{R}^{n}\right)$ to $L^{2}\left(\mathbb{R}^{n}\right)$ that extends to a unitary transformation $\mathcal{F}$ of $L^{2}\left(\mathbb{R}^{n}\right)$. The inverse mapping $\mathcal{F}^{-1}$ of $\mathcal{F}$ sends every square integrable function $\xi \mapsto g(\xi)$ to $\mathcal{F}(\xi \mapsto g(-\xi))$, also expressible by $x \mapsto(\mathcal{F} g)(-x)$. We shall employ the notation $g^{\vee}=\mathcal{F}^{-1} g$ for the inverse Fourier transform.

The Plancherel-Parseval theorem extends to functions $f$ with values in a Euclidean space $F$, giving then an isometry from $L^{2}\left(\mathbb{R}^{n}, F\right)$ to itself. This is clear for instance by looking at coordinates in an orthonormal basis of $F$.

With this normalization of the Fourier transform, we have that

$$
\widehat{\gamma}_{n}(\xi)=\mathrm{e}^{-2 \pi^{2}|\xi|^{2}}, \quad \xi \in \mathbb{R}^{n}
$$

and the Fourier transform of the Poisson kernel $P_{t}$ on $\mathbb{R}^{n}$ is equal to e ${ }^{-2 \pi t|\xi|}$, for every $\xi \in \mathbb{R}^{n}$. Indeed, as the marginals on $\mathbb{R}$ of $P_{t}$ are Cauchy distributions with the same parameter $t$, we find by the residue theorem that

$$
\widehat{P}_{t}(\xi)=\int_{\mathbb{R}} \frac{t \mathrm{e}^{-2 \mathrm{i} \pi s|\xi|}}{\pi\left(t^{2}+s^{2}\right)} \mathrm{d} s=\mathrm{e}^{-2 \pi t|\xi|}
$$

This information on the Fourier transform gives another way of checking the semi-group property $P_{s} * P_{t}=P_{s+t}$ of Poisson distributions. Using the Fourier inversion formula, we notice for future use that the harmonic extension $u(x, t)=\left(P_{t} f\right)(x)$ of $f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$ considered in (1.29) can be written as

$$
\begin{equation*}
u(x, t)=\int_{\mathbb{R}^{n}} \mathrm{e}^{-2 \pi t|\xi|} \widehat{f}(\xi) \mathrm{e}^{2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} \xi, \quad x \in \mathbb{R}^{n}, t>0 \tag{2.1}
\end{equation*}
$$

### 2.1. Littlewood-Paley functions

The Littlewood-Paley function $g(f)$ associated to a function $f$ on $\mathbb{R}^{n}$ is defined by

$$
\forall x \in \mathbb{R}^{n}, \quad g(f)(x)=\left(\int_{0}^{+\infty}|t \nabla u(x, t)|^{2} \frac{\mathrm{~d} t}{t}\right)^{1 / 2}
$$

where $u$ is the harmonic extension of $f$ to the upper half-space in $\mathbb{R}^{n+1}$, and where $\nabla u$ is the gradient of $u$ in $\mathbb{R}^{n+1}$. The classical theory, see for example Zygmund [85, vol. 2] for the circle case in Chap. 14, §3 and Chap. 15, §2, indicates that the norm of $f$ in $L^{p}\left(\mathbb{R}^{n}\right), 1<p<+\infty$, is equivalent to that of $g(f)$. One has that

$$
\begin{equation*}
\kappa_{p}^{-1}\|f\|_{p} \leqslant\|g(f)\|_{p} \leqslant \kappa_{p}\|f\|_{p} \tag{2.2}
\end{equation*}
$$

with a constant $\kappa_{p}$ depending on $p$, but independent of the dimension $n$. A variant of this Littlewood-Paley function is defined by

$$
\begin{equation*}
g_{1}(f)^{2}=\int_{0}^{+\infty}\left|t \frac{\partial}{\partial t} P_{t} f\right|^{2} \frac{\mathrm{~d} t}{t} \tag{2.3}
\end{equation*}
$$

It is clear that $g_{1}(f) \leqslant g(f)$, since $(\partial / \partial t)\left(P_{t} f\right)$ is a coordinate of the vector $\nabla u$. The function $g_{1}$ is one of the variants studied by Stein [73]. More generally, for every integer $k \geqslant 1$, Stein sets

$$
g_{k}(f)^{2}=\int_{0}^{+\infty}\left|t^{k} \frac{\partial^{k}}{\partial t^{k}} P_{t} f\right|^{2} \frac{\mathrm{~d} t}{t}
$$

Let us define $Q_{j}=P_{2^{j}}-P_{2^{j+1}}$, for every $j \in \mathbb{Z}$. Since

$$
\sum_{j \in \mathbb{Z}}\left|Q_{j} f\right|^{2}=\sum_{j \in \mathbb{Z}}\left|\int_{2^{j}}^{2^{j+1}}\left(\frac{\partial}{\partial t} P_{t} f\right) \mathrm{d} t\right|^{2},
$$

we obtain by Cauchy-Schwarz that

$$
\sum_{j \in \mathbb{Z}}\left|Q_{j} f\right|^{2} \leqslant \sum_{j \in \mathbb{Z}} 2^{j} \int_{2^{j}}^{2^{j+1}}\left|\frac{\partial}{\partial t} P_{t} f\right|^{2} \mathrm{~d} t \leqslant \sum_{j \in \mathbb{Z}} \int_{2^{j}}^{2^{j+1}}\left|\frac{\partial}{\partial t} P_{t} f\right|^{2} t \mathrm{~d} t=g_{1}(f)^{2}
$$

The classical result (2.2) on $g(f)$ implies that for $1<p<+\infty$, there exists a constant $\mathrm{q}_{p}$ independent of the dimension $n$ such that

$$
\begin{equation*}
\left\|\left(\sum_{j \in \mathbb{Z}}\left|Q_{j} f\right|^{2}\right)^{1 / 2}\right\|_{L^{p}\left(\mathbb{R}^{n}\right)} \leqslant \mathrm{q}_{p}\|f\|_{L^{p}\left(\mathbb{R}^{n}\right)}, \quad f \in L^{p}\left(\mathbb{R}^{n}\right) \tag{2.4}
\end{equation*}
$$

Observe that the same proof implies that a similar inequality, with a different constant depending on $c>1$, will hold for differences of the form $\widetilde{Q}_{j}=P_{t_{j}}-P_{t_{j-1}}$, where $\left(t_{j}\right)_{j \in \mathbb{Z}}$ is an increasing sequence of positive real numbers, provided that we have $t_{j+1} \leqslant c t_{j}$ for all $j$ s. On the other hand, by Rota's argument (1.15), one can obtain (2.4) from the Burkholder-Gundy inequalities of Theorem 1.6. Inequalities similar to (2.4) would hold for the Gaussian semi-group $\left(G_{t}\right)_{t \geqslant 0}$ defined in (1.19). Let us fix $T>0$. We have seen that $G_{2 t} f, 0 \leqslant t \leqslant T$, is the projection on the $\sigma$-field $\mathcal{G}_{T}$ generated by $B_{T}$ of the member $M_{T-t}$ of the Brownian martingale $M_{s}=\left(P_{T-s} f\right)\left(B_{s}\right)$, $0 \leqslant s \leqslant T$, running under the infinite invariant measure given by the Lebesgue measure on $\mathbb{R}^{n}$. We then apply (1.15). Using Gaussian $Q_{j}$ s would allow us to avoid a few minor technical difficulties later, and this is essentially what Bourgain [13] does for the cube problem, see Section 8 .

Relying on (1.15) and Remark 1.9 gives for the constant $q_{p}$ in (2.4) an upper bound of order $p /(p-1)$ when $p \rightarrow 1$. This can also be obtained if one follows Stein [73, p. 48-51]. When $1<p \leqslant 2$, the proof given there yields $\|g(f)\|_{p} \leqslant(p-1)^{-1 / 2} \mathrm{p}_{p}^{1-p / 2}\|f\|_{p}$ for the right-hand side inequality in (2.2),
where $\mathrm{p}_{p}$ is the constant in the maximal $L^{p}$-inequality for the Poisson semigroup. Since we have $\mathrm{p}_{p} \leqslant p /(p-1)$ by $\left(1.31 . P^{*}\right)$, we get that

$$
\begin{equation*}
\mathrm{q}_{p} \leqslant p /(p-1) \text { when } 1<p \leqslant 2 . \tag{2.5}
\end{equation*}
$$

Looking at the Fourier side, we see that $\sum_{j \in \mathbb{Z}} \widehat{Q}_{j}(\xi)=1$ for every $\xi \neq 0$, since $\widehat{P_{2^{j}}}(\xi)=\mathrm{e}^{-2^{j+1} \pi|\xi|}$ tends to 1 when $j \rightarrow-\infty$ and to 0 when $j \rightarrow+\infty$. It implies for the convolution operators, still denoted by $Q_{j}$, that

$$
\begin{equation*}
\sum_{j \in \mathbb{Z}} Q_{j}=\mathrm{Id} \tag{2.6}
\end{equation*}
$$

### 2.1.1. Littlewood-Paley and maximal functions

Stein [73, Chap. III, §3, p. 75] explains how to get $L^{p}$ estimates for several maximal functions related to semi-groups, by using the Littlewood-Paley functions. Consider a continuous function $\varphi$ on the half-line $[0,+\infty)$, differentiable on $(0,+\infty)$, and denote by $\Phi$ its antiderivative vanishing at 0 . For every $t>0$, one has

$$
t \varphi(t)=\int_{0}^{t}(s \varphi(s))^{\prime} \mathrm{d} s=\int_{0}^{t} \varphi(s) \mathrm{d} s+\int_{0}^{t} s \varphi^{\prime}(s) \mathrm{d} s=\Phi(t)+\int_{0}^{t} s \varphi^{\prime}(s) \mathrm{d} s
$$

Comparing $L^{1}$ and $L^{2}$ norms, one sees that

$$
\int_{0}^{t}\left|s \varphi^{\prime}(s)\right| \frac{\mathrm{d} s}{t} \leqslant\left(\int_{0}^{t}\left|s \varphi^{\prime}(s)\right|^{2} \frac{\mathrm{~d} s}{t}\right)^{1 / 2} \leqslant\left(\int_{0}^{t}\left|s \varphi^{\prime}(s)\right|^{2} \frac{\mathrm{~d} s}{s}\right)^{1 / 2}
$$

Therefore, one has

$$
|\varphi(t)| \leqslant \frac{|\Phi(t)|}{t}+\left(\int_{0}^{+\infty}\left|s \varphi^{\prime}(s)\right|^{2} \frac{\mathrm{~d} s}{s}\right)^{1 / 2}, \quad t>0
$$

One gets that

$$
\sup _{t>0}|\varphi(t)| \leqslant \sup _{t>0} \frac{|\Phi(t)|}{t}+\left(\int_{0}^{+\infty}\left|s \varphi^{\prime}(s)\right|^{2} \frac{\mathrm{~d} s}{s}\right)^{1 / 2}
$$

If $\varphi(s)=\left(P_{s} f\right)(x)$ for a given $x \in \mathbb{R}^{n}$, the upper bound becomes

$$
\sup _{t>0}\left|\left(P_{t} f\right)(x)\right| \leqslant \sup _{t>0} \frac{1}{t}\left|\int_{0}^{t}\left(P_{s} f\right)(x) \mathrm{d} s\right|+g_{1}(f)(x)
$$

One can (again) control the norm in $L^{p}, 1<p<+\infty$, of the maximal function of the Poisson semi-group, by the Hopf maximal inequality and the estimate for the Littlewood-Paley function. This control is easy in $L^{2}$, especially when $L^{2}$ admits an orthonormal basis $\left(f_{j}\right)$ such that $P_{t} f_{j}=\mathrm{e}^{-t \lambda_{j}} f_{j}$ for every $j, \lambda_{j} \geqslant 0$, for example in the case of the Laplacian on a bounded
domain $\Omega \subset \mathbb{R}^{n}$. If $f=\sum_{j} a_{j} f_{j}$ in $L^{2}(\Omega)$, one has $P_{t} f=\sum_{j} a_{j} \mathrm{e}^{-\lambda_{j} t} f_{j}$, and

$$
\begin{aligned}
& \int_{\Omega} g_{1}(f)(x)^{2} \mathrm{~d} x=\int_{0}^{+\infty} \int_{\Omega}\left|\sum_{j} a_{j} t \lambda_{j} \mathrm{e}^{-t \lambda_{j}} f_{j}(x)\right|^{2} \mathrm{~d} x \frac{\mathrm{~d} t}{t} \\
& \quad=\int_{0}^{+\infty}\left(\sum_{j}\left|a_{j}\right|^{2} t^{2} \lambda_{j}^{2} \mathrm{e}^{-2 t \lambda_{j}}\right) \frac{\mathrm{d} t}{t}=\sum_{j}\left|a_{j}\right|^{2} \int_{0}^{+\infty} t^{2} \lambda_{j}^{2} \mathrm{e}^{-2 t \lambda_{j}} \frac{\mathrm{~d} t}{t} \\
& \quad=\left(\int_{0}^{+\infty} u^{2} \mathrm{e}^{-2 u} \frac{\mathrm{~d} u}{u}\right) \sum_{\lambda_{j}>0}\left|a_{j}\right|^{2} \leqslant \frac{\Gamma(2)}{4}\|f\|_{2}^{2}=\frac{1}{4}\|f\|_{2}^{2}
\end{aligned}
$$

For the other Littlewood-Paley functions $g_{k}(f)$, one has in the same way

$$
\begin{aligned}
\int_{\Omega} g_{k}(f)(x)^{2} \mathrm{~d} x & =\int_{0}^{+\infty} \int_{\Omega}\left|\sum_{j} a_{j} t^{k} \lambda_{j}^{k} \mathrm{e}^{-t \lambda_{j}} f_{j}(x)\right|^{2} \mathrm{~d} x \frac{\mathrm{~d} t}{t} \\
& =\sum_{j}\left|a_{j}\right|^{2} \int_{0}^{+\infty} t^{2 k} \lambda_{j}^{2 k} \mathrm{e}^{-2 t \lambda_{j}} \frac{\mathrm{~d} t}{t} \leqslant \frac{\Gamma(2 k)}{4^{k}}\|f\|_{2}^{2}
\end{aligned}
$$

One can also work on $\mathbb{R}^{n}$ by Fourier transform with Parseval. One gets

$$
\begin{aligned}
\int_{\mathbb{R}^{n}} g_{k}(f)(x)^{2} \mathrm{~d} x & =\left.\left.(2 \pi)^{2 k} \int_{0}^{+\infty} \int_{\mathbb{R}^{n}}\left|\widehat{f}(\xi) t^{k}\right| \xi\right|^{k} \mathrm{e}^{-2 \pi t|\xi|}\right|^{2} \mathrm{~d} \xi \frac{\mathrm{~d} t}{t} \\
& =\frac{\Gamma(2 k)}{4^{k}}\|f\|_{2}^{2}
\end{aligned}
$$

We have also other relations like

$$
t^{2} \varphi^{\prime}(t)=\int_{0}^{t}\left(s^{2} \varphi^{\prime}(s)\right)^{\prime} \mathrm{d} s=2 \int_{0}^{t} s \varphi^{\prime}(s) \mathrm{d} s+\int_{0}^{t} s^{2} \varphi^{\prime \prime}(s) \mathrm{d} s
$$

implying that

$$
\sup _{t>0}\left|t \varphi^{\prime}(t)\right| \leqslant 2 \int_{0}^{+\infty}\left|s \varphi^{\prime}(s)\right|^{2} \frac{\mathrm{~d} s}{s}+\int_{0}^{+\infty}\left|s^{2} \varphi^{\prime \prime}(s)\right|^{2} \frac{\mathrm{~d} s}{s} .
$$

This brings back the successive maximal functions associated with each of the expressions $t^{k} \partial^{k} / \partial t^{k}\left(P_{t} f\right), k \geqslant 1$, to quantities that can be estimated or are already estimated, as in

$$
\sup _{t>0}\left|t \frac{\partial}{\partial t}\left(P_{t} f\right)(x)\right| \leqslant 2 g_{1}(f)(x)+g_{2}(f)(x), \quad x \in \mathbb{R}^{n}
$$

### 2.2. Fourier multipliers

We introduce two dilation operators that appear in duality, for instance when dealing with the Fourier transform. Given a function $g$ on $\mathbb{R}^{n}$ and
$\lambda>0$, we use for these operations the notation

$$
\begin{equation*}
g_{(\lambda)}(x)=\lambda^{-n} g\left(\lambda^{-1} x\right), \quad g_{[\lambda]}(x)=g(\lambda x), \quad x \in \mathbb{R}^{n} \tag{2.7}
\end{equation*}
$$

If $g$ already has a subscript, as in $g=g_{1}$, we shall use the heavier notation $\left(g_{1}\right)_{(\lambda)}$ or $\left(g_{1}\right)_{[\lambda]}$. One sees, for example when $g$ is integrable and $h$ bounded, that

$$
\int_{\mathbb{R}^{n}} g_{(\lambda)}(x) h(x) \mathrm{d} x=\int_{\mathbb{R}^{n}} g(y) h_{[\lambda]}(y) \mathrm{d} y, \quad \text { and } \widehat{\left(g_{(\lambda)}\right)}(\xi)=\widehat{g}(\lambda \xi), \quad \xi \in \mathbb{R}^{n}
$$

that is to say, we have $\widehat{\left(g_{(\lambda)}\right)}=(\widehat{g})_{[\lambda]}$. Clearly, $g_{(\lambda \mu)}=\left(g_{(\lambda)}\right)_{(\mu)}$. The $g_{(\lambda)}$ dilation preserves the integral of $g$; it is extended to measures $\mu$ on $\mathbb{R}^{n}$ by setting $\mu_{(\lambda)}(f)=\mu\left(f_{[\lambda]}\right)$, namely

$$
\begin{equation*}
\int_{\mathbb{R}^{n}} f(x) \mathrm{d} \mu_{(\lambda)}(x)=\int_{\mathbb{R}^{n}} f(\lambda x) \mathrm{d} \mu(x) \tag{2.8}
\end{equation*}
$$

for every $f$ in the space $\mathcal{K}\left(\mathbb{R}^{n}\right)$ of continuous and compactly supported functions. The measure $\mu_{(\lambda)}$ is the image of $\mu$ under the mapping $\mathbb{R}^{n} \ni x \mapsto \lambda x$. If $\mathrm{d} \mu(x)=g(x) \mathrm{d} x$, then $g_{(\lambda)}$ is the density of $\mu_{(\lambda)}$.

Let $\xi \mapsto m(\xi)$ belong to $L^{\infty}\left(\mathbb{R}^{n}\right)$. For $f \in L^{2}\left(\mathbb{R}^{n}\right)$, we have $\widehat{f} \in L^{2}\left(\mathbb{R}^{n}\right)$ by Plancherel, $\xi \mapsto m(\xi) \widehat{f}(\xi)$ is also in $L^{2}\left(\mathbb{R}^{n}\right)$ and is therefore the Fourier transform of some function $T_{m} f \in L^{2}\left(\mathbb{R}^{n}\right)$. We thus get a linear operator $T_{m}$ on $L^{2}\left(\mathbb{R}^{n}\right)$ if we define $T_{m} f$, for every $f \in L^{2}\left(\mathbb{R}^{n}\right)$, by means of its Fourier transform, letting

$$
\left(T_{m} f\right)^{\wedge}(\xi)=m(\xi) \widehat{f}(\xi), \quad \xi \in \mathbb{R}^{n}
$$

Let $P_{m}$ be the operator of multiplication by $m$, defined by $P_{m} \varphi=m \varphi$. The operator $T_{m}=\mathcal{F}^{-1} P_{m} \mathcal{F}$ is bounded on $L^{2}\left(\mathbb{R}^{n}\right)$ since by Parseval, one has that

$$
\begin{equation*}
\int_{\mathbb{R}^{n}}\left|\left(T_{m} f\right)(x)\right|^{2} \mathrm{~d} x=\int_{\mathbb{R}^{n}}|m(\xi)|^{2}|\widehat{f}(\xi)|^{2} \mathrm{~d} \xi \leqslant\|m\|_{\infty}^{2}\|f\|_{2}^{2} \tag{2.9}
\end{equation*}
$$

We shall say that $T_{m}$ is the operator associated to the multiplier $m$.
One can ask whether $T_{m}$ also operates as a bounded mapping on certain $L^{p}$ spaces. In this survey, "bounded on $L^{p "}$ will always mean bounded from $L^{p}$ to $L^{p}$. Let $q$ be the conjugate exponent of $p$, defined by $1 / q+1 / p=1$. Assuming that $1<p<+\infty$, we see that $T_{m}$ is bounded on $L^{p}\left(\mathbb{R}^{n}\right)$ if and only if $\int_{\mathbb{R}^{n}} m(\xi) \widehat{\varphi}(\xi) \widehat{\psi}(\xi) \mathrm{d} \xi$ is uniformly bounded when $\varphi, \psi \in \mathcal{S}\left(\mathbb{R}^{n}\right)$ belong to the unit balls of $L^{p}\left(\mathbb{R}^{n}\right)$ and $L^{q}\left(\mathbb{R}^{n}\right)$ respectively, hence $T_{m}$ is then also bounded on $L^{q}\left(\mathbb{R}^{n}\right)$ (and on $L^{2}\left(\mathbb{R}^{n}\right)$ by interpolation, so $m$ has to be a bounded function, see the line after (2.12.P)).

We now observe that the multiplier $m$ and its dilates $m_{[\lambda]}: \xi \mapsto m(\lambda \xi)$, $\lambda>0$, define operators having equal norms on $L^{p}\left(\mathbb{R}^{n}\right)$. We see that

$$
\left(T_{m_{[\lambda]}} f_{(\lambda)}\right)^{\wedge}(\xi)=m(\lambda \xi) \widehat{f}(\lambda \xi)
$$

hence $T_{m_{[\lambda]}} f_{(\lambda)}=\left(T_{m} f\right)_{(\lambda)}$. Consider the operator $S_{\lambda}: f \mapsto f_{(\lambda)}$. For every $p \in[1,+\infty]$ and $1 / q+1 / p=1$, the multiple $S_{\lambda, p}:=\lambda^{n / q} S_{\lambda}$ of $S_{\lambda}$ is an isometric bijection of $L^{p}\left(\mathbb{R}^{n}\right)$ onto itself. The relation $S_{\lambda} \circ T_{m}=T_{m_{[\lambda]}} \circ S_{\lambda}$ becomes

$$
\begin{equation*}
T_{m_{[\lambda]}}=S_{\lambda, p} T_{m} S_{\lambda, p}^{-1} \tag{2.10}
\end{equation*}
$$

and this implies that $T_{m}$ and $T_{m_{[\lambda]}}$ have the same norm on $L^{p}\left(\mathbb{R}^{n}\right)$. More generally, let $\mathbf{m}=\left(m^{(j)}\right)_{j \in J}$ be a family of multipliers and define $T_{\mathbf{m}} f=$ $\sup _{j \in J}\left|T_{m^{(j)}} f\right|$. If we set $\mathbf{m}_{[\lambda]}=\left(m_{[\lambda]}^{(j)}\right)_{j \in J}$, then we have again that

$$
\begin{equation*}
T_{\mathbf{m}_{[\lambda]}}=S_{\lambda, p} T_{\mathbf{m}} S_{\lambda, p}^{-1} \tag{2.11}
\end{equation*}
$$

because $S_{\lambda}$ commutes with $f \mapsto|f|$ and $S_{\lambda}\left(\sup _{j \in J} f_{j}\right)=\sup _{j \in J} S_{\lambda} f_{j}$. Consequently, $T_{\mathbf{m}_{[\lambda]}}$ and $T_{\mathbf{m}}$ also have the same norm on $L^{p}\left(\mathbb{R}^{n}\right)$.

We shall speak of the action on $L^{p}$ of the multiplier $m$ and set

$$
\|m\|_{p \rightarrow p}:=\left\|T_{m}\right\|_{p \rightarrow p} .
$$

If $T_{m}$ is bounded on $L^{p}$, one says that $m$ is a multiplier on $L^{p}$, or a $L^{p}$ multiplier. The next lemma will be useful, it is nothing but a direct consequence of the equality $\left\|m_{[\lambda]}\right\|_{p \rightarrow p}=\|m\|_{p \rightarrow p}$ for every $\lambda>0$, and of the triangle inequality in $L^{p}$.

Lemma 2.1. - Suppose that $1 \leqslant p \leqslant+\infty$ and that $m(\xi)$ is a $L^{p}\left(\mathbb{R}^{n}\right)$ multiplier. If the function $\psi$ is integrable on $(0,+\infty)$, the multiplier $N$ defined by

$$
N(\xi)=\int_{0}^{+\infty} \psi(\lambda) m(\lambda \xi) \mathrm{d} \lambda, \quad \xi \in \mathbb{R}^{n}
$$

is a $L^{p}\left(\mathbb{R}^{n}\right)$-multiplier and $\|N\|_{p \rightarrow p} \leqslant\|\psi\|_{L^{1}(0,+\infty)}\|m\|_{p \rightarrow p}$.
Note that clearly, multiplier operators commute to each other, and commute to translations and differentiations. We will apply many times the easy fact (2.9), which can be written as

$$
\begin{equation*}
\|m\|_{2 \rightarrow 2}=\left\|T_{m}\right\|_{2 \rightarrow 2} \leqslant\|m\|_{L^{\infty}\left(\mathbb{R}^{n}\right)} . \tag{2.12.P}
\end{equation*}
$$

The inequality is actually an equality, since by Parseval, the norm of $T_{m}$ on $L^{2}\left(\mathbb{R}^{n}\right)$ is equal to that of $P_{m}$, the multiplication operator by $m$.

If $K$ is a function integrable on $\mathbb{R}^{n}$, it acts by convolution on $L^{p}\left(\mathbb{R}^{n}\right)$ for all values $1 \leqslant p \leqslant+\infty$, and one gets easily by convexity of the $L^{p}$ norm that

$$
\begin{equation*}
\|K * f\|_{L^{p}\left(\mathbb{R}^{n}\right)} \leqslant\|K\|_{L^{1}\left(\mathbb{R}^{n}\right)}\|f\|_{L^{p}\left(\mathbb{R}^{n}\right)} . \tag{2.13}
\end{equation*}
$$

This is an easy example of operator associated to a multiplier, since convolution of $f$ with $K$ corresponds to multiplication of $\widehat{f}$ by $\widehat{K}$. The Fourier transform $m=\widehat{K}$ of $K$ is thus a multiplier on all spaces $L^{p}\left(\mathbb{R}^{n}\right), 1 \leqslant p \leqslant+\infty$. Consider the Fourier transform $m$ of the convolution kernel $K \in L^{1}\left(\mathbb{R}^{n}\right)$, equal to

$$
m(\xi)=\int_{\mathbb{R}^{n}} K(x) \mathrm{e}^{-2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} x, \quad \xi \in \mathbb{R}^{n}
$$

For $\xi \neq 0$, let $\xi=|\xi| \theta$ and $x=y+s \theta$, where $y$ is in the hyperplane $\theta^{\perp}$ orthogonal to $\theta \in S^{n-1}$, and $s \in \mathbb{R}$. By Fubini, we have for every real number $u$ that

$$
m(u \xi)=\int_{\mathbb{R}}\left(\int_{\theta^{\perp}} K(y+s \theta) \mathrm{d}^{n-1} y\right) \mathrm{e}^{-2 \mathrm{i} \pi s u|\xi|} \mathrm{d} s
$$

where $\mathrm{d}^{n-1} y$ denotes the normalized Lebesgue measure on the Euclidean space $\theta^{\perp} \subset \mathbb{R}^{n}$. In what follows we associate to $K$ and to $\theta$ in the unit sphere $S^{n-1}$ the function $\varphi_{\theta, K}$ defined on $\mathbb{R}$ by

$$
\begin{equation*}
\forall s \in \mathbb{R}, \quad \varphi_{\theta, K}(s):=\int_{\theta^{\perp}} K(y+s \theta) \mathrm{d}^{n-1} y \tag{2.14}
\end{equation*}
$$

so that for $\xi \neq 0$ and $\theta=|\xi|^{-1} \xi$, letting $\varphi_{\theta}=\varphi_{\theta, K}$ we have

$$
\begin{equation*}
m(u \xi)=\int_{\mathbb{R}} \varphi_{\theta}(s) \mathrm{e}^{-2 \mathrm{i} \pi s u|\xi|} \mathrm{d} s=\int_{\mathbb{R}} \frac{1}{|\xi|} \varphi_{\theta}\left(\frac{v}{|\xi|}\right) \mathrm{e}^{-2 \mathrm{i} \pi v u} \mathrm{~d} v \tag{2.15}
\end{equation*}
$$

The function $\mathbb{R} \ni u \mapsto m(u \theta)$ is the Fourier transform (in dimension 1) of $\varphi_{\theta}$.

### 2.2.1. Multipliers "of Laplace type"

We consider a scalar function $F$ on $(0,+\infty)$ that admits an expression of the form

$$
\begin{equation*}
\forall \lambda>0, \quad F(\lambda)=\lambda \int_{0}^{+\infty} \mathrm{e}^{-\lambda t} a(t) \mathrm{d} t \tag{2.16}
\end{equation*}
$$

where $a$ is a measurable function bounded on $(0,+\infty)$. The multiplier $m(\xi)$ "of Laplace type" associated to $F$ is defined by $m(\xi)=F(|\xi|)$, for $\xi \in \mathbb{R}^{n}$. We note that $\|F\|_{\infty} \leqslant\|a\|_{\infty}$, thus by (2.12.P), this multiplier $m$ is bounded on $L^{2}\left(\mathbb{R}^{n}\right)$ with operator norm $\leqslant\|a\|_{\infty}$. Stein proves the following result.

Proposition 2.2 ([73, Theorem $3^{\prime}$, p. 58]). - Let $F$ be defined on $(0,+\infty)$ by $(2.16)$, for some function $a \in L^{\infty}(0,+\infty)$. The operator $T_{m}$ associated to the multiplier $m(\xi)=F(|\xi|)$ is bounded on $L^{p}\left(\mathbb{R}^{n}\right)$ for $1<p<+\infty$ and

$$
\left\|T_{m}\right\|_{p \rightarrow p} \leqslant \lambda_{p}\|a\|_{\infty}
$$

where $\lambda_{p}$ is a constant independent of the dimension $n$.

The identity operator belongs to this class (when $a(t) \equiv 1$ ), we thus see that $\lambda_{p} \geqslant 1$ for every $p$. It follows from the proposition that the imaginary powers of $(-\Delta)^{1 / 2}$ act on the spaces $L^{p}\left(\mathbb{R}^{n}\right)$ when $1<p<+\infty$, with norms bounded independently of the dimension $n$. Indeed, we have the formula of Laplace type

$$
\begin{equation*}
\lambda^{\mathrm{i} b}=\frac{1}{\Gamma(1-\mathrm{i} b)} \lambda \int_{0}^{+\infty} \mathrm{e}^{-\lambda t} t^{-\mathrm{i} b} \mathrm{~d} t, \quad \lambda>0, \quad a(t)=\frac{t^{-\mathrm{i} b}}{\Gamma(1-\mathrm{i} b)}, \tag{2.17}
\end{equation*}
$$

hence $\|a\|_{\infty}=|\Gamma(1-\mathrm{i} b)|^{-1}$, for every $b \in \mathbb{R}$. According to the estimate (3.4) for the Gamma function, we get from Proposition 2.2 that
$\forall b \in \mathbb{R}, \quad\left\||\xi|^{i b}\right\|_{p \rightarrow p} \leqslant \lambda_{p}\left(1+b^{2}\right)^{-1 / 2} \mathrm{e}^{\pi|b| / 2}, \quad 1<p<+\infty$.
Stein's proof of Proposition 2.2 draws on $L^{p}$ inequalities for the LittlewoodPaley functions $g_{1}(f)$ and $g_{2}(f)$, and a comparison $g_{1}\left(T_{m} f\right) \leqslant \kappa g_{2}(f)$. We now sketch another possibility, which invokes martingale inequalities. If $F$ is as in Proposition 2.2 and $m(\xi)=F(|\xi|)$, then $T_{m} f$, for $f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$, can be expressed by

$$
\begin{equation*}
-\left(T_{m_{[2 \pi]}} f\right)=\int_{0}^{+\infty} a(t)\left(\frac{\partial}{\partial t} P_{t} f\right) \mathrm{d} t \tag{2.19}
\end{equation*}
$$

Indeed, we know by (2.1) that

$$
\left(P_{t} f\right)(x)=u(x, t)=\int_{\mathbb{R}^{n}} \mathrm{e}^{-2 \pi t|\xi|} \widehat{f}(\xi) \mathrm{e}^{2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} \xi
$$

and

$$
\begin{aligned}
\left(\int_{0}^{+\infty} a(t)( \right. & \left.\left.\frac{\partial}{\partial t} P_{t} f\right) \mathrm{~d} t\right)(x) \\
& =\int_{0}^{+\infty} a(t)\left(\int_{\mathbb{R}^{n}}(-2 \pi|\xi|) \mathrm{e}^{-2 \pi t|\xi|} \widehat{f}(\xi) \mathrm{e}^{2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} \xi\right) \mathrm{d} t \\
& =-\int_{\mathbb{R}^{n}} F(2 \pi|\xi|) \widehat{f}(\xi) \mathrm{e}^{2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} \xi=-\left(T_{m_{[2 \pi]}} f\right)(x)
\end{aligned}
$$

Suppose that $a$ is a step function supported in $\left[t_{0}, t_{N}\right] \subset[0,+\infty)$. Then

$$
a(t)=\sum_{j=1}^{N} a_{j} \mathbf{1}_{\left[t_{j-1}, t_{j}\right)}(t),
$$

with $0=t_{0}<t_{1}<\ldots<t_{N}$. By (2.19), we obtain that

$$
-T_{m_{[2 \pi]}} f=\sum_{j=1}^{N} a_{j}\left(P_{t_{j}}-P_{t_{j-1}}\right)(f) .
$$

It follows that $T_{m} f$ can be considered as projection of a martingale transform by a conditional expectation $\mathrm{E}_{\mathcal{G}}$. Let $u_{j}=t_{j} / 2, j=0, \ldots, N$, and
$T:=u_{N}$. We have seen in (1.13) that $P_{t_{j}} f=P_{2 u_{j}} f$ is the image under the projection $\mathrm{E}_{\mathcal{G}}$ of the martingale member $M_{T-u_{j}}=\left(P_{u_{j}} f\right)\left(X_{T-u_{j}}\right)$, so letting $L_{i}=M_{T-u_{N-i}}, i=0, \ldots, N$, we see that $T_{m_{[2 \pi]}} f$ is equal to

$$
\mathrm{E}_{\mathcal{G}}\left(\sum_{j=1}^{N} a_{j}\left(M_{T-u_{j-1}}-M_{T-u_{j}}\right)\right)=\mathrm{E}_{\mathcal{G}}\left(\sum_{i=1}^{N} a_{N-i+1}\left(L_{i}-L_{i-1}\right)\right),
$$

which is the transform of the martingale $\left(L_{i}\right)_{i=0}^{N}$ by the bounded nonrandom multipliers $\left(a_{N-i+1}\right)_{i=1}^{N}$. Also, $L_{N}$ is equal to $M_{T}=f\left(X_{T}\right)$ that has the distribution of $f$ with respect to the (infinite) invariant measure, the Lebesgue measure on $\mathbb{R}^{n}$ (see Remark 1.2), hence $\|f\|_{p}=\left\|M_{T}\right\|_{p}$. In this simple case, one deduces Proposition 2.2 from Remark 1.8 about the Burkholder-Gundy inequalities, and it can be easily generalized, first to compactly supported continuous functions $a$. Using Remark 1.9, we find in this way that

$$
\begin{equation*}
\lambda_{p} \leqslant \kappa p^{*}, \quad p^{*}:=\max (p, p /(p-1)), \quad 1<p<+\infty . \tag{2.20}
\end{equation*}
$$

### 2.3. Riesz transforms

In dimension 1 , there is only one Riesz transform $R$, which is called the Hilbert transform $H$. It is defined for $f \in L^{2}(\mathbb{R})$ by

$$
\forall \xi \in \mathbb{R}, \quad(R f)^{\wedge}(\xi)=(H f)^{\wedge}(\xi)=-\frac{\mathrm{i} \xi}{|\xi|} \widehat{f}(\xi)
$$

This is given by a multiplier of constant modulus 1 (almost everywhere), thus the transformation is isometric and invertible on $L^{2}(\mathbb{R})$ by Parseval, and $H$ is a unitary operator on $L^{2}(\mathbb{R})$ with inverse $H^{-1}=-H$. If $\widetilde{u}(x, t)$ denotes the harmonic extension of $H f$ to the upper half-plane, then $u(x, t)+\mathrm{i} \widetilde{u}(x, t)$ is a holomorphic function of the complex variable $z=x+\mathrm{i} t$, because its Fourier transform vanishes for $\xi<0$, implying by inverse Fourier transform that $u(x, t)$ is an integral in $\xi>0$ of the holomorphic functions $\mathrm{e}^{-2 \pi|\xi| t} \mathrm{e}^{2 \mathrm{i} \pi \xi x}=$ $\mathrm{e}^{2 \mathrm{i} \pi \xi(x+\mathrm{i} t)}$. A classical theorem going back to Marcel Riesz [65] states that the Hilbert transform is bounded on $L^{p}(\mathbb{R})$ when $1<p<+\infty$. This is also a consequence of the results on the Littlewood-Paley function $g(f)$, or of martingale inequalities as we shall see below. Some of the first deep connections between Brownian motion and classical Harmonic Analysis can be found in Burkholder-Gundy-Silverstein [18].

The Brownian argument is easier for the Hilbert transform $H_{\mathbb{T}}$ on the unit circle $\mathbb{T} \subset \mathbb{R}^{2}$. Let $\left(B_{t}\right)_{t \geqslant 0}$ be a plane Brownian motion defined on some $(\Omega, \mathcal{F}, P)$, starting from 0 in $\mathbb{R}^{2}$, and let $\tau$ be the first time $t$ when $B_{t}$ hits the circle $\mathbb{T}$. By rotational invariance, the distribution of $B_{\tau}$ is the uniform probability measure on the circle. Let $f$ be a function in
$L^{p}(\mathbb{T})$ and let $u$ be its harmonic extension to the unit disk. Assume that $2 \pi u(0)=\int_{0}^{2 \pi} f(\cos \theta, \sin \theta) \mathrm{d} \theta=0$, and denote by $a \wedge b$ the minimum of $a$ and $b$ real. The random process $\left(M_{t}\right)_{t \geqslant 0}=\left(u\left(B_{t \wedge \tau}\right)\right)_{t \geqslant 0}$ is a Brownian martingale, which can be expressed by the Itō integral

$$
u\left(B_{t \wedge \tau}\right)=\int_{0}^{t \wedge \tau} \nabla u\left(B_{s}\right) \cdot \mathrm{d} B_{s}
$$

Suppose that $1<p<+\infty$. By the continuous version of the BurkholderGundy inequalities, the norm $\|f\|_{L^{p}(\mathbb{T})}=\left\|u\left(B_{\tau}\right)\right\|_{L^{p}(\Omega, \mathcal{F}, P)}$ is equivalent to the norm in $L^{p}(\Omega, \mathcal{F}, P)$ of the square function of the martingale $\left(M_{t}\right)_{t \geqslant 0}$, given by

$$
S(f):=\left(\int_{0}^{\tau}\left|\nabla u\left(B_{s}\right)\right|^{2} \mathrm{~d} s\right)^{1 / 2}
$$

If $\widetilde{f}=H_{\mathbb{T}} f$ denotes the function on $\mathbb{T}$ conjugate to $f$ and $\widetilde{u}$ its harmonic extension to the unit disk, then $|\nabla \widetilde{u}(x)|=|\nabla u(x)|$ for $x$ in the unit disk, according to the Cauchy-Riemann equations for the function $u+\mathrm{i} \widetilde{u}$ holomorphic in the disk. It follows that $S(\widetilde{f})=S(f)$ and the $L^{p}$-boundedness of the Hilbert transform for the circle is established via the BurkholderGundy inequalities of Theorem 1.6. The bound for the norm of $H_{\mathbb{T}}$ obtained in this manner is related to the constants in Burkholder-Gundy. The exact value of the $L^{p}$ norm of $H$ is known, this is due to Pichorides [61], see Remark 2.3 below.

In dimension $n$, there are $n$ Riesz transforms $R_{j}$, defined on $L^{2}\left(\mathbb{R}^{n}\right)$ by

$$
\left(R_{j} f\right)^{\wedge}(\xi)=-\frac{\mathrm{i} \xi_{j}}{|\xi|} \widehat{f}(\xi), \quad j=1, \ldots, n
$$

Since $\left\|\left(\sum_{j=1}^{n}\left|R_{j} f\right|^{2}\right)^{1 / 2}\right\|_{2}^{2}=\sum_{j=1}^{n}\left\|R_{j} f\right\|_{2}^{2}$, one has by Parseval that

$$
\begin{equation*}
\left\|\left(\sum_{j=1}^{n}\left|R_{j} f\right|^{2}\right)^{1 / 2}\right\|_{2}=\|f\|_{2} . \tag{2.21}
\end{equation*}
$$

The Riesz transforms are "collectively bounded" in $L^{p}\left(\mathbb{R}^{n}\right)$, by a constant $\rho_{p}$ independent of the dimension $n$ (Stein [76]), meaning that

$$
\begin{equation*}
\left\|\left(\sum_{j=1}^{n}\left|R_{j} f\right|^{2}\right)^{1 / 2}\right\|_{p} \leqslant \rho_{p}\|f\|_{p}, \quad 1<p<+\infty \tag{2.22}
\end{equation*}
$$

Duoandikoetxea and Rubio de Francia [30] have connected in a few lines this inequality to the properties of the Hilbert transform (see also Pisier [63]).

Proof. - For each nonzero vector $u$ in $\mathbb{R}^{n}$, let us introduce on $L^{2}\left(\mathbb{R}^{n}\right)$ the Hilbert transform $H_{u}$ in the direction $u$ by setting

$$
\forall \xi \in \mathbb{R}^{n}, \quad\left(H_{u} f\right)^{\wedge}(\xi)=-\frac{\mathrm{i} u \cdot \xi}{|u \cdot \xi|} \widehat{f}(\xi)=-\mathrm{i} \operatorname{sign}(u \cdot \xi) \widehat{f}(\xi)
$$

We deduce easily from the one-dimensional case that $H_{u}$ acts on $L^{p}\left(\mathbb{R}^{n}\right)$, with the same norm as that of $H$ on $L^{p}(\mathbb{R})$. It is enough to check the case when $u$ is the first basis vector $\mathbf{e}_{1}$; if one writes the points $x$ in $\mathbb{R}^{n}$ as $x=(t, y), t \in \mathbb{R}, y \in \mathbb{R}^{n-1}$, and if for $f$ belonging to the Schwartz class $\mathcal{S}\left(\mathbb{R}^{n}\right)$ we set $f_{y}(t)=f(t, y)$, we can see that $\left(H_{\mathbf{e}_{1}} f\right)(t, y)=\left(H f_{y}\right)(t)$. Then, applying Fubini's theorem, we obtain

$$
\begin{aligned}
\iint\left|\left(H_{\mathbf{e}_{1}} f\right)(t, y)\right|^{p} \mathrm{~d} t \mathrm{~d} y=\int_{\mathbb{R}^{n-1}}\left(\int_{\mathbb{R}}\left|\left(H f_{y}\right)(t)\right|^{p} \mathrm{~d} t\right) \mathrm{d} y \\
\leqslant\|H\|_{p \rightarrow p}^{p} \int_{\mathbb{R}^{n-1}}\left(\int_{\mathbb{R}}\left|f_{y}(t)\right|^{p} \mathrm{~d} t\right) \mathrm{d} y=\|H\|_{p \rightarrow p}^{p}\|f\|_{p}^{p}
\end{aligned}
$$

We can consider that $\mathcal{R} f=\left(R_{1} f, \ldots, R_{n} f\right)$ is the operator associated to the vector-valued multiplier

$$
\mathbf{m}(\xi)=-\mathrm{i}|\xi|^{-1} \xi \in \mathbb{R}^{n}, \quad \xi \in \mathbb{R}^{n}
$$

that is to say, the operator sending $f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$ to the function $T_{\mathbf{m}} f$ from $\mathbb{R}^{n}$ to $\mathbb{R}^{n}$ whose $\mathbb{R}^{n}$-valued Fourier transform is equal to $\widehat{f}(\xi) \mathbf{m}(\xi)$. For $f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$, let us look at the vector-valued integral

$$
(\mathcal{H} f)(x)=\int_{\mathbb{R}^{n}}\left(H_{u} f\right)(x) u \mathrm{~d} \gamma_{n}(u) \in \mathbb{R}^{n}, \quad x \in \mathbb{R}^{n}
$$

where $\gamma_{n}$ is the Gaussian probability measure from (1.17). The operator $\mathcal{H}$ corresponds to the vector-valued multiplier defined when $\xi \neq 0$ by

$$
-\mathrm{i} \int_{\mathbb{R}^{n}} \operatorname{sign}(u \cdot \xi) u \mathrm{~d} \gamma_{n}(u)=-\mathrm{i}\left(\int_{\mathbb{R}}|v| \mathrm{d} \gamma_{1}(v)\right)|\xi|^{-1} \xi=-\mathrm{i} \sqrt{\frac{2}{\pi}}|\xi|^{-1} \xi
$$

This can be seen by integrating on affine hyperplanes orthogonal to $\xi$. The "normalized" partial integral on the hyperplane $\xi^{\perp}+v|\xi|^{-1} \xi, v \in \mathbb{R}$, is equal to

$$
\int_{\xi^{\perp}} \operatorname{sign}(v)\left(w+v|\xi|^{-1} \xi\right) \mathrm{d} \gamma_{\xi^{\perp}}(w)=|v||\xi|^{-1} \xi
$$

It follows that $\mathcal{R} f=\sqrt{\pi / 2} \mathcal{H} f$. For $x$ fixed, the norm of $(\mathcal{H} f)(x)$ is the supremum of scalar products with vectors $\theta \in S^{n-1}$, and letting $1 / q+1 / p=1$, one has that

$$
\begin{aligned}
|(\mathcal{H} f)(x)| & \leqslant \sup _{\theta \in S^{n-1}} \int_{\mathbb{R}^{n}}|\theta \cdot u|\left|\left(H_{u} f\right)(x)\right| \mathrm{d} \gamma_{n}(u) \\
& \leqslant\left(\int_{\mathbb{R}}|v|^{q} \mathrm{~d} \gamma_{1}(v)\right)^{1 / q}\left(\int_{\mathbb{R}^{n}}\left|\left(H_{u} f\right)(x)\right|^{p} \mathrm{~d} \gamma_{n}(u)\right)^{1 / p}
\end{aligned}
$$

Using the notation $g_{q}$ of (1.18) for the Gaussian moments, we get

$$
\begin{align*}
\left(\int_{\mathbb{R}^{n}}|(\mathcal{R} f)(x)|^{p} \mathrm{~d} x\right)^{1 / p} & =\sqrt{\frac{\pi}{2}}\left(\int_{\mathbb{R}^{n}}|(\mathcal{H} f)(x)|^{p} \mathrm{~d} x\right)^{1 / p} \\
& \leqslant \sqrt{\frac{\pi}{2}} g_{q}\left(\int_{\mathbb{R}^{n}} \int_{\mathbb{R}^{n}}\left|\left(H_{u} f\right)(x)\right|^{p} \mathrm{~d} \gamma_{n}(u) \mathrm{d} x\right)^{1 / p} \\
& \leqslant \sqrt{\frac{\pi}{2}} g_{q}\|H\|_{p \rightarrow p}\|f\|_{p} \tag{2.23}
\end{align*}
$$

This argument yields $\rho_{p} \leqslant \sqrt{\pi / 2} g_{q}\|H\|_{p \rightarrow p}$ for the constant $\rho_{p}$ in (2.22). When $p=2$, this gives $\rho_{2} \leqslant \sqrt{\pi / 2}$ instead of the correct value $\rho_{2}=1$ of (2.21). When $p$ tends to 1 , we obtain by (1.18) that

$$
\left(\int_{\mathbb{R}^{n}}|(\mathcal{R} f)(x)|^{p} \mathrm{~d} x\right)^{1 / p} \leqslant \kappa \sqrt{q}\|H\|_{p \rightarrow p}\|f\|_{p}
$$

Remark 2.3. - The value $g_{q}=\left(\int_{\mathbb{R}}|v|^{q} \mathrm{~d} \gamma_{1}(v)\right)^{1 / q}$ tends to $\sqrt{2 / \pi}$ when $p$ tends to $+\infty$, and the asymptotic result $\rho_{p} \simeq\|H\|_{p \rightarrow p}$ obtained from (2.23) in this case is essentially best possible. Indeed, Iwaniec and Martin [47] have shown that the operator norm on $L^{p}\left(\mathbb{R}^{n}\right)$ of each individual Riesz transform $R_{j}, j=1, \ldots, n$, is equal to the one of the Hilbert transform $H$ on $L^{p}(\mathbb{R})$, hence $1 \leqslant\|H\|_{p \rightarrow p} \leqslant \rho_{p}$. According to Pichorides [61], the norm of the Hilbert transform is given by

$$
\|H\|_{p \rightarrow p}=\cot \left(\frac{\pi}{2 p^{*}}\right), \text { with } p^{*}=\max (p, p /(p-1))
$$

Iwaniec and Martin [47] also bound the "collective" norm in (2.22) by $\sqrt{2} H_{p}(1)$, where $H_{p}(1)$ is the norm on $L^{p}(\mathbb{C}) \simeq L^{p}\left(\mathbb{R}^{2}\right)$ of the "complex Hilbert transform", which corresponds to the multiplier $\mathbb{C} \ni \xi \mapsto \mathrm{i}|\xi|^{-1} \xi$, in other words, the operator $R_{1}+\mathrm{i} R_{2}$ on $L^{p}\left(\mathbb{R}^{2}\right)$. Iwaniec and Sbordone [48, Appendix] add a few lines and give $H_{p}(1) \leqslant \frac{\pi}{2}\|H\|_{p \rightarrow p}$ so that finally

$$
\begin{equation*}
\rho_{p} \leqslant \sqrt{2} H_{p}(1) \leqslant \frac{\pi}{\sqrt{2}}\|H\|_{p \rightarrow p} \leqslant \kappa p^{*} \tag{2.24}
\end{equation*}
$$

Remark. - The proof from [30] is in the spirit of the method of rotations, which uses integration in polar coordinates to get directional operators in its radial part, see also Section 4.1. With this method, one can relate to the Hilbert transform not only the Riesz transforms, but also more general singular integrals with odd kernel, see [39, Section 5.2] for example.

## 3. Analytic tools

### 3.1. Some known facts about the Gamma function

From Euler's formula

$$
\forall z \in \mathbb{C} \backslash(-\mathbb{N}), \quad \Gamma(z)=\lim _{n \rightarrow \infty} \frac{n!n^{z}}{z(z+1) \ldots(z+n)}
$$

one passes to the Weierstrass infinite product for $1 / \Gamma$, stating that

$$
\frac{1}{\Gamma(z+1)}=\frac{1}{z \Gamma(z)}=\mathrm{e}^{\gamma z} \prod_{n=1}^{\infty}\left((1+z / n) \mathrm{e}^{-z / n}\right)
$$

where $\gamma$ is the Euler-Mascheroni constant. It follows that $1 / \Gamma$ is an entire function, with simple zeroes $z=0,-1,-2, \ldots$. For the interpolation arguments to come, we need upper estimates on the modulus of $1 / \Gamma(\sigma+\mathrm{i} \tau)$ for $\sigma, \tau$ real. From the preceding formula and from $\Gamma(\bar{z})=\overline{\Gamma(z)}$, we infer that

$$
\begin{equation*}
\left|\frac{1}{\Gamma(1+\mathrm{i} \tau)}\right|^{2}=\prod_{n=1}^{\infty}\left(1+\frac{\tau^{2}}{n^{2}}\right)=\frac{\sinh (\pi \tau)}{\pi \tau}, \quad \tau \in \mathbb{R} \tag{3.1}
\end{equation*}
$$

according to another result due to Euler, the famous formula

$$
\begin{equation*}
\frac{\sin (\pi z)}{\pi z}=\prod_{n=1}^{\infty}\left(1-\frac{z^{2}}{n^{2}}\right) \tag{3.2.E}
\end{equation*}
$$

The connoisseur has seen that we just came upon a special case of the "Euler reflection formula", stating that $\Gamma(z)^{-1} \Gamma(1-z)^{-1}=\sin (\pi z) / \pi$ for every $z \in \mathbb{C}$, or equivalently $\Gamma(1+z)^{-1} \Gamma(1-z)^{-1}=\sin (\pi z) /(\pi z)$. For every $x$ real, one has

$$
\begin{equation*}
\frac{\sinh (\pi x)}{\pi x} \leqslant \frac{\mathrm{e}^{\pi|x|}}{1+\pi|x|} \leqslant \frac{\mathrm{e}^{\pi|x|}}{\left(1+x^{2}\right)^{1 / 2}} \tag{3.3}
\end{equation*}
$$

The right-hand inequality is evident, the left-hand one is equivalent to saying that for every $y \geqslant 0$, we have $(1+y) \sinh (y) \leqslant y \mathrm{e}^{y}$ or $h(y):=(y-1) \mathrm{e}^{2 y}+y+$ $1 \geqslant 0$, which is true because $h(0)=h^{\prime}(0)=0$ and $h^{\prime \prime}(y)=4 y \mathrm{e}^{2 y} \geqslant 0$ when $y \geqslant 0$. Using (3.1) and (3.3), we get in particular that

$$
\begin{equation*}
\forall \tau \in \mathbb{R}, \quad\left|\frac{1}{\Gamma(1+\mathrm{i} \tau)}\right| \leqslant\left(\sqrt{1+\tau^{2}}\right)^{-1 / 2} \mathrm{e}^{\pi|\tau| / 2} \tag{3.4}
\end{equation*}
$$

More generally than in (3.1), for $\sigma \in[0,1]$, let us write

$$
\begin{aligned}
\left|\frac{1}{\Gamma(1+\sigma+\mathrm{i} \tau)}\right|^{2} & =\mathrm{e}^{2 \gamma \sigma} \prod_{n \geqslant 1}\left(\left[(1+\sigma / n)^{2}+(\tau / n)^{2}\right] \mathrm{e}^{-2 \sigma / n}\right) \\
& =\mathrm{e}^{2 \gamma \sigma} \prod_{n \geqslant 1}\left((1+\sigma / n)^{2} \mathrm{e}^{-2 \sigma / n}\right) \prod_{n \geqslant 1}\left(1+\left(\frac{\tau}{n+\sigma}\right)^{2}\right) \\
& =\Gamma(1+\sigma)^{-2} \prod_{n \geqslant 1}\left(1+\left(\frac{\tau}{n+\sigma}\right)^{2}\right)
\end{aligned}
$$

We have by convexity of $\ln \left(1+x^{-2}\right)$ for $x>0$ that

$$
\begin{aligned}
\prod_{n \geqslant 1}\left(1+\left(\frac{\tau}{n+\sigma}\right)^{2}\right) & \leqslant\left(\prod_{n \geqslant 1}\left(1+\left(\frac{\tau}{n}\right)^{2}\right)\right)^{1-\sigma}\left(\prod_{n \geqslant 1}\left(1+\left(\frac{\tau}{n+1}\right)^{2}\right)\right)^{\sigma} \\
& =\left(1+\tau^{2}\right)^{-\sigma} \prod_{n \geqslant 1}\left(1+\left(\frac{\tau}{n}\right)^{2}\right) \\
& =\left(1+\tau^{2}\right)^{-\sigma} \frac{\sinh (\pi \tau)}{\pi \tau}
\end{aligned}
$$

It follows that

$$
\left|\frac{1}{\Gamma(1+\sigma+\mathrm{i} \tau)}\right|^{2} \leqslant \Gamma(1+\sigma)^{-2}\left(1+\tau^{2}\right)^{-\sigma} \frac{\sinh (\pi \tau)}{\pi \tau}
$$

and applying (3.3) we obtain

$$
\begin{equation*}
\left|\frac{1}{\Gamma(1+\sigma+\mathrm{i} \tau)}\right| \leqslant \Gamma(1+\sigma)^{-1}\left(\sqrt{1+\tau^{2}}\right)^{1 / 2-1-\sigma} \mathrm{e}^{\pi|\tau| / 2} \tag{3.5}
\end{equation*}
$$

We extend this bound by using the functional equation $z \Gamma(z)=\Gamma(z+1)$.
When $z=k+1+\sigma+\mathrm{i} \tau$, with $\sigma \in(0,1)$ and $k \geqslant 1$ an integer, we have

$$
\begin{align*}
\left|\frac{1}{\Gamma(k+1+\sigma+\mathrm{i} \tau)}\right| & =\left(\prod_{j=1}^{k}\left((j+\sigma)^{2}+\tau^{2}\right)^{-1 / 2}\right)\left|\frac{1}{\Gamma(1+\sigma+\mathrm{i} \tau)}\right| \\
& \leqslant\left(\sqrt{1+\tau^{2}}\right)^{-k}\left|\frac{1}{\Gamma(1+\sigma+\mathrm{i} \tau)}\right| \tag{3.6}
\end{align*}
$$

Letting $a \wedge b=\min (a, b)$ for $a, b \in \mathbb{R}$, we see that

$$
\begin{aligned}
& \Gamma(1+\sigma)=\int_{0}^{+\infty} u^{\sigma} \mathrm{e}^{-u} \mathrm{~d} u \\
& \quad \geqslant \int_{0}^{+\infty}(u \wedge 1) \mathrm{e}^{-u} \mathrm{~d} u=\int_{0}^{1} \mathrm{e}^{-u} \mathrm{~d} u=1-\mathrm{e}^{-1}>\frac{1}{2} \\
& \quad-49-
\end{aligned}
$$

Let us mention that the actual minimal value of $\Gamma$ on $(0,+\infty)$ is reached at

$$
\begin{equation*}
x_{\Gamma}=1.46163 \ldots \text { and that } \Gamma\left(x_{\Gamma}\right)>0.88 . \tag{3.7}
\end{equation*}
$$

Note that on $(0,+\infty)$, the function $x \mapsto \ln \Gamma(x)$ is convex and $\ln \Gamma(1)=$ $\ln \Gamma(2)=0$, hence $\Gamma(x) \leqslant 1$ when $x \in[1,2]$ and $\Gamma(x) \geqslant 1$ on $(0,1]$ and $[2,+\infty)$.

We get consequently by (3.5) and (3.6) that

$$
\begin{equation*}
\left|\frac{1}{\Gamma(k+1+\sigma+\mathrm{i} \tau)}\right| \leqslant 2\left(\sqrt{1+\tau^{2}}\right)^{1 / 2-k-1-\sigma} \mathrm{e}^{\pi|\tau| / 2} . \tag{3.8}
\end{equation*}
$$

When $z=-k+\sigma+\mathrm{i} \tau$, with $k \geqslant 0$ an integer, we obtain by the functional equation

$$
\left|\frac{1}{\Gamma(z)}\right|=\left|\frac{1}{\Gamma(1+\sigma+\mathrm{i} \tau)}\right| \prod_{j=-k}^{0}\left((j+\sigma)^{2}+\tau^{2}\right)^{1 / 2}
$$

For $j=0,-1$, the factors in the product are $\leqslant\left(1+\tau^{2}\right)^{1 / 2}$, thus

$$
\begin{equation*}
\left|\frac{1}{\Gamma(-k+\sigma+\mathrm{i} \tau)}\right| \leqslant\left(1+\tau^{2}\right)^{(k+1) / 2}\left|\frac{1}{\Gamma(1+\sigma+\mathrm{i} \tau)}\right| \quad \text { when } k=0,1 \tag{3.9}
\end{equation*}
$$

and when $j \leqslant-2$, we have $\left((j+\sigma)^{2}+\tau^{2}\right)^{1 / 2} \leqslant(|j|-\sigma)\left(1+\tau^{2}\right)^{1 / 2}$. It follows for $z=-k+\sigma+\mathrm{i} \tau, k \geqslant 2$, that

$$
\begin{equation*}
\left|\frac{1}{\Gamma(z)}\right| \leqslant(k-\sigma)(k-1-\sigma) \ldots(2-\sigma)\left(1+\tau^{2}\right)^{(k+1) / 2}\left|\frac{1}{\Gamma(1+\sigma+\mathrm{i} \tau)}\right| \tag{3.10}
\end{equation*}
$$

By the functional equation and the convexity of $\ln \Gamma$ on $(0,+\infty)$, we have

$$
\begin{aligned}
& \Gamma(1+\sigma)^{-1}(k-\sigma)(k-1-\sigma) \ldots(2-\sigma)=\frac{\Gamma(k+1-\sigma)}{\Gamma(2-\sigma) \Gamma(1+\sigma)} \\
& \leqslant \frac{\Gamma(k+1-\sigma)}{\Gamma(3 / 2)}=\frac{2}{\sqrt{\pi}} \Gamma(k+1-\sigma)<2 \Gamma(k+1-\sigma)
\end{aligned}
$$

Coming back to (3.10) and using (3.5), we conclude when $k \geqslant 2$ that

$$
\begin{equation*}
\left|\frac{1}{\Gamma(-k+\sigma+\mathrm{i} \tau)}\right| \leqslant 2 \Gamma(k-\sigma+1)\left(\sqrt{1+\tau^{2}}\right)^{1 / 2+k-\sigma} \mathrm{e}^{\pi|\tau| / 2} \tag{3.11}
\end{equation*}
$$

When $\operatorname{Re} z \geqslant-1$, it follows from (3.8) and (3.9) that

$$
\left|\frac{1}{\Gamma(z)}\right| \leqslant 2\left(\sqrt{1+(\operatorname{Im} z)^{2}}\right)^{1 / 2-\operatorname{Re} z} \mathrm{e}^{\pi|\operatorname{Im} z| / 2}
$$

so, in every half-plane of the form $\operatorname{Re} z \geqslant a$, one has by (3.11) an upper bound

$$
\begin{equation*}
\left|\frac{1}{\Gamma(z)}\right| \leqslant \beta_{a}\left(\sqrt{1+|\operatorname{Im} z|^{2}}\right)^{1 / 2-a} \mathrm{e}^{\pi|\operatorname{Im} z| / 2} \tag{3.12.Г}
\end{equation*}
$$

with $\beta_{a}=2 \Gamma(|a|+1)$ when $a \leqslant-1$, and $\beta_{a}=2$ otherwise.

Remark. - The rather crude estimate (3.12. Г) is sufficient for our purposes. In [73], Stein refers to Titchmarsh [82, p. 259], for an exact asymptotic estimate. When $\sigma$ is fixed and $|\tau| \rightarrow+\infty$, one has

$$
|\Gamma(\sigma+\mathrm{i} \tau)| \simeq \sqrt{2 \pi} \mathrm{e}^{-\pi|\tau| / 2}|\tau|^{\sigma-1 / 2}
$$

When $\sigma \geqslant 1$, the preceding proof gives a lower bound $2^{-1} \sqrt{2 \pi} \mathrm{e}^{-\pi|\tau| / 2}|\tau|^{\sigma-1 / 2}$ for every $\tau$. We can see it by replacing the inequality (3.3) with the evident inequality $\sinh (\pi x) /(\pi x) \leqslant(2 \pi|x|)^{-1} \mathrm{e}^{\pi|x|}$. It is not possible to replace $\sqrt{1+|\operatorname{Im} z|^{2}}$ by $|\operatorname{Im} z|$ in $(3.12 . \Gamma)$ when $\operatorname{Re} z \leqslant-1$, because the zeroes $-1,-2, \ldots$ of $1 / \Gamma$ are simple. For more results on the Gamma function, we refer to Andrews-Askey-Roy [2].

### 3.2. The interpolation scheme

We begin with the classical three lines lemma, an easier version of which is the Hadamard three-circle theorem. After this, we shall turn to interpolation of holomorphic families of linear operators.

### 3.2.1. The three lines lemma

Lemma 3.1. - Let $S$ denote the open strip $\{z: 0<\operatorname{Re} z<1\}$ in the complex plane. Let $f$ be a function holomorphic in $S$ and continuous on the closure of $S$. Assume that $f$ is bounded in $S$ and that

$$
|f(0+\mathrm{i} \tau)| \leqslant C_{0}, \quad|f(1+\mathrm{i} \tau)| \leqslant C_{1}
$$

for all $\tau \in \mathbb{R}$. Then, for every $\theta \in(0,1)$, one has that $|f(\theta)| \leqslant C_{0}^{1-\theta} C_{1}^{\theta}$.
Remark 3.2. - Of course $f(\theta+\mathrm{i} \tau)$ admits the same bound for every $\tau \in \mathbb{R}$, by translating $f$ vertically. The somewhat strange assumption that $f$ must be bounded on the whole strip by a value which does not appear in the final result is not the finest assumption that makes the conclusion valid, see a better criterion below. However, when Lemma 3.1 applies, the function $f$ is bounded at last. It is well known that some restriction on the size of $f$ inside the strip is needed for the lemma to hold true. Indeed, in the strip $S_{\pi}=\{z:|\operatorname{Re} z| \leqslant \pi / 2\}$, the function $f(z)=\mathrm{e}^{\cos z}$ has modulus one on the two lines $\operatorname{Re} z= \pm \pi / 2$, but it is "very big" when $\operatorname{Re} z=0$, since $|f(\mathrm{i} \tau)|=\mathrm{e}^{\cosh (\tau)}$. For a function $f$ holomorphic in an open vertical strip $S$, continuous on the closure and bounded by 1 on the two boundary lines, either $f$ is bounded by 1 on $S$, or else $\sup _{|\operatorname{Im} z|=|\tau|}|f(z)|$ must become extremely large when $|\tau|$ tends to infinity. This is the typical situation with the theorems of Phragmén-Lindelöf type, see [69, Chap 12, 12.7] for example.

Here is a sufficient criterion ensuring that $|f|$ is bounded by its supremum on the boundary $\partial S_{w}$ of a vertical strip $S_{w}$ of width $w$. If $f$ is holomorphic on $S_{w}$, continuous on the closure with $|f| \leqslant 1$ on $\partial S_{w}$, and if for some $a<\pi / w$ one has

$$
|f(z)|=O\left(\exp \left(\mathrm{e}^{a|\operatorname{Im} z|}\right)\right)
$$

when $z$ tends to infinity in $S_{w}$, then $|f|$ is bounded by 1 on the strip. Let us prove it assuming $\ln |f(z)| \leqslant \kappa \mathrm{e}^{a|\operatorname{Im} z|}$ in $S_{\pi}=\{|\operatorname{Re} z| \leqslant \pi / 2\}$, for an $a<1=\pi / w$. Set $g_{\varepsilon}(z)=\mathrm{e}^{-\varepsilon \cos (b z)}$, with $\varepsilon>0$ and $a<b<1$. If $z=\sigma+\mathrm{i} \tau$ and $|\sigma| \leqslant \pi / 2$, we have

$$
\begin{aligned}
\left|g_{\varepsilon}(z)\right| & =\exp (-\varepsilon \operatorname{Re} \cos (b z))=\exp (-\varepsilon \cos (b \sigma) \cosh (b \tau)) \\
& \leqslant \exp (-\varepsilon \cos (b \pi / 2) \cosh (b \tau)) \leqslant \exp \left(-B_{\varepsilon} \mathrm{e}^{b|\tau|}\right) \leqslant 1, \quad B_{\varepsilon}>0
\end{aligned}
$$

hence $\left|f(z) g_{\varepsilon}(z)\right| \leqslant 1$ on $\partial S_{\pi}$, and if $|\tau|=|\operatorname{Im} z| \geqslant(b-a)^{-1} \ln \left(\kappa / B_{\varepsilon}\right)$ we get

$$
\begin{equation*}
\ln \left|f(z) g_{\varepsilon}(z)\right| \leqslant \kappa \mathrm{e}^{a|\tau|}-B_{\varepsilon} \mathrm{e}^{b|\tau|} \leqslant 0 . \tag{3.13}
\end{equation*}
$$

Given any $z_{0} \in S_{\pi}$, we can find a rectangle $R_{\varepsilon}=\{|\operatorname{Re} z| \leqslant \pi / 2,|\operatorname{Im} z| \leqslant$ $\left.\tau_{0}(\varepsilon)\right\}$ containing $z_{0}$ such that $\left|f(z) g_{\varepsilon}(z)\right| \leqslant 1$ on $\partial R_{\varepsilon}$. We then have $\left|f\left(z_{0}\right) g_{\varepsilon}\left(z_{0}\right)\right| \leqslant 1$ by the maximum principle, $\left|f\left(z_{0}\right)\right| \leqslant\left|\mathrm{e}^{\varepsilon \cos \left(b z_{0}\right)}\right|$ for every $\varepsilon>0$, thus $\left|f\left(z_{0}\right)\right| \leqslant 1$.

Several times later on, we encounter situations where the function $f$ is not bounded on the two lines limiting a vertical strip $S$, but has instead a growth exponential in $|\tau|=|\operatorname{Im} z|$. The next lemma generalizes the preceding. Our proof and estimate are not the "correct" ones, as we shall explain below after Corollary 3.4, but they give a reasonable explicit bound. In these Notes, we shall say that a function $f$ defined on a vertical strip $S$ has an admissible growth in the strip if for some $\kappa>0$, the function $f$ admits in $S$ a bound of the form $|f(z)| \leqslant \kappa \mathrm{e}^{\kappa|\operatorname{Im} z|}$.

Lemma 3.3. - Let $f$ be a function holomorphic in the strip $S=\{z$ : $0<\operatorname{Re} z<1\}$, with admissible growth in $S$ and continuous on the closure of $S$. Assume that there exist real numbers $a_{0}, a_{1} \geqslant 0$ and $b_{0}, b_{1}$ such that for every $\tau \in \mathbb{R}$, one has

$$
|f(0+\mathrm{i} \tau)| \leqslant \mathrm{e}^{a_{0}|\tau|+b_{0}}, \quad|f(1+\mathrm{i} \tau)| \leqslant \mathrm{e}^{a_{1}|\tau|+b_{1}}
$$

For every $\theta \in(0,1)$, it follows that

$$
|f(\theta)| \leqslant \exp \left(\sqrt{\theta(1-\theta)} \sqrt{(1-\theta) a_{0}^{2}+\theta a_{1}^{2}}+(1-\theta) b_{0}+\theta b_{1}\right)
$$

Proof. - We introduce the holomorphic function $g(z):=\mathrm{e}^{c z^{2} / 2+d z}$, with $c>0$ and $d$ real. If $z=\sigma+\mathrm{i} \tau$, we see that $|g(z)|=\mathrm{e}^{c\left(\sigma^{2}-\tau^{2}\right) / 2+d \sigma}$. On the vertical side $\operatorname{Re} z=0$ of $S$, we have that

$$
|f(\mathrm{i} \tau) g(\mathrm{i} \tau)| \leqslant \mathrm{e}^{a_{0}|\tau|+b_{0}-c \tau^{2} / 2} \leqslant \mathrm{e}^{a_{0}^{2} /(2 c)+b_{0}}=: E_{0}
$$

and when $\operatorname{Re} z=1$, we get the upper bound

$$
|f(1+\mathrm{i} \tau) g(1+\mathrm{i} \tau)| \leqslant \mathrm{e}^{a_{1}|\tau|+b_{1}-c \tau^{2} / 2+c / 2+d} \leqslant \mathrm{e}^{a_{1}^{2} /(2 c)+b_{1}+c / 2+d}=: E_{1}
$$

We choose $d$ so that $E_{0}=E_{1}$, and we need not mention the value of $d$.
It follows from the assumption $|f(z)| \leqslant \kappa \mathrm{e}^{\kappa|\operatorname{Im} z|}=\kappa \mathrm{e}^{\kappa|\tau|}$ that $f(z) g(z)$ tends to zero at infinity in $S$. Let us fix $\theta \in(0,1)$. If $f(\theta) \neq 0$, there exists $\tau_{0}>0$ such that $|f(z) g(z)|<|f(\theta) g(\theta)|$ when $|\operatorname{Im} z| \geqslant \tau_{0}$. By the maximum principle for the compact rectangle $R=\left\{0 \leqslant \operatorname{Re} z \leqslant 1,|\operatorname{Im} z| \leqslant \tau_{0}\right\}$, we know that the maximum of $|f(z) g(z)|$ is reached at the boundary of $R$, but it cannot be on the horizontal sides $|\operatorname{Im} z|=\tau_{0}$. Hence $|f(\theta) g(\theta)| \leqslant E_{0}=$ $E_{1}=E_{0}^{1-\theta} E_{1}^{\theta}$, we get therefore

$$
\begin{aligned}
& |f(\theta)| \leqslant \mathrm{e}^{-c \theta^{2} / 2-d \theta} E_{0}^{1-\theta} E_{1}^{\theta} \\
& \quad=\exp \left(\frac{(1-\theta) a_{0}^{2}+\theta a_{1}^{2}}{2 c}+(1-\theta) b_{0}+\theta b_{1}+c \theta(1-\theta) / 2\right)
\end{aligned}
$$

and after optimizing in $c>0$, we conclude that

$$
|f(\theta)| \leqslant \exp \left(\sqrt{(1-\theta) a_{0}^{2}+\theta a_{1}^{2}} \sqrt{\theta(1-\theta)}+(1-\theta) b_{0}+\theta b_{1}\right)
$$

Corollary 3.4. - Let $f$ be a function holomorphic in $S=\left\{z: \alpha_{0}<\right.$ $\left.\operatorname{Re} z<\alpha_{1}\right\}$, with admissible growth in the strip $S$ and continuous on the closure of $S$. Assume that there exist real numbers $u_{0}, u_{1} \geqslant 0$ and $v_{0}, v_{1}$ such that

$$
\left|f\left(\alpha_{0}+\mathrm{i} \tau\right)\right| \leqslant \mathrm{e}^{u_{0}|\tau|+v_{0}}, \quad\left|f\left(\alpha_{1}+\mathrm{i} \tau\right)\right| \leqslant \mathrm{e}^{u_{1}|\tau|+v_{1}}
$$

for every $\tau \in \mathbb{R}$. Let $\theta \in[0,1]$, set $\alpha_{\theta}=(1-\theta) \alpha_{0}+\theta \alpha_{1}, u_{\theta}=(1-\theta) u_{0}+\theta u_{1}$ and $v_{\theta}=(1-\theta) v_{0}+\theta v_{1}$. For every $\tau \in \mathbb{R}$, one has

$$
\left|f\left(\alpha_{\theta}+\mathrm{i} \tau\right)\right| \leqslant E_{w, \theta}\left(u_{0}, u_{1}\right) \mathrm{e}^{u_{\theta}|\tau|+v_{\theta}}
$$

where $w=\alpha_{1}-\alpha_{0}$ denotes the width of the strip $S$ and where

$$
E_{w, \theta}\left(u_{0}, u_{1}\right):=\exp \left(w \sqrt{\theta(1-\theta)} \sqrt{(1-\theta) u_{0}^{2}+\theta u_{1}^{2}}\right)
$$

Notice that $\sqrt{\theta(1-\theta)} \leqslant 1 / 2$ for every $\theta \in[0,1]$. When $0 \leqslant u_{0}, u_{1} \leqslant u$, one can always employ the simpler bound $E_{w, \theta}(u, u) \leqslant \mathrm{e}^{w u / 2}$.

Proof. - We begin with $S_{1}:=\{0<\operatorname{Re} z<1\}$. We bound the modulus of $f\left(\theta+\mathrm{i} \tau_{0}\right)$ for $\tau_{0}$ in $\mathbb{R}$ by performing a vertical translation of $f$, then invoking Lemma 3.3. The function $F(z)=f\left(z+\mathrm{i} \tau_{0}\right)$ satisfies $|F(j+\mathrm{i} \tau)| \leqslant$ $\mathrm{e}^{u_{j}|\tau|+\left(u_{j}\left|\tau_{0}\right|+v_{j}\right)}, j=0,1$, and the bound for $|F(\theta)|$ given at Lemma 3.3 implies that

$$
\begin{equation*}
\left|f\left(\theta+\mathrm{i} \tau_{0}\right)\right| \leqslant E_{1, \theta}\left(u_{0}, u_{1}\right) \mathrm{e}^{u_{\theta}\left|\tau_{0}\right|+v_{\theta}}, \quad \tau_{0} \in \mathbb{R} \tag{3.14}
\end{equation*}
$$

It is easy to pass to $S=\left\{\alpha_{0}<\operatorname{Re} z<\alpha_{1}\right\}$ with the transform that replaces $f(z)$, defined for $z \in S$, by $F(Z)=f\left(\alpha_{0}+Z w\right)$ for $Z \in S_{1}$. If $\left|f\left(\alpha_{j}+\mathrm{i} \tau\right)\right| \leqslant$ $\mathrm{e}^{u_{j}|\tau|+v_{j}}, j=0,1$, then $|F(j+\mathrm{i} \tau)| \leqslant \mathrm{e}^{w u_{j}|\tau|+v_{j}}$ and by (3.14) we have that

$$
\begin{aligned}
& \left|f\left(\alpha_{\theta}+\mathrm{i} \tau_{0}\right)\right|=\left|F\left(\theta+\mathrm{i} \tau_{0} / w\right)\right| \\
& \quad \leqslant E_{1, \theta}\left(w u_{0}, w u_{1}\right) \mathrm{e}^{\left(w u_{\theta}\right)\left|\tau_{0} / w\right|+v_{\theta}}=E_{w, \theta}\left(u_{0}, u_{1}\right) \mathrm{e}^{u_{\theta}\left|\tau_{0}\right|+v_{\theta}}
\end{aligned}
$$

Applying Corollary 3.4 in the case where $u_{0}=u_{1}=u \geqslant 0$ and $v_{j}=0$, one sees that when $f$ has an admissible growth in $S$, the hypothesis $\left|f\left(\alpha_{j}+\mathrm{i} \tau\right)\right| \leqslant$ $\mathrm{e}^{u|\tau|}$ for all $\tau \in \mathbb{R}$ and $j=0,1$, implies $|f(\alpha+\mathrm{i} \tau)| \leqslant \mathrm{e}^{w u / 2} \mathrm{e}^{u|\tau|}$ in the strip. It is not possible to replace the "bounding factor" $\mathrm{e}^{w u / 2}$ by 1 , as we shall understand below.

The "correct" proof of Lemma 3.3 uses a lemma given by Hirschman [45, Lemma 1], cited by Stein [70]. In our case, we consider the function $U$, harmonic in the open strip $S_{1}=\{0<\operatorname{Re} z<1\}$ and continuous on the closed strip, equal to $a_{j}|\tau|+b_{j}$ at each boundary point $j+\mathrm{i} \tau$, with $a_{j} \geqslant 0, \tau \in \mathbb{R}$ and $j=0,1$. Let $V$ be the harmonic conjugate of $U$ in $S_{1}$, defined up to an additive constant by the fact that $\nabla V(z)$, for $z \in S_{1}$, is equal to $R \nabla U(z)$ where $R$ is the rotation of angle $+\pi / 2$ in $\mathbb{R}^{2} \simeq \mathbb{C}$. Let us set $V(1 / 2)=0$ in order to fix $V$ entirely. Since $U$ is harmonic, the 1-form $-U_{y} d x+U_{x} d y$ is closed and $V(z)=\int_{0}^{1} R \nabla U(\gamma(s))$. $\gamma^{\prime}(s) \mathrm{d} s$ for any $C^{1}$ path $\gamma$ in $S_{1}$ such that $\gamma(0)=1 / 2$ and $\gamma(1)=z$. Then $U+\mathrm{i} V$ is holomorphic, by the Cauchy-Riemann equations. Consider the holomorphic outer function

$$
g(z)=\exp (-U(z)-\mathrm{i} V(z)), \quad z \in S_{1}
$$

for which $|g(z)|=\exp (-U(z))$ and $|g(z)| \leqslant \mathrm{e}^{-\left(b_{0} \wedge b_{1}\right)}$ in $S_{1}$. If $f$ is as in Lemma 3.3, then $|f g| \leqslant 1$ at the boundary of $S_{1}$ and $f g$ has an admissible growth. It follows from an easy variation of Lemma 3.1 that $|(f g)(\theta)| \leqslant 1$ thus $|f(\theta)| \leqslant \mathrm{e}^{U(\theta)}$, and it remains to express $U(\theta)$, with the help of the harmonic measure at $\theta$ for $S_{1}$.

We shall obtain the harmonic measures for $S_{\pi}=\{z \in \mathbb{C}:|\operatorname{Re} z|<$ $\pi / 2\}$ from the case of the open unit disk $D$, by a conformal mapping (see also [39, proof of Lemma 1.3.8]). Let $\sigma$ belong to $I_{\pi}=(-\pi / 2, \pi / 2)=S_{\pi} \cap \mathbb{R}$. The Poisson probability measure $\mu_{\sigma}$ at $\sigma$ relative to $S_{\pi}$ can be written as $\mu_{\sigma}=\mu_{\sigma, 0}+\mu_{\sigma, 1}$, where $\mu_{\sigma, 0}$ is supported on $B_{0}=-\pi / 2+\mathrm{i} \mathbb{R}$ and $\mu_{\sigma, 1}$ on $B_{1}=\pi / 2+\mathrm{i} \mathbb{R}$. If $h$ is real, harmonic in $S_{\pi}$, bounded and continuous on the closure of $S_{\pi}$, the value of $h$ at $\sigma$ is equal to

$$
\begin{equation*}
h(\sigma)=\int_{\partial S_{\pi}} h \mathrm{~d} \mu_{\sigma}=\int_{B_{0}} h \mathrm{~d} \mu_{\sigma, 0}+\int_{B_{1}} h \mathrm{~d} \mu_{\sigma, 1} \tag{3.15}
\end{equation*}
$$

The Poisson probability measure $\nu_{r}$ for $D$ at $r \in(-1,1)$ has density $g_{r}\left(\mathrm{e}^{\mathrm{i} \beta}\right)=$ $\left(1-r^{2}\right) /\left(1-2 r \cos \beta+r^{2}\right)$ with respect to the invariant probability measure
on the unit circle $\mathbb{T}$. Let $\Phi$ be the holomorphic bijection from $S_{\pi}$ onto $D$ given by $\Phi(z)=\tan (z / 2)$ when $z \in S_{\pi}$, extended to $|\operatorname{Re} z|=\pi / 2$ by the same formula. Then $\partial S_{\pi}$ is sent to $\mathbb{T} \backslash\{\mathrm{i},-\mathrm{i}\}$ and if $\Phi(\pi / 2+\mathrm{i} \tau)=\mathrm{e}^{\mathrm{i} \beta}$, we have $\beta \in(-\pi / 2, \pi / 2)$ and $\tanh (\tau / 2)=\tan (\beta / 2)$. For $r=\tan (\sigma / 2)$ we see that $\nu_{r}=\Phi_{\#} \mu_{\sigma}$ and

$$
\int_{B_{1}} h \mathrm{~d} \mu_{\sigma, 1}=\int_{\mathbb{R}} h(\pi / 2+\mathrm{i} \tau) f_{\sigma}(\tau) \mathrm{d} \tau \quad \text { with } f_{\sigma}(\tau)=\frac{\cos \sigma}{2 \pi(\cosh \tau-\sin \sigma)}
$$

while $\int_{B_{0}} h \mathrm{~d} \mu_{\sigma, 0}=\int_{\mathbb{R}} h(-\pi / 2+\mathrm{i} \tau) f_{-\sigma}(\tau) \mathrm{d} \tau$. One finds $\left\|\mu_{\sigma, 1}\right\|_{1}=\theta:=$ $\sigma / \pi+1 / 2$ and $\left\|\mu_{\sigma, 0}\right\|_{1}=1-\theta$ by harmonicity of $h(z)=\operatorname{Re} z$. When $\sigma$ tends to $\pi / 2$, the density $f_{\sigma}$ resembles the Cauchy kernel $P_{\varepsilon}^{(1)}$ in (1.33.C) with $\varepsilon=\pi / 2-\sigma$, since

$$
f_{\sigma}(\tau)=\frac{1}{2 \pi} \frac{\sin \varepsilon}{\cosh \tau-\cos \varepsilon} \simeq \frac{\varepsilon}{\pi\left(\tau^{2}+\varepsilon^{2}\right)}
$$

One can also comprehend $f_{\sigma}$ as sum of the alternate series of Cauchy kernels

$$
\begin{aligned}
f_{\sigma}=P_{\pi / 2-\sigma}^{(1)}-P_{\pi+\pi / 2+\sigma}^{(1)}+P_{2 \pi+\pi / 2-\sigma}^{(1)}- & P_{2 \pi+\pi+\pi / 2+\sigma}^{(1)} \\
& +P_{4 \pi+\pi / 2-\sigma}^{(1)}-P_{4 \pi+\pi+\pi / 2+\sigma}^{(1)}+\cdots
\end{aligned}
$$

indeed, if $\varphi_{\sigma}$ denotes the sum of the series above and if $g$ belongs to $\mathcal{K}(\mathbb{R})$, then $G(\sigma+\mathrm{i} \tau)=\left(\varphi_{\sigma} * g\right)(\tau)$ is harmonic in $S_{\pi}$, tends to $g(\tau)$ when $\sigma \rightarrow \pi / 2$ and to 0 when $\sigma \rightarrow-\pi / 2$, the same properties as for $\left(f_{\sigma} * g\right)(\tau)$.

Let $h_{*}$ be a continuous function on $\partial S_{\pi}$, and suppose that the two functions $t \mapsto \mathrm{e}^{-|t|} h_{*}( \pm \pi / 2+\mathrm{i} t)$ are Lebesgue-integrable on the real line. Then, writing

$$
\begin{equation*}
\widetilde{h}_{*}(z)=\int_{\mathbb{R}}\left(h_{*}(\pi / 2+\mathrm{i}(\tau-t)) f_{\sigma}(t)+h_{*}(-\pi / 2+\mathrm{i}(\tau-t)) f_{-\sigma}(t)\right) \mathrm{d} t \tag{3.16}
\end{equation*}
$$

for every $z=\sigma+\mathrm{i} \tau \in S_{\pi}$, one defines a harmonic function $\widetilde{h}_{*}$ in $S_{\pi}$, continuous on the closure if one sets $\widetilde{h}_{*}\left(z_{*}\right)=h_{*}\left(z_{*}\right)$ for $z_{*} \in \partial S_{\pi}$. Let $\mathcal{H}_{c}\left(S_{\pi}\right)$ denote the class of functions harmonic in $S_{\pi}$ and continuous on the closure. Not every $h \in \mathcal{H}_{c}\left(S_{\pi}\right)$ can be expressed by (3.16) from its restriction $h_{*}=\left.h\right|_{\partial S_{\pi}}$. First, $h_{*}$ must be $\mu_{\sigma}$-integrable, but even then, $h(z)=$ Re $\cos (z)=\cos (\sigma) \cosh (\tau)$, for which $h_{*}=0$, is a counterexample.

Let us say here that $g$ defined on $S_{\pi}$, resp. $\partial S_{\pi}$, is moderate if there is $a<1$ such that $g(z)=O\left(\mathrm{e}^{a|\operatorname{Im} z|}\right)$ for $z \in S_{\pi}$, resp. $\partial S_{\pi}$. If $h_{*}$ is moderate and continuous on $\partial S_{\pi}$, the extension $\widetilde{h}_{*}$ in (3.16) is in $\mathcal{H}_{c}\left(S_{\pi}\right)$, and it is
moderate because

$$
\begin{aligned}
\left|\widetilde{h}_{*}(\sigma+\mathrm{i} \tau)\right| & \leqslant \kappa \int_{\mathbb{R}} \mathrm{e}^{a|\tau-t|}\left(f_{\sigma}+f_{-\sigma}\right)(t) \mathrm{d} t \\
& \leqslant \kappa\left(\mathrm{e}^{a}+\int_{1}^{+\infty} \frac{\mathrm{e}^{a|t|}}{\cosh t-1} \mathrm{~d} t\right) \mathrm{e}^{a|\tau|}
\end{aligned}
$$

Lemma 3.5 (after [45]). - If $h \in \mathcal{H}_{c}\left(S_{\pi}\right)$ is moderate and $h_{*}=\left.h\right|_{\partial S_{\pi}}$, then $h=\widetilde{h}_{*}$.

If one replaces $S_{\pi}$ by a strip $S_{w}$ of width $w$, then clearly the moderation condition in $S_{w}$ must be formulated for $z \in S_{w}$ as $g(z)=O\left(\mathrm{e}^{a|\operatorname{Im} z|}\right)$ with $a<\pi / w$.

Proof. - We have that $h_{*}$ is moderate on $\partial S_{\pi}$, hence $U=h-\widetilde{h}_{*}$ is moderate on $S_{\pi}$ and vanishes on $\partial S_{\pi}$. Given $z_{0} \in S_{\pi}, a<1$ such that $U=O\left(\mathrm{e}^{a|\tau|}\right), \varepsilon>0$ and $b \in(a, 1)$, we see as in (3.13) that $U-\varepsilon \operatorname{Re} \cos (b z)$ is $\leqslant 0$ on the boundary of a rectangle containing $z_{0}$, hence $U\left(z_{0}\right) \leqslant \varepsilon \operatorname{Re} \cos \left(z_{0}\right)$ by the maximum principle. Doing it also with $-U$ and letting $\varepsilon \rightarrow 0$ we conclude that $h-\widetilde{h}_{*}=0$.

We now study the function $h_{1}$ defined by $h_{1}(\pi / 2+\mathrm{i} \tau)=|\tau|, h_{1}(-\pi / 2+$ $\mathrm{i} \tau)=0$ for every $\tau \in \mathbb{R}$ and its (moderate) harmonic extension given at $\sigma \in I_{\pi}$ by

$$
h_{1}(\sigma)=\int_{\mathbb{R}}|\tau| f_{\sigma}(\tau) \mathrm{d} \tau=\frac{2}{\pi} \int_{0}^{+\infty} \arctan \left(\frac{\cos \sigma}{\mathrm{e}^{\tau}-\sin \sigma}\right) \mathrm{d} \tau .
$$

Recall that $\left\|f_{\sigma}\right\|_{L^{1}(\mathbb{R})}=\theta=\sigma / \pi+1 / 2$. When $\sigma=0$, we have the easy bound

$$
h_{1}(0)=\frac{2}{\pi} \int_{0}^{+\infty} \arctan \left(\mathrm{e}^{-\tau}\right) \mathrm{d} \tau<\frac{2}{\pi} \int_{0}^{+\infty} \mathrm{e}^{-\tau} \mathrm{d} \tau=\frac{2}{\pi} .
$$

One can find $h_{1}(0)$ by writing the power series expansion of $\arctan (x)$, letting then $x=\mathrm{e}^{-\tau}$ and integrating in $\tau \in(0,+\infty)$. One gets $h_{1}(0)=$ $2 G / \pi<0.584$, where $G=\sum_{k=0}^{+\infty}(-1)^{k}(2 k+1)^{-2}$ is the Catalan constant, $0.915<G<0.916$. One has

$$
h_{1}^{\prime}(\sigma)=\frac{2}{\pi} \int_{0}^{+\infty} \frac{\mathrm{e}^{-2 \tau}-\mathrm{e}^{-\tau} \sin \sigma}{\mathrm{e}^{-2 \tau}-2 \mathrm{e}^{-\tau} \sin \sigma+1} \mathrm{~d} \tau=\frac{1}{\pi} \ln (2-2 \sin \sigma)
$$

thus $h_{1}$ is concave on $I_{\pi}$ and maximal when $\sigma=\pi / 6$. One can find numerically that $0.646<h_{1}(\pi / 6)<0.647$. By concavity, we obtain for each $\sigma \in I_{\pi}$ that

$$
\begin{equation*}
h_{1}(\sigma)=h_{1}(\sigma)-h_{1}(-\pi / 2) \leqslant h_{1}^{\prime}(-\pi / 2)(\sigma+\pi / 2)=\theta \ln 4 \tag{3.17}
\end{equation*}
$$

One has $h_{1}(\pi / 2)=0$, the behavior of $h_{1}(\sigma)$ when $\varepsilon=\pi / 2-\sigma \rightarrow 0$ is given by

$$
\begin{align*}
& h_{1}(\sigma)=-\frac{1}{\pi} \int_{0}^{\varepsilon} \ln (2-2 \cos s) \mathrm{d} s \\
& \simeq-\frac{1}{\pi} \int_{0}^{\varepsilon} \ln \left(s^{2}\right) \mathrm{d} s=\frac{2}{\pi}(\varepsilon \ln (1 / \varepsilon)+\varepsilon) \tag{3.18}
\end{align*}
$$

Since $h_{1}(\cdot+\mathrm{i} \tau)-h_{1}(\cdot)$ is bounded by $|\tau|$ on $B_{1}$ and vanishes on $B_{0}$, we have

$$
\begin{equation*}
0<h_{1}(\sigma+\mathrm{i} \tau) \leqslant h_{1}(\sigma)+\theta|\tau| \leqslant h_{1}(\pi / 6)+\theta|\tau|, \quad \sigma \in I_{\pi}, \quad \tau \in \mathbb{R} \tag{3.19}
\end{equation*}
$$

If $S_{w}=\left\{z \in \mathbb{C}: \alpha_{0}<\operatorname{Re} z<\alpha_{1}\right\}$ has width $w=\alpha_{1}-\alpha_{0}$ and if $\lambda=w / \pi$, we may associate to $h$, harmonic on $S_{w}$, the harmonic function $H(Z)=h\left(\alpha_{1 / 2}+\lambda Z\right)$ for $Z \in S_{\pi}$, where $\alpha_{t}=(1-t) \alpha_{0}+t \alpha_{1}$ when $t \in[0,1]$. If we set $h_{1, w}\left(\alpha_{1}+\mathrm{i} \tau\right)=|\tau|$ and $h_{1, w}=0$ on $\alpha_{0}+\mathrm{i} \mathbb{R}$, then $H_{1, w}=\lambda h_{1}$, and we get from (3.19) that

$$
\begin{aligned}
h_{1, w}\left(\alpha_{\theta}+\mathrm{i} \tau\right)=h_{1, w}\left(\alpha_{1 / 2}+\lambda \sigma+\mathrm{i} \tau\right) & =\lambda h_{1}(\sigma+\mathrm{i} \tau / \lambda) \\
& \leqslant w h_{1}(\pi / 6) / \pi+\theta|\tau|
\end{aligned}
$$

We now comment on Corollary 3.4. If $f$ is holomorphic in $S_{w}$ with admissible growth, satisfies $\left|f\left(\alpha_{j}+\mathrm{i} \tau\right)\right| \leqslant \mathrm{e}^{u_{j}|\tau|}$ on $\partial S_{w}, u_{j} \geqslant 0, j=0,1$, the "correct" bound at $z \in S_{w}$ for $f$ is $\mathrm{e}^{U_{\mathbf{u}, w}(z)}$ where $U_{\mathbf{u}, w}=u_{0} h_{0, w}+u_{1} h_{1, w}$, with $h_{0, w}\left(\alpha_{1 / 2}+\zeta\right)=h_{1, w}\left(\alpha_{1 / 2}-\zeta\right)$. One gets in particular $U_{\mathbf{u}, w}\left(\alpha_{1 / 2}\right)=$ $2 \lambda\left(u_{0}+u_{1}\right) G / \pi$. When $u_{0}=u_{1}=1$, this finer method gives at $\alpha_{1 / 2}$ a bounding factor $\mathrm{e}^{(4 G / \pi)(w / \pi)}$ instead of $E_{w, 1 / 2}(1,1)=\mathrm{e}^{w / 2}$, and $4 G / \pi^{2}<$ $0.3713<1 / 2$.

Let $V_{\mathbf{u}, w}$ be the harmonic conjugate of $U_{\mathbf{u}, w}$. Our first method in Corollary 3.4 applied to $f_{0}(z)=\mathrm{e}^{U_{\mathbf{u}, w}(z)+\mathrm{i} V_{\mathbf{u}, w}(z)}$ yields

$$
\begin{equation*}
U_{\mathbf{u}, w}\left(\alpha_{\theta}+\mathrm{i} \tau\right) \leqslant \ln E_{w, \theta}\left(u_{0}, u_{1}\right)+u_{\theta}|\tau| . \tag{3.20}
\end{equation*}
$$

If $u_{0}=u_{1}=u>0$, we get $U_{\mathbf{u}, w}\left(\alpha_{\theta}\right)=u\left(h_{0, w}+h_{1, w}\right)\left(\alpha_{\theta}\right) \leqslant w \sqrt{\theta(1-\theta)} u$. This estimate (3.20) has the right order of magnitude in $w$ and $u$, but not in $\theta$ when $\theta$ tends to 0 or 1 . The correct order when $\theta \rightarrow 0$ is $\kappa \theta \log (1 / \theta)$, according to (3.18).

Remark 3.6. - We shall have to deal with cases where the bounds on the lines limiting the strip $S_{w}=\left\{z \in \mathbb{C}: \alpha_{0}<\operatorname{Re} z<\alpha_{1}\right\}, w=\alpha_{1}-\alpha_{0}$, have the form

$$
\left|f\left(\alpha_{j}+\mathrm{i} \tau\right)\right| \leqslant\left(1+\tau^{2}\right)^{c_{j}} \mathrm{e}^{u_{j}|\tau|+v_{j}}, \quad c_{j}, u_{j} \geqslant 0, \quad j=0,1 .
$$

It is obviously possible to "absorb" the polynomial factor by replacing $u_{j}$ in the exponential with $u_{j}+\varepsilon, \varepsilon>0$ arbitrary, and modifying $v_{j}$ accordingly, but one can work a little more carefully as follows.

Let $\ell_{1, w}$ be the moderate harmonic function on $S_{w}$ such that $\ell_{1, w}\left(\alpha_{1}+\right.$ $\mathrm{i} \tau)=\ln \left(1+\tau^{2}\right)$ for $\tau \in \mathbb{R}$ and $\ell_{1, w}=0$ on $\alpha_{0}+\mathrm{i} \mathbb{R}$. Let $\alpha_{\theta}=(1-\theta) \alpha_{0}+\theta \alpha_{1}$, $\lambda=w / \pi, \sigma=\pi \theta-\pi / 2$ and $L_{1, w}(Z)=\ell_{1, w}\left(\alpha_{1 / 2}+\lambda Z\right)$. By Lemma 3.5 and (3.16), we get

$$
\begin{aligned}
& \ell_{1, w}\left(\alpha_{\theta}+\mathrm{i} \tau\right)=\ell_{1, w}\left(\alpha_{1 / 2}+\lambda \sigma+\mathrm{i} \tau\right)=L_{1, w}(\sigma+\mathrm{i} \tau / \lambda) \\
& =\int_{\mathbb{R}} L_{1, w}(\pi / 2+\mathrm{i} \tau / \lambda-\mathrm{i} t) f_{\sigma}(t) \mathrm{d} t \leqslant \int_{\mathbb{R}} \ln \left(1+(\lambda|t|+|\tau|)^{2}\right) f_{\sigma}(t) \mathrm{d} t
\end{aligned}
$$

Applying Jensen's inequality to the probability density $\widetilde{f}_{\sigma}=\theta^{-1} f_{\sigma}$, one sees that

$$
\begin{aligned}
& \exp \left(\int_{\mathbb{R}} \ln \left(\left[1+(\lambda|t|+|\tau|)^{2}\right]^{1 / 2}\right) \widetilde{f}_{\sigma}(t) \mathrm{d} t\right) \\
& \quad \leqslant \int_{\mathbb{R}}\left[1+(\lambda|t|+|\tau|)^{2}\right]^{1 / 2} \widetilde{f}_{\sigma}(t) \mathrm{d} t \leqslant\left(1+\tau^{2}\right)^{1 / 2}+\lambda \int_{\mathbb{R}}|t| \widetilde{f}_{\sigma}(t) \mathrm{d} t
\end{aligned}
$$

bounded by $\left(1+\tau^{2}\right)^{1 / 2}+\lambda \ln 4$ by (3.17). For every $\tau \in \mathbb{R}$, one has therefore

$$
\begin{equation*}
0<\ell_{1, w}\left(\alpha_{\theta}+\mathrm{i} \tau\right)<2 \theta \ln \left(\left(1+\tau^{2}\right)^{1 / 2}+\lambda \ln 4\right) \tag{3.21}
\end{equation*}
$$

Define a harmonic function $U$ in $S_{w}$, continuous on the closure, by $U=$ $c_{0} \ell_{0, w}+c_{1} \ell_{1, w}$, where $\ell_{0, w}(z)=\ell_{1, w}\left(2 \alpha_{1 / 2}-z\right)$, so that $U\left(\alpha_{j}+\mathrm{i} \tau\right)=$ $c_{j} \ln \left(1+\tau^{2}\right)$. Let $V$ be conjugate to $U$ in $S_{w}$. Then $g=\mathrm{e}^{-U-\mathrm{i} V}$ is holomorphic in $S_{w}$ and $\left|(f g)\left(\alpha_{j}+\mathrm{i} \tau\right)\right| \leqslant \mathrm{e}^{u_{j}|\tau|+v_{j}}$ on $\partial S_{w}$. By (3.21), we can bound $|f(z)|$ at $z=\alpha_{\theta}+\mathrm{i} \tau$ by multiplying the inside bound of Corollary 3.4 for $f g$ with the additional factor

$$
\mathrm{e}^{U\left(\alpha_{\theta}+\mathrm{i} \tau\right)} \leqslant\left(\left(1+\tau^{2}\right)^{1 / 2}+\ln (4) w / \pi\right)^{2 c_{\theta}} \leqslant(1+\ln (4) w / \pi)^{2 c_{\theta}}\left(1+\tau^{2}\right)^{c_{\theta}}
$$ where $c_{\theta}=(1-\theta) c_{0}+\theta c_{1}$. Since $\ln (4) / \pi<1 / 2$, we may remember that

$$
\begin{equation*}
\left|f\left(\alpha_{\theta}+\mathrm{i} \tau\right)\right| \leqslant(1+w / 2)^{2 c_{\theta}} E_{w, \theta}\left(u_{0}, u_{1}\right)\left(1+\tau^{2}\right)^{c_{\theta}} \mathrm{e}^{u_{\theta}|\tau|+v_{\theta}} \tag{3.22}
\end{equation*}
$$

### 3.2.2. Interpolation of holomorphic families of linear operators

We now recall the classical complex interpolation method for bounding in the norm of $L^{p}(X, \Sigma, \mu)$, when $1<p<+\infty$, a linear operator $T_{\alpha}$ that is a member of a holomorphic family of operators $\left(T_{z}\right)$, for $z$ in a vertical strip $S$ containing $\alpha$. We consider a linear space $\mathcal{E}$ which is a common subspace of all $L^{r}(X, \Sigma, \mu), 1 \leqslant r \leqslant+\infty$, and which is dense in $L^{r}(X, \Sigma, \mu)$ when $1 \leqslant r<+\infty$. This space $\mathcal{E}$ can be the space of simple $\Sigma$-measurable and $\mu$-integrable functions, or for the specific spaces $L^{r}\left(\mathbb{R}^{n}\right)$, it can be $\mathcal{S}\left(\mathbb{R}^{n}\right)$ or the space $\mathcal{K}\left(\mathbb{R}^{n}\right)$. We consider a closed strip $\alpha_{0} \leqslant \operatorname{Re} z \leqslant \alpha_{1}$ in $\mathbb{C}$, with $\alpha_{0}<\alpha<\alpha_{1}$. We assume that each $T_{z}$, for $z$ in this closed strip, is defined on $\mathcal{E}$ and linear with values in $L^{1}(X, \Sigma, \mu)+L^{\infty}(X, \Sigma, \mu)$. The holomorphy assumption means that for $f, g \in \mathcal{E}$, the function $z \mapsto\left\langle T_{z} f, g\right\rangle$ is holomorphic
in the open strip $\alpha_{0}<\operatorname{Re} z<\alpha_{1}$, but one also assumes that it extends as a continuous function on the closed strip. The above bracket is bilinear, given by $\int_{X}\left(T_{z} f\right) g \mathrm{~d} \mu$. Later in these Notes, we shall abuse slightly and speak about holomorphic family of linear operators in the closed strip $\alpha_{0} \leqslant \operatorname{Re} z \leqslant$ $\alpha_{1}$.

We consider $1 \leqslant p_{0}, p_{1} \leqslant+\infty$ and $p$ between $p_{0}$ and $p_{1}$, so that $1<$ $p<+\infty$. We assume that when $\operatorname{Re} z=\alpha_{j}, j=0,1$, the $T_{z}$ s are uniformly bounded from $\mathcal{E}$, equipped with the $L^{p_{j}}$ norm, to $L^{p_{j}}(X, \Sigma, \mu)$, and we assume that for a certain $\theta \in(0,1)$, we have both

$$
\frac{1}{p}=\frac{1-\theta}{p_{0}}+\frac{\theta}{p_{1}} \quad \text { and } \quad \alpha=(1-\theta) \alpha_{0}+\theta \alpha_{1} .
$$

We want to show that $T_{\alpha}$ is bounded from $\mathcal{E}$, equipped with the $L^{p}$ norm, to $L^{p}(X, \Sigma, \mu)$. Then, by the density of $\mathcal{E}$, we will be able to extend to $L^{p}(X, \Sigma, \mu)$ the bound obtained for the functions in $\mathcal{E}$.

We must of course bound $\left\langle T_{\alpha} f, g\right\rangle$, uniformly for $f$ in the intersection of $\mathcal{E}$ with the unit ball of $L^{p}(X, \Sigma, \mu)$ and for $g$ in the unit ball of the dual $L^{q}(X, \Sigma, \mu), 1 / p+1 / q=1$. Denote by $q_{0}$ the conjugate of $p_{0}$ and by $q_{1}$ that of $p_{1}$. Observe that we have also $1 / q=(1-\theta) / q_{0}+\theta / q_{1}$. We write $f(x)=$ $u(x)|f(x)|, g(x)=v(x)|g(x)|$ for every $x \in X$, with $|u(x)|=|v(x)|=1$. Next, for each $z \in \mathbb{C}$, we set

$$
\begin{equation*}
f_{z}(x)=u(x)|f(x)|^{p(s z+t)}, \quad g_{z}(x)=v(x)|g(x)|^{q(1-s z-t)}, \quad x \in X \tag{3.23}
\end{equation*}
$$

where $s, t$ real are chosen such that $s \alpha_{0}+t=1 / p_{0}$ and $s \alpha_{1}+t=1 / p_{1}$. This yields $s \alpha+t=1 / p$. We see that $f_{\alpha}=f, g_{\alpha}=g$ and we also see that the exponents have been chosen so that the assumptions $\|f\|_{p} \leqslant 1$ and $\|g\|_{q} \leqslant 1$ imply

$$
\forall \tau \in \mathbb{R}, \quad\left\|f_{\alpha_{0}+\mathrm{i} \tau}\right\|_{p_{0}} \leqslant 1, \quad\left\|f_{\alpha_{1}+\mathrm{i} \tau}\right\|_{p_{1}} \leqslant 1, \quad . \quad\left\|g_{\alpha_{0}+\mathrm{i} \tau}\right\|_{q_{0}} \leqslant 1, \quad\left\|g_{\alpha_{1}+\mathrm{i} \tau}\right\|_{q_{1}} \leqslant 1 .
$$

We notice for future reference that if $f$ and $g$ are bounded by $M$ on $X$, then

$$
\begin{equation*}
\left|f_{z}\right| \leqslant \max \left(M^{p / p_{0}}, M^{p / p_{1}}\right), \quad\left|g_{z}\right| \leqslant \max \left(M^{q / q_{0}}, M^{q / q_{1}}\right) \tag{3.24}
\end{equation*}
$$

when $\alpha_{0} \leqslant \operatorname{Re} z \leqslant \alpha_{1}$, because $\operatorname{Re}(s z+t)$ stays between $1 / p_{0}$ and $1 / p_{1}$ and $\operatorname{Re}(1-s z-t)$ between $1 / q_{0}$ and $1 / q_{1}$ when $z \in S$. We now apply the three lines Lemma 3.1 for bounding the value $H(\alpha)=\left\langle T_{\alpha} f, g\right\rangle$ of the holomorphic function

$$
\begin{equation*}
H: z \mapsto\left\langle T_{z} f_{z}, g_{z}\right\rangle, \quad z \in S \tag{3.25}
\end{equation*}
$$

from the bounds on the lines $\operatorname{Re} z=\alpha_{0}$ and $\operatorname{Re} z=\alpha_{1}$. When $\operatorname{Re} z=\alpha_{j}$, we get

$$
|H(z)|=\left|\left\langle T_{z} f_{z}, g_{z}\right\rangle\right| \leqslant\left\|T_{z}\right\|_{p_{j} \rightarrow p_{j}}\left\|f_{z}\right\|_{p_{j}}\left\|g_{z}\right\|_{q_{j}} \leqslant\left\|T_{z}\right\|_{p_{j} \rightarrow p_{j}}
$$

for $j=0,1$. In addition, the holomorphic function $H$ must be bounded on the strip, see Remark 3.2 above. If true, we know by Lemma 3.1 that

$$
|H(\alpha)|=\left|\left\langle T_{\alpha} f, g\right\rangle\right| \leqslant\left(\sup _{\tau \in \mathbb{R}}\left\|T_{\alpha_{0}+\mathrm{i} \tau}\right\|_{p_{0} \rightarrow p_{0}}\right)^{1-\theta}\left(\sup _{\tau \in \mathbb{R}}\left\|T_{\alpha_{1}+\mathrm{i} \tau}\right\|_{p_{1} \rightarrow p_{1}}\right)^{\theta}
$$

and by taking the supremum over $f$ and $g$, we obtain

$$
\begin{equation*}
\left\|T_{\alpha}\right\|_{p \rightarrow p} \leqslant\left(\sup _{\tau \in \mathbb{R}}\left\|T_{\alpha_{0}+\mathrm{i} \tau}\right\|_{p_{0} \rightarrow p_{0}}\right)^{1-\theta}\left(\sup _{\tau \in \mathbb{R}}\left\|T_{\alpha_{1}+\mathrm{i} \tau}\right\|_{p_{1} \rightarrow p_{1}}\right)^{\theta} \tag{3.26}
\end{equation*}
$$

Finally, we can extend $T_{\alpha}$ from the dense subspace $\mathcal{E}$ to $L^{p}(X, \Sigma, \mu)$. Sometimes, rather than looking for extension, one obtains in this way a sharper estimate for the norm of an operator $T_{\alpha}$ already known to be bounded on $L^{p}(X, \Sigma, \mu)$.

This complex method, introduced for $L^{p}$ spaces by Thorin [80, 81] for one linear operator, extended by Stein [70] to families, can also be extended (see [6]) to spaces of the form $L^{p}\left(L^{r}\right)$ and more generally, by the abstract complex interpolation method due to Calderón [19], to a pair of the form $\left(L^{p_{0}}\left(A_{0}\right), L^{p_{1}}\left(A_{1}\right)\right)$. One then obtains estimates in $L^{p}\left(A_{\theta}\right)$, where $A_{\theta}$ is the space associated to the pair $\left(A_{0}, A_{1}\right)$ by Calderón's method with parameter $\theta \in(0,1)$.

In many cases later on, the norms of the operators $\left(T_{z}\right)_{z \in S}$ are not uniformly bounded on the boundary lines, but obey for some $\lambda>0$ estimates of the form

$$
\left\|T_{\alpha_{0}+\mathrm{i} \tau}\right\|_{p_{0} \rightarrow p_{0}} \leqslant C_{0} \mathrm{e}^{\lambda|\tau|}, \quad\left\|T_{\alpha_{1}+\mathrm{i} \tau}\right\|_{p_{1} \rightarrow p_{1}} \leqslant C_{1} \mathrm{e}^{\lambda|\tau|}, \quad \tau \in \mathbb{R}
$$

Using Corollary 3.4, we can handle this situation. We must simply check that the above function $H(z)=H_{f, g}(z)$ in (3.25) has an admissible growth in the strip. We have to find an ad hoc argument giving such a growth for each choice of $f$ and $g$ in suitable dense subsets, growth depending on $f, g$. Indeed, in general, we do not know yet bounds on the norm $\left\|T_{z}\right\|_{p_{z} \rightarrow p_{z}}$ for $z \in S$, where $w / p_{z}=\left(\alpha_{1}-\operatorname{Re} z\right) / p_{0}+\left(\operatorname{Re} z-\alpha_{0}\right) / p_{1}$ and where $w=\alpha_{1}-\alpha_{0}$ is the width of $S$. If each function $H_{f, g}$ has an admissible growth in $S$, we obtain here at last that

$$
\left\|T_{\alpha}\right\|_{p \rightarrow p} \leqslant C_{0}^{1-\theta} C_{1}^{\theta} \mathrm{e}^{\lambda w \sqrt{\theta(1-\theta)}}
$$

If an additional polynomial factor is present in the bound of $\left\|T_{\alpha_{j}+\mathrm{i} \tau}\right\|_{p_{j} \rightarrow p_{j}}$, $j=0,1$, then we make use of Remark 3.6 and of the estimate (3.22).

### 3.3. On the definition of maximal functions

Let us consider a family $\left(K_{t}\right)_{t>0}$ of integrable functions on $\mathbb{R}^{n}$ and define a related maximal function by the formula

$$
\begin{equation*}
\mu(f)=\sup _{t>0}\left|K_{t} * f\right| \tag{3.27}
\end{equation*}
$$

for $f \in L^{p}\left(\mathbb{R}^{n}\right), 1 \leqslant p<+\infty$. We are faced with a standard difficulty of processes with continuous time parameter. In this generality, the convolution $K_{t} * f$ is only defined almost everywhere, for each $t>0$, and the preceding supremum is not a well defined equivalence class of measurable functions. However, if $D$ is a countable subset of $(0,+\infty)$, there is no problem in considering

$$
\mu_{D}(f)=\sup _{t \in D}\left|K_{t} * f\right|
$$

and a classical workaround for defining $\mu(f)$ consists in introducing the essential supremum: there is a countable subset $D_{0} \subset(0,+\infty)$ such that $\mu_{D}(f)=\mu_{D_{0}}(f)$ almost everywhere, whenever $D \supset D_{0}$. In other words, for every $t>0$, we then have $\left|K_{t} * f\right| \leqslant \mu_{D_{0}}(f)$ almost everywhere. The essential supremum is defined to be the equivalence class of $\mu_{D_{0}}(f)$. It is also the least upper bound of the family $\left(\left|K_{t} * f\right|\right)_{t>0}$ in the Banach lattice $L^{p}\left(\mathbb{R}^{n}\right)$.

Most often, we shall have the specific problem where one considers an integrable kernel $K$ on $\mathbb{R}^{n}$ and defines a maximal function using the dilates of $K$, by

$$
\mu(f)=\sup _{t>0}\left|K_{(t)} * f\right|
$$

If $f \in L^{p}\left(\mathbb{R}^{n}\right)$ and if $K$ belongs to $L^{q}\left(\mathbb{R}^{n}\right)$, with $q<+\infty$ and $1 / q+1 / p=1$, then $K_{(t)} * f$ is defined pointwise and $t \mapsto K_{(t)}$ is continuous from $(0,+\infty)$ to $L^{q}$. It follows that $t \mapsto\left(K_{(t)} * f\right)(x)$ is continuous for every $f \in L^{p}\left(\mathbb{R}^{n}\right)$, $x \in \mathbb{R}^{n}$, and the aforementioned problem disappears. If $K \in L^{1}\left(\mathbb{R}^{n}\right)$ and $f \in L^{p}\left(\mathbb{R}^{n}\right)$ are nonnegative, then $\left(K_{(t)} * f\right)(x)$ is a definite value in $[0,+\infty]$ for every $x \in \mathbb{R}^{n}$, but it is not immediately clear that a direct application of (3.27) gives what we want. However, we can find an increasing sequence $\left(f_{k}\right)_{k \geqslant 0}$ of bounded nonnegative Borel functions tending almost everywhere to $f$. Then for every $x \in \mathbb{R}^{n}$ and $k \geqslant 0$, the map $t \mapsto\left(K_{(t)} * f_{k}\right)(x)$ is continuous from $(0,+\infty)$ to $[0,+\infty)$, because $t \mapsto K_{(t)}$ is continuous from $(0,+\infty)$ to $L^{1}\left(\mathbb{R}^{n}\right)$. It follows that $t \mapsto\left(K_{(t)} * f\right)(x)$ is lower semi-continuous, since it is an increasing limit of continuous functions. For every countable dense set $D$ one has thus

$$
\mu_{D}(f)(x)=\sup _{s \in D}\left(K_{(s)} * f\right)(x)=\sup _{t>0}\left(K_{(t)} * f\right)(x)
$$

This argument does not apply to kernels that can also assume negative values, and it is precisely the case that will appear later.

We will have to investigate maximal functions such as $\mu(f)=$ $\sup _{t>0}\left|K_{(t)} * f\right|$, usually when $K \in L^{1}\left(\mathbb{R}^{n}\right)$, but also more generally when $K$ is a bounded measure on $\mathbb{R}^{n}$. It will be often convenient to start the study with nice functions, for example functions $\varphi$ belonging to the Schwartz class $\mathcal{S}\left(\mathbb{R}^{n}\right)$, for which $\mu(\varphi)$ is clearly defined. If a function $f \in L^{p}\left(\mathbb{R}^{n}\right)$ is given and since $\mathcal{S}\left(\mathbb{R}^{n}\right)$ is dense in $L^{p}\left(\mathbb{R}^{n}\right)$, we may find for every $\varepsilon>0$ a sequence $\left(\varphi_{k}\right)_{k \geqslant 0}$ in $\mathcal{S}\left(\mathbb{R}^{n}\right)$ such that

$$
f=\sum_{k=0}^{+\infty} \varphi_{k} \text { in } L^{p}\left(\mathbb{R}^{n}\right), \quad \text { and } \quad \sum_{k=0}^{+\infty}\left\|\varphi_{k}\right\|_{p}<\|f\|_{p}+\varepsilon
$$

Since the convolution with $K_{(t)}$ is linear and continuous on $L^{p}\left(\mathbb{R}^{n}\right)$, we have

$$
K_{(t)} * f=\sum_{k=0}^{+\infty} K_{(t)} * \varphi_{k} \text { in } L^{p}\left(\mathbb{R}^{n}\right), \quad \text { and } \sum_{k=0}^{+\infty}\left\|K_{(t)} * \varphi_{k}\right\|_{p}<+\infty
$$

so the series $\sum_{k=0}^{+\infty} K_{(t)} * \varphi_{k}$ converges also almost everywhere to $K_{(t)} * f$, and we have almost everywhere

$$
\left|K_{(t)} * f\right| \leqslant \sum_{k=0}^{+\infty}\left|K_{(t)} * \varphi_{k}\right| \leqslant \sum_{k=0}^{+\infty} \mu\left(\varphi_{k}\right)
$$

For any countable subset $D \subset(0,+\infty)$ we get $\mu_{D}(f) \leqslant \sum_{k=0}^{+\infty} \mu\left(\varphi_{k}\right)$, implying that $\mu(f)$, defined as essential supremum, is bounded by $\sum_{k=0}^{+\infty} \mu\left(\varphi_{k}\right)$. If we know that there exists $\kappa$ such that $\|\mu(\varphi)\|_{p} \leqslant \kappa\|\varphi\|_{p}$ when $\varphi \in \mathcal{S}\left(\mathbb{R}^{n}\right)$, it follows that

$$
\|\mu(f)\|_{p} \leqslant \sum_{k=0}^{+\infty}\left\|\mu\left(\varphi_{k}\right)\right\|_{p} \leqslant \kappa \sum_{k=0}^{+\infty}\left\|\varphi_{k}\right\|_{p} \leqslant \kappa\left(\|f\|_{p}+\varepsilon\right)
$$

for every $\varepsilon>0$. In order to bound $\mu(f)$ in $L^{p}\left(\mathbb{R}^{n}\right)$, it is therefore enough to obtain a uniform bound for Schwartz functions. Clearly, any dense linear subspace of $L^{p}\left(\mathbb{R}^{n}\right)$ consisting of nice functions can be used instead of $\mathcal{S}\left(\mathbb{R}^{n}\right)$.

The classical maximal function $\mathrm{M} f$, as well as $\mathrm{M}_{C} f$ in (0.3.M), is actually defined by means of $\sup _{t>0} K_{(t)} *|f|$. This makes sense whenever the kernel $K$ is nonnegative, but not for a general $K$. We shall distinguish

$$
\mathrm{M}_{K} f:=\sup _{t>0} K_{(t)} *|f| \quad \text { and } \quad M_{K} f:=\sup _{t>0}\left|K_{(t)} * f\right|
$$

by the tiny notational difference between the slanted or unslanted letter M. When the kernel $K$ is nonnegative, we have obviously $M_{K} f \leqslant \mathrm{M}_{K} f=$ $M_{K}(|f|)$.

## 4. The results of Stein for Euclidean balls

We prove here the remarkable fact due to Stein [75] that for $p>1$, the maximal operator associated to Euclidean balls, i.e., the classical HardyLittlewood maximal operator $M$ defined in (0.1), may be bounded in $L^{p}\left(\mathbb{R}^{n}\right)$ independently of the dimension $n$. Full details appeared in [77]. Other proofs have appeared since then, let us mention Auscher and Carro [4] who found the simple explicit bound $2+\sqrt{2}$ in $L^{2}\left(\mathbb{R}^{n}\right)$, extended by interpolation as $(2+$ $\sqrt{2})^{2 / p}$ for $p \geqslant 2$. It is not known whether or not the weak $(1,1)$ norm of the maximal operator M is also bounded independently of the dimension. Even if we shall not develop this weak type aspect mentioned in our introduction, let us recall that the best upper estimate that is known for the weak $(1,1)$ norm of M is the Stein-Strömberg $\mathrm{O}(n)$ bound [77].

ThEOREM 4.1 (Stein [75]). - Let $1<p \leqslant+\infty$. For every integer $n \geqslant 1$ and all functions $f \in L^{p}\left(\mathbb{R}^{n}\right)$, one has that

$$
\|\mathrm{M} f\|_{L^{p}\left(\mathbb{R}^{n}\right)} \leqslant C(p)\|f\|_{L^{p}\left(\mathbb{R}^{n}\right)}
$$

where $C(p)$ is a constant independent of the dimension $n$.

### 4.1. Proof of Theorem 4.1

The main tool in the proof is the spherical maximal operator $\mathcal{M}$ defined by

$$
(\mathcal{M} f)(x)=\left(M_{\sigma} f\right)(x)=\sup _{r>0}\left|\int_{S^{n-1}} f(x-r \theta) \mathrm{d} \sigma(\theta)\right|, \quad x \in \mathbb{R}^{n}
$$

where $\sigma$ is the normalized Haar measure on the unit sphere $S^{n-1}$. It is clear that $\mathcal{M} f$ is well defined when $f$ is regular, but not when $f \in L_{\mathrm{loc}}^{1}\left(\mathbb{R}^{n}\right)$. Theorem 4.2 below means in particular that for suitable $p$ and $n, \mathcal{M} f$ can be defined when $f \in L^{p}\left(\mathbb{R}^{n}\right)$, for example by the method described at the end of Section 3.3. The maximal function $\mathcal{M}(|f|)$ controls $\mathrm{M} f$ pointwise, as one sees easily by using polar coordinates. The maximal operator $\mathcal{M}$ is bounded in $L^{p}\left(\mathbb{R}^{N}\right)$ for some $p$ and $N$, with a bound depending on the dimension $N$, according to the following theorem also due to Stein. An extension by Bourgain of this result can be found in [8].

Theorem 4.2 (Stein [74]). - Let $N \geqslant 3$ and assume that $N /(N-1)<$ $p \leqslant+\infty$. There exists a constant $C(N, p)$ such that for every function $f \in$ $L^{p}\left(\mathbb{R}^{N}\right)$, one has

$$
\|\mathcal{M} f\|_{L^{p}\left(\mathbb{R}^{N}\right)} \leqslant C(N, p)\|f\|_{L^{p}\left(\mathbb{R}^{N}\right)}
$$

The condition $p>N /(N-1)$ can be easily seen necessary, and the case $p=+\infty$ is obvious, with $C(N, \infty)=1$. We postpone the proof of this theorem to the next section. It requires a number of harmonic analysis methods, including square function, multipliers and Littlewood-Paley decomposition.

In order to prove Theorem 4.1, we first introduce the following weighted maximal operator, depending on a parameter $k \in \mathbb{N}$. For $f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$, let

$$
\left(\mathrm{M}_{n, k} f\right)(x)=\sup _{r>0} \frac{\int_{|y| \leqslant r}|f(x-y)||y|^{k} \mathrm{~d} y}{\int_{|y| \leqslant r}|y|^{k} \mathrm{~d} y}, \quad x \in \mathbb{R}^{n}
$$

where $|y|$ denotes the Euclidean norm of $y \in \mathbb{R}^{n}$. Taking polar coordinates gives us the pointwise inequality

$$
\left(\mathrm{M}_{n, k} f\right)(x) \leqslant(\mathcal{M}|f|)(x), \quad x \in \mathbb{R}^{n}
$$

from which we can deduce by applying Theorem 4.2 that for every integer $N \geqslant 3$, for $p$ such that $N /(N-1)<p \leqslant+\infty$ and for every $f$ in $L^{p}\left(\mathbb{R}^{N}\right)$, we have

$$
\begin{equation*}
\left\|\mathrm{M}_{N, k} f\right\|_{L^{p}\left(\mathbb{R}^{N}\right)} \leqslant C(N, p)\|f\|_{L^{p}\left(\mathbb{R}^{N}\right)} \tag{4.1}
\end{equation*}
$$

where $C(N, p)$ is the constant in Theorem 4.2. We shall obtain Theorem 4.1 by lifting to $\mathbb{R}^{n}$ the inequality (4.1) obtained in a lower dimension $N=n-k$. This is done by integrating over the Grassmannian of $(n-k)$-planes in $\mathbb{R}^{n}$. This method of descent is in the spirit of the Calderón-Zygmund method of rotations.

We write $\mathbb{R}^{n}=\mathbb{R}^{n-k} \times \mathbb{R}^{k}$ and $x=\left(x_{1}, x_{2}\right)$ accordingly, for every $x \in \mathbb{R}^{n}$, with $x_{1} \in \mathbb{R}^{n-k}$ and $x_{2} \in \mathbb{R}^{k}$. For each $U$ in the orthogonal group $\mathcal{O}(n)$, we introduce the auxiliary maximal operator

$$
\left(\mathrm{M}_{k}^{U} f\right)(x)=\sup _{r>0} \frac{\int_{\mathbb{R}^{n-k}} \mathbf{1}_{\left\{\left|y_{1}\right| \leqslant r\right\}}\left|f\left(x-U\left(y_{1}, 0\right)\right)\right|\left|y_{1}\right|^{k} \mathrm{~d} y_{1}}{\int_{\mathbb{R}^{n-k}} \mathbf{1}_{\left\{\left|y_{1}\right| \leqslant r\right\}}\left|y_{1}\right|^{k} \mathrm{~d} y_{1}}, \quad x \in \mathbb{R}^{n}
$$

We need two lemmas.
Lemma 4.3. - Let $n \geqslant k+3$ and $p>(n-k) /(n-k-1)$. Then for all $f \in L^{p}\left(\mathbb{R}^{n}\right)$ and $U \in \mathcal{O}(n)$, we have

$$
\left\|\mathrm{M}_{k}^{U} f\right\|_{L^{p}\left(\mathbb{R}^{n}\right)} \leqslant C(n-k, p)\|f\|_{L^{p}\left(\mathbb{R}^{n}\right)}
$$

where $C(n-k, p)$ is the constant appearing in Theorem 4.2.
Proof. - Let us set $f_{[U]}(x)=f(U x)$, for every $x \in \mathbb{R}^{n}$. Since $U \in \mathcal{O}(n)$, the mapping $S_{U}: f \mapsto f_{[U]}$ is an isometry of $L^{p}\left(\mathbb{R}^{n}\right)$. Observe that

$$
\int_{\left|y_{1}\right| \leqslant r}\left|f\left(U x-U\left(y_{1}, 0\right)\right)\right|\left|y_{1}\right|^{k} \mathrm{~d} y_{1}=\int_{\left|y_{1}\right| \leqslant r}\left|f_{[U]}\left(x-\left(y_{1}, 0\right)\right)\right|\left|y_{1}\right|^{k} \mathrm{~d} y_{1}
$$

hence we have that $\left(\mathrm{M}_{k}^{U} f\right)(U x)=\left(\mathrm{M}_{k}^{\mathrm{Id}} f_{[U]}\right)(x)$, for every $x \in \mathbb{R}^{n}$. This means that $S_{U} \mathrm{M}_{k}^{U}=\mathrm{M}_{k}^{\mathrm{Id}} S_{U}$. It follows that we need only consider $\mathrm{M}_{k}^{\mathrm{Id}}$. Now, for every $x=\left(x_{1}, x_{2}\right) \in \mathbb{R}^{n}$, we have

$$
\begin{aligned}
\left(\mathrm{M}_{k}^{\mathrm{Id}} f\right)\left(x_{1}, x_{2}\right) & =\sup _{r>0} \frac{\int_{\mathbb{R}^{n-k}} \mathbf{1}_{\left\{\left|y_{1}\right| \leqslant r\right\}}\left|f\left(x_{1}-y_{1}, x_{2}\right)\right|\left|y_{1}\right|^{k} \mathrm{~d} y_{1}}{\int_{\mathbb{R}^{n-k}} \mathbf{1}_{\left\{\left|y_{1}\right| \leqslant r\right\}}\left|y_{1}\right|^{k} \mathrm{~d} y_{1}} \\
& =\left(\mathrm{M}_{n-k, k} f_{x_{2}}\right)\left(x_{1}\right)
\end{aligned}
$$

with $f_{x_{2}}\left(x_{1}\right)=f\left(x_{1}, x_{2}\right)$. Applying (4.1) to $\mathrm{M}_{n-k, k}$ for each $x_{2} \in \mathbb{R}^{k}$ gives

$$
\int_{\mathbb{R}^{n-k}}\left|\left(\mathrm{M}_{k}^{\mathrm{Id}} f\right)\left(x_{1}, x_{2}\right)\right|^{p} \mathrm{~d} x_{1} \leqslant C(n-k, p)^{p} \int_{\mathbb{R}^{n-k}}\left|f_{x_{2}}\left(x_{1}\right)\right|^{p} \mathrm{~d} x_{1}
$$

therefore

$$
\begin{aligned}
\left\|\mathrm{M}_{k}^{\mathrm{Id}} f\right\|_{L^{p}\left(\mathbb{R}^{n}\right)}^{p} & \leqslant C(n-k, p)^{p} \int_{\mathbb{R}^{k}}\left(\int_{\mathbb{R}^{n-k}}\left|f_{x_{2}}\left(x_{1}\right)\right|^{p} \mathrm{~d} x_{1}\right) \mathrm{d} x_{2} \\
& =C(n-k, p)^{p}\|f\|_{L^{p}\left(\mathbb{R}^{n}\right)}^{p}
\end{aligned}
$$

Lemma 4.4. - For every locally integrable function $f$ on $\mathbb{R}^{n}$ and $1 \leqslant$ $k \leqslant n$, one has the pointwise inequality

$$
(\mathrm{M} f)(x) \leqslant \int_{\mathcal{O}(n)}\left(\mathrm{M}_{k}^{U} f\right)(x) \mathrm{d} \mu_{n}(U), \quad x \in \mathbb{R}^{n}
$$

where $\mu_{n}$ denotes the normalized Haar measure on $\mathcal{O}(n)$.
Proof. - The desired pointwise inequality follows from the next equality, true for every nonnegative Borel function $g$ on $\mathbb{R}^{n}$, stating that

$$
\begin{equation*}
\frac{\int_{|y| \leqslant r} g(y) \mathrm{d} y}{\int_{|y| \leqslant r} \mathrm{~d} y}=\frac{\int_{\mathcal{O}(n)} \int_{\mathbb{R}^{n-k}} \mathbf{1}_{\left\{\left|y_{1}\right| \leqslant r\right\}} g\left(U\left(y_{1}, 0\right)\right)\left|y_{1}\right|^{k} \mathrm{~d} y_{1} \mathrm{~d} \mu_{n}(U)}{\int_{\mathbb{R}^{n-k}} \mathbf{1}_{\left\{\left|y_{1}\right| \leqslant r\right\}}\left|y_{1}\right|^{k} \mathrm{~d} y_{1}} \tag{4.2}
\end{equation*}
$$

Indeed, for each $r>0$ and $x \in \mathbb{R}^{n}$, the previous equality allows us to write

$$
\begin{aligned}
\frac{1}{\left|B_{r}\right|} \int_{B_{r}}|f(x-y)| \mathrm{d} y & =\frac{\int_{\mathcal{O}(n)} \int_{\left|y_{1}\right| \leqslant r} \mid f\left(\left.\left(x-U\left(y_{1}, 0\right)\right)| | y_{1}\right|^{k} \mathrm{~d} y_{1} \mathrm{~d} \mu_{n}(U)\right.}{\int_{\left|y_{1}\right| \leqslant r}\left|y_{1}\right|^{k} \mathrm{~d} y_{1}} \\
& \leqslant \int_{\mathcal{O}(n)}\left(\mathrm{M}_{k}^{U} f\right)(x) \mathrm{d} \mu_{n}(U)
\end{aligned}
$$

and we conclude by taking the supremum over all $r>0$.
It remains to check (4.2). By standard measure-theoretic arguments about classes of functions generating the Borel $\sigma$-algebra of $\mathbb{R}^{n}$, we can suppose that $g$ has the form $g(x)=g_{0}(|x|) g_{1}\left(x^{\prime}\right)$, with $x=|x| x^{\prime}$ and $x^{\prime} \in S^{n-1}$. By taking polar coordinates, we see that the left-hand side of (4.2) is equal to

$$
\frac{n}{r^{n}}\left(\int_{0}^{r} g_{0}(t) t^{n-1} \mathrm{~d} t\right)\left(\int_{S^{n-1}} g_{1}\left(y^{\prime}\right) \mathrm{d} \sigma_{n-1}\left(y^{\prime}\right)\right)
$$

where $\sigma_{n-1}$ is the invariant probability measure on $S^{n-1}$. The right-hand side is

$$
\frac{n}{r^{n}}\left(\int_{0}^{r} g_{0}(t) t^{n-1} \mathrm{~d} t\right)\left(\int_{\mathcal{O}(n)} \int_{S^{n-k-1}} g_{1}\left(U\left(y_{1}^{\prime}, 0\right)\right) \mathrm{d} \sigma_{n-k-1}\left(y_{1}^{\prime}\right) \mathrm{d} \mu_{n}(U)\right)
$$

Observe that for every $\theta_{0} \in S^{n-1}$, we have

$$
\int_{\mathcal{O}(n)} g_{1}\left(U \theta_{0}\right) \mathrm{d} \mu_{n}(U)=\int_{S^{n-1}} g_{1}(\theta) \mathrm{d} \sigma_{n-1}(\theta)
$$

since the left-hand side of this equality defines a probability measure on $S^{n-1}$, namely $\mathcal{B}_{S^{n-1}} \ni A \mapsto \int_{\mathcal{O}(n)} \mathbf{1}_{A}\left(U \theta_{0}\right) \mathrm{d} \mu_{n}(U)$, which is invariant under the left-action of $\mathcal{O}(n)$, hence equal to $\sigma_{n-1}$. We have therefore

$$
\begin{aligned}
\int_{\mathcal{O}(n)} \int_{S^{n-k-1}} g_{1} & \left(U\left(y_{1}^{\prime}, 0\right)\right) \mathrm{d} \sigma_{n-k-1}\left(y_{1}^{\prime}\right) \mathrm{d} \mu_{n}(U) \\
& =\int_{S^{n-k-1}}\left(\int_{\mathcal{O}(n)} g_{1}\left(U\left(y_{1}^{\prime}, 0\right)\right) \mathrm{d} \mu_{n}(U)\right) \mathrm{d} \sigma_{n-k-1}\left(y_{1}^{\prime}\right) \\
& =\int_{S^{n-k-1}}\left(\int_{S^{n-1}} g_{1}(\theta) \mathrm{d} \sigma_{n-1}(\theta)\right) \mathrm{d} \sigma_{n-k-1}\left(y_{1}^{\prime}\right) \\
& =\int_{S^{n-1}} g_{1}\left(y^{\prime}\right) \mathrm{d} \sigma_{n-1}\left(y^{\prime}\right)
\end{aligned}
$$

completing the proof.
Proof of Theorem 4.1. - Let $1<p \leqslant+\infty$. There is obviously nothing to do if $n \leqslant 2$. When $n \leqslant p /(p-1)$, the "bad" Vitali-bound $C(n)=3^{n}$ in the classical maximal inequality (ST) is less than a function of $p$ alone, namely $3^{p /(p-1)}$. We can therefore assume that both inequalities $n>p /(p-1)$ and $n \geqslant 3$ hold. We then write $n=(n-k)+k$ with $n-k=\lfloor\max (p /(p-1), 2)\rfloor+1$, and the result follows from Lemma 4.3 and Lemma 4.4 since with this choice, the bound $C(n-k, p)$ in Lemma 4.3 is now a function of $p$ alone.

### 4.2. Boundedness of the spherical maximal operator

In this section, we prove Theorem 4.2 following the approach of Rubio de Francia [68], see also Grafakos [39]. Let $n \geqslant 2$. The spherical maximal operator is expressed by

$$
(\mathcal{M} f)(x)=\sup _{r>0}\left|\left[m_{\sigma}(r \cdot) \widehat{f}(\cdot)\right]^{\vee}(x)\right|=\sup _{r>0}\left|\left(\sigma_{(r)} * f\right)(x)\right|, \quad x \in \mathbb{R}^{n}
$$

where $h^{\vee}(x)=\widehat{h}(-x)$ denotes the inverse Fourier transform of a function $h$, $m_{\sigma}$ is the Fourier transform of the uniform probability measure $\sigma$ on the unit
sphere $S^{n-1}$, and $\sigma_{(r)}$ is the dilated probability measure defined in (2.8). It is known that

$$
\begin{equation*}
m_{\sigma}(\xi)=\widehat{\sigma}(\xi)=(2 \pi|\xi|)^{-(n-2) / 2} \mathrm{~J}_{(n-2) / 2}(2 \pi|\xi|), \quad \xi \in \mathbb{R}^{n} \tag{4.3}
\end{equation*}
$$

with $\mathrm{J}_{\nu}$ the Bessel function of order $\nu$. This equality follows from the fact that the two functions $t \mapsto t^{-(n-2) / 2} \mathrm{~J}_{(n-2) / 2}(t)$ and

$$
t \mapsto F(t)=\int_{S^{n-1}} \mathrm{e}^{\mathrm{i} t x_{1}} \mathrm{~d} \sigma(x)=\frac{2 s_{n-2}}{s_{n-1}} \int_{0}^{1}\left(1-s^{2}\right)^{(n-3) / 2} \cos (s t) \mathrm{d} s
$$

are entire functions $g$ satisfying $g(0)=1$ and $t^{2}\left(g^{\prime \prime}(t)+g(t)\right)=-(n-1) t g^{\prime}(t)$.
We shall rely on the Littlewood-Paley theory, decomposing multipliers into dyadic pieces with localized frequencies. More precisely, we shall dominate $\mathcal{M}$ by a series of maximal operators $\sum_{\ell=0}^{+\infty} M_{K_{\ell}}$, where each kernel $K_{\ell}$ is radial with a well localized Fourier transform $m_{\ell}$. We establish that $M_{K_{\ell}}$ is of strong type when $p=2$ and of weak type $(1,1)$. Then, we get an $L^{p}$ bound for $M_{K_{\ell}}$ by interpolation, and the range of $p$ in Theorem 4.2 is chosen for making the series of bounds convergent. For the case $p=2$, we mainly use for $m_{\ell}$ both the decay at infinity and a support property, together with a precise upper bound for the $L^{2}\left(\mathbb{R}^{n}\right)$ norm of a related square function. When $p=1$, we invoke the usual Hardy-Littlewood theorem. Before giving the proof of Theorem 4.2, we introduce the dyadic decomposition of $m_{\sigma}=\widehat{\sigma}$.

Let $\varphi_{0}$ be a smooth radial function on $\mathbb{R}^{n}$ satisfying for every $\xi \in \mathbb{R}^{n}$ that

$$
\varphi_{0}(\xi)= \begin{cases}1 & \text { if }|\xi| \leqslant 1 \\ 0 & \text { if }|\xi| \geqslant 2\end{cases}
$$

Let $\psi(\xi)=\varphi_{0}(\xi)-\varphi_{0}(2 \xi)$ for $\xi$ in $\mathbb{R}^{n}$. This function $\psi$ is supported in the annulus $\{1 / 2 \leqslant|\xi| \leqslant 2\}$. For every integer $\ell \geqslant 1$ we define

$$
\varphi_{\ell}(\xi)=\varphi_{0}\left(2^{-\ell} \xi\right)-\varphi_{0}\left(2^{1-\ell} \xi\right)=\psi_{\left[2^{-l}\right]}(\xi), \quad \xi \in \mathbb{R}^{n}
$$

and for every $\ell \geqslant 0$, we consider the dyadic radial piece $m_{\ell}=\varphi_{\ell} m_{\sigma}$ associated to the multiplier $m_{\sigma}$. We can check that $\sum_{\ell=0}^{+\infty} \varphi_{\ell}=1$, thus $m_{\sigma}=\sum_{\ell=0}^{+\infty} m_{\ell}$. For every $\ell \geqslant 0$, we introduce the integrable kernel $K_{\ell}=$ $m_{\ell}^{\vee}=\varphi_{\ell}^{\vee} * \sigma$ and we set
$\left(M_{K_{\ell}} f\right)(x)=\sup _{r>0}\left|\left[m_{\ell}(r \cdot) \widehat{f}(\cdot)\right]^{\vee}(x)\right|=\sup _{r>0}\left|\left[\left(\varphi_{\ell}^{\vee}\right)_{(r)} * \sigma_{(r)} * f\right](x)\right|, \quad x \in \mathbb{R}^{n}$,
when $f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$. In particular, we have $M_{K_{0}} f=\sup _{r>0}\left|\left(\varphi_{0}^{\vee}\right)_{(r)} * \sigma_{(r)} * f\right|$ and

$$
M_{K_{\ell}} f=\sup _{r>0}\left|\left(\psi^{\vee}\right)_{\left(2^{-\ell} r\right)} * \sigma_{(r)} * f\right|, \quad \ell \geqslant 1
$$

For every $x \in \mathbb{R}^{n}$ and $r>0$, we see that $\left(\sigma_{(r)} * f\right)(x)=\sum_{\ell=0}^{+\infty}\left(\left(K_{\ell}\right)_{(r)} * f\right)(x)$ and we get the pointwise inequality

$$
\begin{equation*}
(\mathcal{M} f)(x) \leqslant \sum_{\ell=0}^{+\infty}\left(M_{K_{\ell}} f\right)(x) \tag{4.4}
\end{equation*}
$$

In a first subsection, we present some useful results on this type of maximal operators and associated square functions. Then, we shall prove that each $M_{K_{\ell}}$, for $\ell \geqslant 0$, is of strong type when $p=2$ and of weak type when $p=1$, and we give the proof of Theorem 4.2 in a third subsection.

### 4.2.1. Maximal operator and square function

Let $m(\xi)$ be a multiplier that is a bounded continuous function on $\mathbb{R}^{n}$, vanishing at 0 , with $|m(\xi)|=O(|\xi|)$ in a neighborhood of 0 . For $f$ in the Schwartz class $\mathcal{S}\left(\mathbb{R}^{n}\right)$ and for $x \in \mathbb{R}^{n}$, set

$$
\begin{aligned}
\left(g_{m} f\right)(x) & =\left(\int_{0}^{+\infty}\left|\left(T_{m_{[u]}} f\right)(x)\right|^{2} \frac{\mathrm{~d} u}{u}\right)^{1 / 2} \\
& =\left(\int_{0}^{+\infty}\left|\int_{\mathbb{R}^{n}} m(u \xi) \widehat{f}(\xi) \mathrm{e}^{2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} \xi\right|^{2} \frac{\mathrm{~d} u}{u}\right)^{1 / 2}
\end{aligned}
$$

We obtain the Littlewood-Paley function $g_{1}(f)$ of (2.3) when $m(\xi)=$ $2 \pi|\xi| \mathrm{e}^{-2 \pi|\xi|}$.

Lemma 4.5. - Assume that the multiplier $m(\xi)$ is a bounded function of $\xi \in \mathbb{R}^{n}$, supported in an annulus of the form $\{a \leqslant|\xi| \leqslant r a\}$, $a>0$ and $r>1$. For every function $f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$, one has that

$$
\left\|g_{m} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)} \leqslant \sqrt{\ln r}\|m\|_{L^{\infty}\left(\mathbb{R}^{n}\right)}\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)}
$$

Proof. - According to the Fubini theorem, followed by Parseval, Fubini again and setting finally $v=u|\xi|$, we have

$$
\begin{aligned}
& \int_{\mathbb{R}^{n}}\left|\left(g_{m} f\right)(x)\right|^{2} \mathrm{~d} x \\
&=\int_{0}^{+\infty}\left\|T_{m_{[u]}} f\right\|_{2}^{2} \frac{\mathrm{~d} u}{u}=\int_{0}^{+\infty} \int_{\mathbb{R}^{n}}|m(u \xi)|^{2}|\widehat{f}(\xi)|^{2} \frac{\mathrm{~d} u}{u} \mathrm{~d} \xi \\
& \leqslant\|m\|_{\infty}^{2} \int_{\mathbb{R}^{n}}\left(\int_{a}^{r a} \frac{\mathrm{~d} v}{v}\right)|\widehat{f}(\xi)|^{2} \mathrm{~d} \xi=\|m\|_{\infty}^{2} \ln (r)\|f\|_{2}^{2}
\end{aligned}
$$

Lemma 4.6. - Assume that $m(\xi)$ is of class $C^{1}$ on $\mathbb{R}^{n}$ and vanishes outside a compact subset of $\mathbb{R}^{n} \backslash\{0\}$. For every $t>0$ and $f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$, we have

$$
\left|\int_{\mathbb{R}^{n}} m(t \xi) \widehat{f}(\xi) \mathrm{e}^{2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} \xi\right|^{2} \leqslant 2\left(g_{m} f\right)(x)\left(g_{m^{*}} f\right)(x), \quad x \in \mathbb{R}^{n}
$$

where we have set $m^{*}(\xi)=\xi \cdot \nabla m(\xi)$ for every $\xi \in \mathbb{R}^{n}$.
Proof. - For each $s \geqslant 0$ let us set

$$
\left(g_{m, s} f\right)(x)=\left(T_{m_{[s]}} f\right)(x)=\int_{\mathbb{R}^{n}} m(s \xi) \widehat{f}(\xi) \mathrm{e}^{2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} \xi, \quad x \in \mathbb{R}^{n}
$$

We note that
$s \frac{\mathrm{~d}}{\mathrm{~d} s}\left(g_{m, s} f\right)(x)=\int_{\mathbb{R}^{n}} s \xi \cdot \nabla m(s \xi) \widehat{f}(\xi) \mathrm{e}^{2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} \xi=\int_{\mathbb{R}^{n}} m^{*}(s \xi) \widehat{f}(\xi) \mathrm{e}^{2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} \xi$, which allows us to see this quantity as $\left(g_{m^{*}, s} f\right)(x)$. Since $m$ vanishes in a neighborhood of 0 , one has $\left(g_{m, 0} f\right)(x)=0$, thus

$$
\begin{aligned}
\left|\left(g_{m, t} f\right)(x)\right|^{2} & =\int_{0}^{t} \frac{\mathrm{~d}}{\mathrm{~d} s}\left|\left(g_{m, s} f\right)(x)\right|^{2} \mathrm{~d} s \\
& =2 \operatorname{Re} \int_{0}^{t} \overline{\left(g_{m, s} f\right)(x)} s \frac{\mathrm{~d}}{\mathrm{~d} s}\left(g_{m, s} f\right)(x) \frac{\mathrm{d} s}{s} \\
& =2 \operatorname{Re} \int_{0}^{t} \overline{\left(g_{m, s} f\right)(x)}\left(g_{m^{*}, s} f\right)(x) \frac{\mathrm{d} s}{s}
\end{aligned}
$$

By Cauchy-Schwarz, and bounding the integral on $[0, t]$ by the integral on $[0,+\infty)$, we obtain that

$$
\begin{aligned}
\left|\left(g_{m, t} f\right)(x)\right|^{2} & \leqslant 2\left(\int_{0}^{+\infty}\left|\left(g_{m, s} f\right)(x)\right|^{2} \frac{\mathrm{~d} s}{s}\right)^{1 / 2}\left(\int_{0}^{+\infty}\left|\left(g_{m^{*}, s} f\right)(x)\right|^{2} \frac{\mathrm{~d} s}{s}\right)^{1 / 2} \\
& =2\left(g_{m} f\right)(x)\left(g_{m^{*}} f\right)(x)
\end{aligned}
$$

Lemma 4.7. - Let $K$ be an integrable kernel on $\mathbb{R}^{n}$. Suppose that $m$, the Fourier transform of $K$, is of class $C^{1}$ on $\mathbb{R}^{n}$ and supported in an annulus of the form $\{a \leqslant|\xi| \leqslant r a\}, a>0$ and $r>1$. For every function $f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$, one has that

$$
\begin{aligned}
\left\|M_{K} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}^{2}=\| \sup _{t>0} \mid K_{(t)} & * f \mid \|_{L^{2}\left(\mathbb{R}^{n}\right)}^{2} \\
& \leqslant 2 \ln (r)\|m\|_{L^{\infty}\left(\mathbb{R}^{n}\right)}\left\|m^{*}\right\|_{L^{\infty}\left(\mathbb{R}^{n}\right)}\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)}^{2}
\end{aligned}
$$

where $m^{*}(\xi)=\xi \cdot \nabla m(\xi)$ for $\xi \in \mathbb{R}^{n}$.
Proof. - By Lemma 4.6, we have for every $x \in \mathbb{R}^{n}$ and $t>0$ that

$$
\left|\left(K_{(t)} * f\right)(x)\right|^{2}=\left|\int_{\mathbb{R}^{n}} m(t \xi) \widehat{f}(\xi) \mathrm{e}^{2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} \xi\right|^{2} \leqslant 2\left(g_{m} f\right)(x)\left(g_{m^{*}} f\right)(x)
$$

This upper bound is independent of $t$, thus

$$
\left(\left(M_{K} f\right)(x)\right)^{2} \leqslant 2\left(g_{m} f\right)(x)\left(g_{m^{*}} f\right)(x)
$$

and by Cauchy-Schwarz we get

$$
\left\|M_{K} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}^{2} \leqslant 2\left\|g_{m} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}\left\|g_{m^{*}} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}
$$

According to Lemma 4.5, we conclude that

$$
\left\|M_{K} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}^{2} \leqslant 2 \ln (r)\|m\|_{L^{\infty}\left(\mathbb{R}^{n}\right)}\left\|m^{*}\right\|_{L^{\infty}\left(\mathbb{R}^{n}\right)}\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)}^{2}
$$

The following proposition is nearly obvious.
Proposition 4.8. - Let $K \in \mathcal{S}\left(\mathbb{R}^{n}\right)$ be a radial kernel. For every $p$ in $(1,+\infty]$, the maximal operator $M_{K}$ is bounded on $L^{p}\left(\mathbb{R}^{n}\right)$.

One also gets the weak type $(1,1)$ for $M_{K}$, but we shall not use it.
Proof. - Since $K$ is a Schwartz radial function, we can find an integrable function $\Omega$, radial and radially decreasing, such that $|K| \leqslant \Omega$. It implies that

$$
\sup _{r>0}\left|\left[K_{(r)} * f\right](x)\right| \leqslant \sup _{r>0}\left(\Omega_{(r)} *|f|\right)(x), \quad x \in \mathbb{R}^{n}
$$

and $\Omega$ being radial and radially decreasing, we classically have

$$
\begin{equation*}
\sup _{r>0}\left(\Omega_{(r)} *|f|\right)(x) \leqslant\|\Omega\|_{L^{1}\left(\mathbb{R}^{n}\right)}(\mathrm{M} f)(x), \quad x \in \mathbb{R}^{n} \tag{4.5}
\end{equation*}
$$

By Theorem 0.1, the usual maximal theorem for M , we get the conclusion.
For proving (4.5), it suffices to show that

$$
\begin{equation*}
|(\Omega * f)(x)| \leqslant\|\Omega\|_{L^{1}\left(\mathbb{R}^{n}\right)}(\mathrm{M} f)(x), \quad x \in \mathbb{R}^{n} . \tag{4.6}
\end{equation*}
$$

Suppose that $\Omega \leqslant 1$ for simplicity, and consider for each integer $k \geqslant 1$ the set

$$
A_{k}=\left\{x \in \mathbb{R}^{n}: \Omega(x)>2^{-k}\right\}
$$

This set $A_{k}$ is a Euclidean ball, and if we define $g=\sum_{k \geqslant 1} 2^{-k} \mathbf{1}_{A_{k}}$, we can check that $g / 2 \leqslant \Omega \leqslant g$. We rewrite $g$ as

$$
g=\sum_{k \geqslant 1} a_{k} \frac{\mathbf{1}_{A_{k}}}{\left|A_{k}\right|}
$$

with $a_{k}>0$ for every $k \geqslant 1$. Since $\Omega$ is integrable, $g$ is also integrable and

$$
\sum_{k \geqslant 0} a_{k}=\int_{\mathbb{R}^{n}} g(x) \mathrm{d} x \leqslant 2 \int_{\mathbb{R}^{n}} \Omega(x) \mathrm{d} x=2\|\Omega\|_{L^{1}\left(\mathbb{R}^{n}\right)} .
$$

We have for every $x \in \mathbb{R}^{n}$ that

$$
\begin{aligned}
|(\Omega * f)(x)| & \leqslant(g *|f|)(x)=\sum_{k \geqslant 1} \frac{a_{k}}{\left|A_{k}\right|} \int_{x+A_{k}}|f(y)| \mathrm{d} y \\
& \leqslant\left(\sum_{k \geqslant 0} a_{k}\right)(\mathrm{M} f)(x) \leqslant 2\|\Omega\|_{L^{1}\left(\mathbb{R}^{n}\right)}(\mathrm{M} f)(x) .
\end{aligned}
$$

The inequality with constant 1 can be reached by refining the partition, replacing the values $2^{-k}$ by $(1+\varepsilon)^{-k}$, with $\varepsilon>0$ tending to 0 . One can also give a direct proof involving integration by parts, or the Fubini theorem and level sets of $\Omega$.

### 4.2.2. Strong and weak type results for $M_{K_{\ell}}, \ell \geqslant 1$

We begin with the strong type result, when $p=2$.
Proposition 4.9. - For every integer $\ell \geqslant 1$ and every $f \in L^{2}\left(\mathbb{R}^{n}\right)$ one has that

$$
\left\|M_{K_{\ell}} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)} \leqslant C(n) 2^{-\ell(n-2) / 2}\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)}
$$

where $C(n)$ is a constant independent of $\ell$.
Proof. - For each $\ell \geqslant 1$, the multiplier $m_{\ell}=\widehat{K_{\ell}}$ is $C^{1}$, supported in the annulus

$$
I_{\ell}=\left\{\xi \in \mathbb{R}^{n}: 2^{\ell-1} \leqslant|\xi| \leqslant 2^{\ell+1}\right\}
$$

Applying Lemma 4.7 to $K_{\ell}$, with $m_{\ell}^{*}(\xi)=\xi \cdot \nabla m_{\ell}(\xi)$ and $r=4$, we obtain

$$
\left\|M_{K_{\ell}} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}^{2} \leqslant 2 \ln (4)\left\|m_{\ell}\right\|_{L^{\infty}\left(\mathbb{R}^{n}\right)}\left\|m_{\ell}^{*}\right\|_{L^{\infty}\left(\mathbb{R}^{n}\right)}\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)}^{2}
$$

The desired result will be consequence of the inequalities

$$
\begin{equation*}
\left\|m_{\ell}\right\|_{L^{\infty}\left(\mathbb{R}^{n}\right)} \leqslant C_{1}(n) 2^{-\ell(n-1) / 2}, \quad\left\|m_{\ell}^{*}\right\|_{L^{\infty}\left(\mathbb{R}^{n}\right)} \leqslant C_{2}(n) 2^{-\ell(n-3) / 2} \tag{4.7}
\end{equation*}
$$

that we establish now, with $C_{1}(n)$ and $C_{2}(n)$ independent of $\ell$. Thanks to well-known properties of Bessel functions (see for instance [2, p. 238]), we have

$$
\begin{equation*}
\sup _{t \geqslant 1} t^{1 / 2}\left|\mathrm{~J}_{\alpha}(t)\right|<+\infty, \quad \text { and } \frac{\mathrm{d}}{\mathrm{~d} t} \mathrm{~J}_{\alpha}(t)=\frac{1}{2}\left(\mathrm{~J}_{\alpha-1}(t)-\mathrm{J}_{\alpha+1}(t)\right) \tag{4.8}
\end{equation*}
$$

The first property follows from the fact that $u_{\alpha}(t)=\sqrt{t} \mathrm{~J}_{\alpha}(t)$ satisfies a differential equation $u_{\alpha}^{\prime \prime}(t)+\left(1+\kappa_{\alpha} t^{-2}\right) u_{\alpha}(t)=0$ for $t>0$, hence $v_{\alpha}(t):=$ $\left(u_{\alpha}(t)^{2}+u_{\alpha}^{\prime}(t)^{2}\right) / 2$ satisfies $v_{\alpha}^{\prime}(t)=-\kappa_{\alpha} t^{-2} u_{\alpha}(t) u_{\alpha}^{\prime}(t) \leqslant\left|\kappa_{\alpha}\right| t^{-2} v_{\alpha}(t)$, yielding $v_{\alpha}(t) \leqslant \mathrm{e}^{\left|\kappa_{\alpha}\right|} v_{\alpha}(1)$ for every $t \geqslant 1$. The second property can be checked on the coefficients of the power series $\sum_{m \geqslant 0}(-1)^{m}(m!\Gamma(m+\alpha+$ $1))^{-1}(t / 2)^{2 m}$ of $t^{-\alpha} \mathrm{J}_{\alpha}(t)$, and when $\alpha=n \in \mathbb{N}$, it is even simpler to see it on the integral expression $2 \pi \mathrm{~J}_{n}(t)=\int_{0}^{2 \pi} \mathrm{e}^{\mathrm{i}(t \sin s-n s)} \mathrm{d} s$.

Since $m_{\ell}$ and $m_{\ell}^{*}$ are supported in the annulus $I_{\ell}$, we need only bound $m_{\ell}(\xi)$ and $m_{\ell}^{*}(\xi)$ when $1 \leqslant 2^{\ell-1} \leqslant|\xi| \leqslant 2^{\ell+1}$ (we have $\ell \geqslant 1$ ). We then obtain (4.7) by recalling (4.3) and by applying (4.8) to $t=2 \pi|\xi|>1$, which give that

$$
\left|m_{\ell}(\xi)\right| \leqslant c_{1}(n)|\xi|^{-n / 2+1 / 2} \quad \text { and } \quad\left|m_{\ell}^{*}(\xi)\right| \leqslant c_{2}(n)|\xi|^{-n / 2+3 / 2}
$$

We state in the next proposition a crucial weak type estimate for $M_{K_{\ell}}$.
Proposition 4.10. - Let $\ell \geqslant 1$. For all $f \in L^{1}\left(\mathbb{R}^{n}\right)$ and every $\lambda>0$, one has that

$$
\left|\left\{x \in \mathbb{R}^{n}:\left(M_{K_{\ell}} f\right)(x)>\lambda\right\}\right| \leqslant C(n) \frac{2^{\ell}}{\lambda}\|f\|_{L^{1}\left(\mathbb{R}^{n}\right)}
$$

where $C(n)$ is a constant independent of $\ell$ and $\lambda$.
Proof. - We claim that it is enough to prove that for each $\ell \geqslant 1$, we have

$$
\begin{equation*}
\left|K_{\ell}(x)\right| \leqslant C(n) \frac{2^{\ell}}{(1+|x|)^{n+1}}, \quad x \in \mathbb{R}^{n} \tag{4.9}
\end{equation*}
$$

Indeed, since $(1+|x|)^{-n-1}$ is radial, radially decreasing and integrable, we will have for all $x \in \mathbb{R}^{n}$, as in (4.6), that

$$
\sup _{r>0}\left|\left[\left(K_{\ell}\right)_{(r)} * f\right](x)\right| \leqslant \widetilde{C}(n) 2^{\ell}(\mathrm{M} f)(x)
$$

The result of Proposition 4.10 follows then from the weak estimate in Theorem 0.1 , the standard maximal theorem. We now turn to the proof of (4.9). We want a bound for $K_{\ell}=\varphi_{\ell}^{\vee} * \sigma$, for $\ell \geqslant 1$, where $\sigma$ is the uniform probability measure on $S^{n-1}$ and $\varphi_{\ell}^{\vee}=\left(\psi^{\vee}\right)_{\left(2^{-\ell}\right)}$. Since $\psi^{\vee}$ belongs to the Schwartz class, we can bound $\left|\psi^{\vee}\right|$ by a multiple $c_{n} g$ of the radial and radially decreasing integrable function $g(x)=(1+|x|)^{-n-1}$. In order to bound $K_{\ell}$, we shall prove that

$$
\begin{aligned}
c_{n}^{-1}\left|K_{\ell}(x)\right| \leqslant\left(g_{\left(2^{-\ell}\right)} * \sigma\right)(x) & =\int_{S^{n-1}} g_{\left(2^{-\ell}\right)}(x-z) \mathrm{d} \sigma(z) \\
& \leqslant C(n) 2^{\ell}(1+|x|)^{-n-1} .
\end{aligned}
$$

This is easy when $|x|>2$, because for each $z$ in $S^{n-1}$, we have then $|x-z| \geqslant$ $|x|-1 \geqslant|x| / 2$ and $1+|x| \leqslant 2|x|$. Recalling $g_{\left(2^{-\ell}\right)}(y)=2^{n \ell} g\left(2^{\ell} y\right)$, we get

$$
\begin{aligned}
G_{\ell}(x):=\left(g_{\left(2^{-\ell}\right)} * \sigma\right)(x) & \leqslant \max _{z \in S^{n-1}} g_{\left(2^{-\ell}\right)}(x-z) \leqslant 2^{n \ell}\left(1+2^{\ell}|x| / 2\right)^{-n-1} \\
& \leqslant 2^{n \ell 2^{-(\ell-1)(n+1)}|x|^{-n-1}=2^{n+1-\ell}|x|^{-n-1}} \\
& \leqslant 2^{2 n+1}(1+|x|)^{-n-1},
\end{aligned}
$$

even better than required. Suppose now that $|x| \leqslant 2$. It is enough to prove that $G_{\ell}(x) \leqslant C(n) 2^{\ell}$, since we have $1+|x| \leqslant 3$ in this second case, hence it
will follow that $C(n) 2^{\ell} \leqslant\left[C(n) 3^{n+1}\right] 2^{\ell}(1+|x|)^{-n-1}$. For $y \in \mathbb{R}^{n}$, we write $y=(v, t)$ with $v \in \mathbb{R}^{n-1}$ and $t$ real. By the rotational invariance, we may restrict the study to $x=(0, s), s \geqslant 0$. We write each $z \in S^{n-1}$ as $z=(v, t)$, and thus $x-z=(-v, s-t)$. Let $\pi_{0}$ be the orthogonal projection of $\mathbb{R}^{n}$ onto the hyperplane of vectors $(w, 0), w \in \mathbb{R}^{n-1}$. Since $g_{\left(2^{-\ell}\right)}$ is radial and radially decreasing, we see that $g_{\left(2^{-\ell}\right)}(x-z) \leqslant g_{\left(2^{-\ell}\right)}\left(\pi_{0}(x-z)\right)=g_{\left(2^{-\ell}\right)}(-v, 0)$. This yields

$$
\begin{aligned}
G_{\ell}(x)=G_{\ell}(0, s) & =\int_{S^{n-1}} g_{\left(2^{-\ell}\right)}(x-z) \mathrm{d} \sigma(z) \\
& \leqslant \int_{S^{n-1}} g_{\left(2^{-\ell}\right)}\left(\pi_{0}(x-z)\right) \mathrm{d} \sigma(z) \\
& =\int_{\mathbb{R}^{n-1}} g_{\left(2^{-\ell}\right)}(-v, 0) \mathrm{d} \nu(v)
\end{aligned}
$$

where $\nu$ is the projection on $\mathbb{R}^{n-1}$ of the probability measure $\sigma$. We have that

$$
\mathrm{d} \nu(v)=\frac{2}{s_{n-1}} \frac{\mathbf{1}_{\{|v|<1\}}}{\sqrt{1-|v|^{2}}} \mathrm{~d} v=C(n) \frac{\mathbf{1}_{\{|v|<1\}}}{\sqrt{1-|v|^{2}}} \mathrm{~d} v
$$

where $s_{n-1}$ is the measure of $S^{n-1}$ recalled in (1.34). We cut the integral with respect to $\nu$ into two parts, according to $|v|<1 / 2$ or not. In the part $E_{1}$ corresponding to $|v|<1 / 2$, we have $1-|v|^{2} \geqslant 3 / 4$, hence

$$
E_{1} \leqslant \sqrt{\frac{4}{3}} C(n) \int_{|v|<1 / 2} g_{\left(2^{-\ell}\right)}(v, 0) \mathrm{d} v \leqslant 2 C(n) \int_{\mathbb{R}^{n-1}} g_{\left(2^{-\ell}\right)}(v, 0) \mathrm{d} v
$$

We are integrating on $\mathbb{R}^{n-1}$ the function $g_{\left(2^{-\ell}\right)}$ that is normalized for a change of variable in dimension $n$. This implies that

$$
E_{1} \leqslant 2 C(n) 2^{\ell} 2^{(n-1) \ell} \int_{\mathbb{R}^{n-1}} g\left(2^{\ell} v, 0\right) \mathrm{d} v=2 C(n) 2^{\ell} \int_{\mathbb{R}^{n-1}} g(u, 0) \mathrm{d} u
$$

a bound of the expected form. In the second case, we have $|v|>1 / 2$ and

$$
g_{\left(2^{-\ell}\right)}(v, 0)=2^{n \ell}\left(1+2^{\ell}|v|\right)^{-n-1} \leqslant 2^{n \ell} 2^{-(\ell-1)(n+1)} \leqslant 2^{n}
$$

It follows that the integral $E_{2}$ limited to $|v|>1 / 2$, with respect to the probability measure $\nu$, is bounded by a function of $n$.

### 4.2.3. Conclusion

Proof of Theorem 4.2. - Thanks to the results of the previous subsection, the proof is easy. Using the Marcinkiewicz theorem (see Zygmund [85, Chap. XII], or [64, Theorem 5.60]), we shall interpolate between the weak type $(1,1)$ and the strong type (2,2). We apply Proposition 4.9, Proposition 4.10 in $\mathbb{R}^{N}$ and interpolation with parameter $\theta=2-2 / p$, where
$1<p \leqslant 2$. For all $\ell \geqslant 1$ and all $f \in L^{p}\left(\mathbb{R}^{N}\right)$, since the chosen interpolation parameter $\theta$ satisfies $(1-\theta) / 1+\theta / 2=1 / p$, we have

$$
\left\|M_{K_{\ell}} f\right\|_{L^{p}\left(\mathbb{R}^{N}\right)} \leqslant \kappa(1,2, p) C(N)\left(2^{\ell}\right)^{-1+2 / p}\left(2^{-\ell(N-2) / 2}\right)^{2-2 / p}\|f\|_{L^{p}\left(\mathbb{R}^{N}\right)}
$$

where $\kappa(1,2, p)$ is independent of $N$ and $\ell$. We have thus obtained that

$$
\left\|M_{K_{\ell}} f\right\|_{L^{p}\left(\mathbb{R}^{N}\right)} \leqslant C^{\prime}(N, p) 2^{\ell[N / p-(N-1)]}\|f\|_{L^{p}\left(\mathbb{R}^{N}\right)}
$$

For $p>N /(N-1)$, the series $\sum_{\ell \geqslant 1} 2^{\ell[N / p-(N-1)]}$ converges. Moreover, we know by Proposition 4.8 that $M_{K_{0}}$ maps $L^{p}\left(\mathbb{R}^{N}\right)$ to itself for all $1<p<+\infty$. Therefore, in view of (4.4), we obtain that $\mathcal{M}$ is bounded on $L^{p}\left(\mathbb{R}^{N}\right)$ for every real number $p$ such that $N /(N-1)<p \leqslant 2$. For $p>2$, we proceed by interpolation between the $L^{2}\left(\mathbb{R}^{N}\right)$ case and the trivial $L^{\infty}\left(\mathbb{R}^{N}\right)$ case.

## 5. The $L^{2}$ result of Bourgain

In an article published in 1986, Bourgain has generalized the $L^{2}$ case of the Stein result presented in Section 4. This $L^{2}$ case for Euclidean balls only required Proposition 4.9 and the "method of rotations". The maximal operator $\mathrm{M}_{C}$ associated to a symmetric convex body $C$ was defined in (0.3.M).

ThEOREM 5.1 (Bourgain [9]). - There exists a universal constant $\kappa_{2}$ such that for every integer $n \geqslant 1$ and every symmetric convex body $C \subset \mathbb{R}^{n}$, one has

$$
\forall f \in L^{2}\left(\mathbb{R}^{n}\right), \quad\left\|\mathrm{M}_{C} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)} \leqslant \kappa_{2}\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)}
$$

The rest of this section is devoted to the proof of this maximal theorem, together with the description of the general framework concerning maximal functions associated to convex sets. We shall in particular establish some geometric inequalities for log-concave distributions that will be applied in the subsequent sections.

### 5.1. The general setting

Let $C$ be a symmetric convex body in $\mathbb{R}^{n}$. Throughout these Notes, we let $K_{C}$ be the density of the uniform probability measure $\mu_{C}$ on $C$, and $m_{C}$ denotes the Fourier transform of $K_{C}$ or of $\mu_{C}$. Hence, we have

$$
K_{C}(x)=\frac{1}{|C|} \mathbf{1}_{C}(x), \quad \mathrm{d} \mu_{C}(x)=K_{C}(x) \mathrm{d} x, \quad m_{C}(\xi)=\widehat{K_{C}}(\xi)=\widehat{\mu_{C}}(\xi)
$$

for all $x, \xi \in \mathbb{R}^{n}$. Notice that $K_{\lambda C}=\left(K_{C}\right)_{(\lambda)}$ and $m_{\lambda C}(\xi)=m_{C}(\lambda \xi)$ for each $\lambda>0$ and $\xi \in \mathbb{R}^{n}$. We already know that the maximal operator $\mathrm{M}_{C}$
acts boundedly on $L^{p}\left(\mathbb{R}^{n}\right), 1<p \leqslant+\infty$, but the bounds we have so far depend on $n$.

This $L^{p}$ result comes from the weak type estimate (0.4) given by the Vitali covering lemma. Except for the value of the constant, it is clear that this weak type $(1,1)$ result for $\mathrm{M}_{C}$ is optimal, as we can see by taking for $f$ the indicator $\mathbf{1}_{C}$ of the symmetric convex body $C \subset \mathbb{R}^{n}$. Let $C$ have volume 1, so that $\|f\|_{1}=1$. For any given $r>0$ and $x \in r C$, we see that $x+(r+1) C$ contains $C$, therefore

$$
\left(\mathrm{M}_{C} f\right)(x) \geqslant|(r+1) C|^{-1} \int_{x+(r+1) C} \mathbf{1}_{C}(y) \mathrm{d} y=|(r+1) C|^{-1}=(r+1)^{-n}
$$

and $\left\{\mathrm{M}_{C} f \geqslant(r+1)^{-n}\right\} \supset r C$. Every value $c$ in the interval $\left(0,2^{-n}\right]$ can be written as $c=(r+1)^{-n}$ for some $r \geqslant 1$, hence

$$
\forall c \in\left(0,2^{-n}\right], \quad\left|\left\{\mathrm{M}_{C} f \geqslant c\right\}\right| \geqslant|r C|=\frac{(r+1)^{-n}}{c} r^{n} \geqslant \frac{2^{-n}}{c}
$$

The maximal function $\mathrm{M}_{C} \mathbf{1}_{C}$ is not integrable. It belongs to the space $L^{1, \infty}\left(\mathbb{R}^{n}\right)$, the so-called weak- $L^{1}$ space, and nothing better: any bounded radial and radially decreasing function belonging to $L^{1, \infty}\left(\mathbb{R}^{n}\right)$ is smaller than a multiple of $\mathrm{M}_{C} \mathbf{1}_{C}$.

The maximal function $\mathrm{M}_{C} f$ is given by $\mathrm{M}_{C} f=\sup _{t>0}\left(K_{C}\right)_{(t)} *|f|$, where $\left(K_{C}\right)_{(t)}$ is the dilate from (2.7). More generally, let $K$ be a probability density on $\mathbb{R}^{n}$, resp. an integrable kernel $K$. We define the maximal function $\mathrm{M}_{K}$ or $M_{K}$ by

$$
\mathrm{M}_{K} f=\sup _{t>0} K_{(t)} *|f|, \quad \text { resp. } \quad M_{K} f=\sup _{t>0}\left|K_{(t)} * f\right|
$$

If $A$ is linear and bijective on $\mathbb{R}^{n}$, we can see that the maximal operators $\mathrm{M}_{C}$ and $\mathrm{M}_{A C}$ have the same norm on $L^{p}\left(\mathbb{R}^{n}\right)$. For a function $f$ on $\mathbb{R}^{n}$ we define $f_{(A)}$ by

$$
\forall x \in \mathbb{R}^{n}, \quad f_{(A)}(x)=|\operatorname{det} A|^{-1} f\left(A^{-1} x\right)
$$

We have $|f|_{(A)}=\left|f_{(A)}\right|,\left(\sup _{i} f_{i}\right)_{(A)}=\sup _{i}\left(f_{i}\right)_{(A)}$, and $(f * g)_{(A)}=f_{(A)} *$ $g_{(A)}$ since
$\int_{\mathbb{R}^{n}}|\operatorname{det} A|^{-2} f\left(A^{-1}(x-y)\right) g\left(A^{-1} y\right) \mathrm{d} y=|\operatorname{det} A|^{-1} \int_{\mathbb{R}^{n}} f\left(A^{-1} x-z\right) g(z) \mathrm{d} z$.
It is clear that $\left(f_{(A)}\right)_{(t)}=f_{(t A)}=\left(f_{(t)}\right)_{(A)}$. If $S_{A}$ is the mapping $f \mapsto f_{(A)}$, then $S_{A, p}:=|\operatorname{det} A|^{1 / q} S_{A}$, with $q$ conjugate to $p$, is an onto isometry of $L^{p}\left(\mathbb{R}^{n}\right)$.

The density $K_{A C}$ is equal to $\left(K_{C}\right)_{(A)}$. For every integrable kernel $K$ on $\mathbb{R}^{n}$, we see now that $K$ and $K_{(A)}$ produce maximal functions that are
conjugate by the isometry $S_{A, p}$ of $L^{p}\left(\mathbb{R}^{n}\right)$, and have therefore the same norm on $L^{p}\left(\mathbb{R}^{n}\right)$. We have

$$
\begin{aligned}
M_{K_{(A)}} f_{(A)} & =\sup _{t>0}\left|\left(K_{(A)}\right)_{(t)} * f_{(A)}\right|=\sup _{t>0}\left|\left(K_{(t)}\right)_{(A)} * f_{(A)}\right| \\
& =\sup _{t>0}\left|\left(K_{(t)} * f\right)_{(A)}\right|=\left(M_{K} f\right)_{(A)}
\end{aligned}
$$

It follows that $M_{K_{(A)}} \circ S_{A, p}=S_{A, p} \circ M_{K}$. This remark allows us to assume that $C$ is in isotropic position: one says that a symmetric convex body $C$ is in isotropic position if the quadratic form

$$
Q_{C}: \xi \mapsto Q_{C}(\xi)=\int_{C}(\xi \cdot x)^{2} \mathrm{~d} x, \quad \xi \in \mathbb{R}^{n}
$$

is a multiple of the square $\xi \mapsto|\xi|^{2}$ of the Euclidean norm on $\mathbb{R}^{n}$. Since $Q_{C}$ is positive definite for every symmetric convex body $C$, we can bring it to the form $\xi \mapsto \lambda|\xi|^{2}, \lambda>0$, by a suitable linear change of coordinates. For an isotropic symmetric convex set $C_{0}$ of volume 1 , one defines the isotropy constant $L\left(C_{0}\right)$ by

$$
L\left(C_{0}\right)^{2}=\int_{C_{0}}\left(\mathbf{e}_{1} \cdot x\right)^{2} \mathrm{~d} x, \quad \text { and one has then } \int_{C_{0}}(\xi \cdot x)^{2} \mathrm{~d} x=L\left(C_{0}\right)^{2}|\xi|^{2}
$$

for every $\xi \in \mathbb{R}^{n}$. For $C_{*}$ isotropic of the form $C_{*}=r C_{0}, r>0$, we get $\left|C_{*}\right|=r^{n}$ and for every $\xi \in \mathbb{R}^{n}$, we have

$$
\begin{align*}
\int_{\mathbb{R}^{n}}(\xi \cdot x)^{2} K_{C_{*}}(x) \mathrm{d} x & =\frac{1}{\left|C_{*}\right|} \int_{C_{*}}(\xi \cdot x)^{2} \mathrm{~d} x=r^{-n} \int_{C_{0}}(\xi \cdot r u)^{2} r^{n} \mathrm{~d} u  \tag{5.1}\\
& =r^{2} L\left(C_{0}\right)^{2}|\xi|^{2}=\left|C_{*}\right|^{2 / n} L\left(C_{0}\right)^{2}|\xi|^{2}
\end{align*}
$$

Let $A$ linear and invertible put $C_{*}$ in another isotropic position $A C_{*}$, so that $Q_{A C_{*}}(\xi)=\lambda|\xi|^{2}$ for some $\lambda>0$ and all $\xi \in \mathbb{R}^{n}$. Letting $\nu=\lambda\left|A C_{*}\right|^{-1}$ we get

$$
\begin{aligned}
\nu|\xi|^{2}=\int_{\mathbb{R}^{n}}(\xi \cdot y)^{2} K_{A C_{*}}(y) \mathrm{d} y & =\int_{\mathbb{R}^{n}}(\xi \cdot A x)^{2} K_{C_{*}}(x) \mathrm{d} x \\
& =\left|C_{*}\right|^{2 / n} L\left(C_{0}\right)^{2}\left|A^{T} \xi\right|^{2}
\end{aligned}
$$

hence $A$ is a multiple $\rho U$ of an isometry $U,|\operatorname{det} A|=\rho^{n}$ and $\nu=$ $\left|C_{*}\right|^{2 / n} L\left(C_{0}\right)^{2} \rho^{2}=\left|A C_{*}\right|^{2 / n} L\left(C_{0}\right)^{2}$, thus $\left|A C_{*}\right|^{-2 / n} \int_{\mathbb{R}^{n}}(\theta \cdot y)^{2} K_{A C_{*}}(y) \mathrm{d} y=$ $L\left(C_{0}\right)^{2}$ for every $\theta \in S^{n-1}$.

When $C$ is isotropic, it follows that $L(C):=L\left(C_{0}\right)$ is well defined by

$$
\begin{align*}
L(C)^{2} & =|C|^{-2 / n} \int_{\mathbb{R}^{n}}(\theta \cdot x)^{2} K_{C}(x) \mathrm{d} x  \tag{5.2}\\
& =|C|^{-1-2 / n} \int_{C}(\theta \cdot x)^{2} \mathrm{~d} x, \quad \theta \in S^{n-1}
\end{align*}
$$

A well-known open question (see [58]) is to decide whether the isotropy constant is bounded above by a universal constant valid for all symmetric convex bodies and every $n$. The best upper bound that is known so far, due to Klartag [49] improving Bourgain [12], is $L(C) \leqslant \kappa n^{1 / 4}$ in dimension $n$. It is known that $L(C)$ is bounded below by a universal constant. However, neither this known fact nor the unsolved problem will interfere with the treatment of the maximal function problem.

Clearly, $K_{C}$ and $\left(K_{C}\right)_{(\lambda)}$ have the same maximal function for every $\lambda>0$, so we can choose any multiple among isotropic positions of $C$. Here, we do not follow Bourgain [9] who chooses the isotropic position of volume 1, we prefer the isotropic position such that $\mu_{C}$ has covariance matrix $\mathrm{I}_{n}$. We thus assume that

$$
\begin{equation*}
\forall \theta \in S^{n-1}, \quad \int_{\mathbb{R}^{n}}(\theta \cdot x)^{2} \mathrm{~d} \mu_{C}(x)=\frac{1}{|C|} \int_{C}(\theta \cdot x)^{2} \mathrm{~d} x=1 \tag{5.3}
\end{equation*}
$$

This means that the one-dimensional marginals of $\mu_{C}$, images of $\mu_{C}$ by $x \mapsto \theta \cdot x$ for $\theta \in S^{n-1}$, have all variance 1 . We shall say in this case that $C$ is isotropic and normalized by variance. We have then in addition that

$$
\int_{C}|x|^{2} \mathrm{~d} x=n|C| \quad \text { and } \quad|C|=L(C)^{-n}
$$

If we look for a (centrally symmetric) Euclidean ball in $\mathbb{R}^{n}$ normalized by variance, its radius $r=r_{n, V}$ must therefore satisfy $\int_{0}^{r} t^{n+1} s_{n-1} \mathrm{~d} t=$ $n \int_{0}^{r} t^{n-1} s_{n-1} \mathrm{~d} t$, giving

$$
\begin{equation*}
r_{n, V}=\sqrt{n+2} \tag{5.4}
\end{equation*}
$$

In the same way, we can bring to isotropy a symmetric probability density $K$ on $\mathbb{R}^{n}$, i.e., such that $K(-x)=K(x)$ for $x \in \mathbb{R}^{n}$, by a linear change to $K_{(A)}$ for some $A$ linear and invertible. When $K$ is isotropic, there exists $\sigma>0$ such that

$$
\int_{\mathbb{R}^{n}}(\xi \cdot x)^{2} K(x) \mathrm{d} x=\sigma^{2}|\xi|^{2}, \quad \xi \in \mathbb{R}^{n}
$$

which means that all one-dimensional marginals of $K$ have the same variance $\sigma^{2}$. We shall then say for brevity that $K$ is isotropic with variance $\sigma^{2}$. The dilated density $K_{(1 / \sigma)}: x \mapsto \sigma^{n} K(\sigma x)$ is normalized by variance. For example, the standard Gaussian $\gamma_{n}$ in (1.17) is normalized by variance. For the study of maximal functions, we can always assume that $K$ is normalized by variance.

L. Deleaval, O. Guédon and B. Maurey

### 5.2. On the volume of sections

We have seen in (2.14) that the Fourier transform $m$ of a kernel $K \in$ $L^{1}\left(\mathbb{R}^{n}\right)$ can be expressed as

$$
m(u \xi)=\int_{\mathbb{R}} \varphi_{\theta, K}(s) \mathrm{e}^{-2 \mathrm{i} \pi s u|\xi|} \mathrm{d} s, \quad u \in \mathbb{R}, \xi \in \mathbb{R}^{n} \backslash\{0\}
$$

where one has set $\theta=|\xi|^{-1} \xi$ and $\varphi_{\theta}(s)=\varphi_{\theta, K}(s)=\int_{\theta \perp} K(y+s \theta) \mathrm{d}^{n-1} y$ for every $s \in \mathbb{R}$. When $K$ is the kernel $K_{C}$ corresponding to a symmetric convex body $C$, the function $\varphi_{\theta}$ is the "normalized" function of $(n-1)$-dimensional volumes of hyperplane sections parallel to $\theta^{\perp}$, defined by

$$
\varphi_{\theta, C}(s)=\int_{\theta^{\perp}} K_{C}(y+s \theta) \mathrm{d}^{n-1} y=\frac{\left|C \cap\left(\theta^{\perp}+s \theta\right)\right|_{n-1}}{|C|_{n}}
$$

We know by the Brunn-Minkowski inequality [37, Theorem 4.1] that $\varphi_{\theta, C}$ is $\log$-concave on $\mathbb{R}$. Indeed, a form of this inequality states that

$$
|(1-\lambda) A+\lambda B| \geqslant|A|^{1-\lambda}|B|^{\lambda}
$$

whenever $A, B$ are compact subsets of $\mathbb{R}^{n}$ and $\lambda \in[0,1]$. Recall that a function $K \geqslant 0$ on $\mathbb{R}^{n}$ is log-concave when

$$
K\left((1-\alpha) x_{0}+\alpha x_{1}\right) \geqslant K\left(x_{0}\right)^{1-\alpha} K\left(x_{1}\right)^{\alpha}, \quad x_{0}, x_{1} \in \mathbb{R}^{n}, \alpha \in[0,1]
$$

in other words, when $\log K$ is concave on the convex set $\{K>0\}$.
More generally than Brunn-Minkowski, the Prékopa-Leindler inequality [37, Theorem 7.1] implies that the function $\varphi_{\theta, K}$ defined in (2.14) is a log-concave probability density on the real line if $K$ is a log-concave probability density on $\mathbb{R}^{n}$. The statement of Prékopa-Leindler is as follows: if $\alpha$ is in $(0,1)$, if $f_{0}, f_{1}, f_{\alpha}$ nonnegative and integrable Borel functions on $\mathbb{R}^{n}$ are such that

$$
f_{\alpha}\left((1-\alpha) x_{0}+\alpha x_{1}\right) \geqslant f_{0}\left(x_{0}\right)^{1-\alpha} f_{1}\left(x_{1}\right)^{\alpha}
$$

for all $x_{0}, x_{1} \in \mathbb{R}^{n}$, then

$$
\int_{\mathbb{R}^{n}} f_{\alpha}(x) \mathrm{d} x \geqslant\left(\int_{\mathbb{R}^{n}} f_{0}(x) \mathrm{d} x\right)^{1-\alpha}\left(\int_{\mathbb{R}^{n}} f_{1}(x) \mathrm{d} x\right)^{\alpha}
$$

Given $\theta \in S^{n-1}, s_{0}, s_{1}$ real and letting $f_{j}(y)=K\left(y+s_{j} \theta\right)$ for $y \in \theta^{\perp}$ and $j=0,1, s_{\alpha}=(1-\alpha) s_{0}+\alpha s_{1}$ and $f_{\alpha}(y)=K\left(y+s_{\alpha} \theta\right)$, we obtain that $\varphi_{\theta, K}$ is log-concave by applying Prékopa-Leindler on $\theta^{\perp} \simeq \mathbb{R}^{n-1}$ to these functions $f_{0}, f_{1}$ and $f_{\alpha}$. Similarly, one shows that convolutions of log-concave densities are log-concave. Without more effort, Bourgain's proof also gives the following theorem.

Theorem 5.2. - There exists a constant $\kappa_{2}<140$ such that for every integer $n \geqslant 1$ and every symmetric log-concave probability density $K$ on $\mathbb{R}^{n}$, one has

$$
\forall f \in L^{2}\left(\mathbb{R}^{n}\right), \quad\left\|\mathrm{M}_{K} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)} \leqslant \kappa_{2}\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)}
$$

We turn to the proof of the main inequalities about log-concave functions, which will be used throughout our Notes. We introduce the right maximal function $f_{r}^{*}$ of a locally integrable function $f$ on an interval $[\tau,+\infty)$ of the line by setting

$$
\begin{equation*}
f_{r}^{*}(x)=\sup _{t>0} \frac{1}{t} \int_{x}^{x+t}|f(s)| \mathrm{d} s, \quad x \geqslant \tau \tag{5.5}
\end{equation*}
$$

One sees that $f_{r}^{*} \leqslant f^{*} \leqslant 2 \mathrm{M} f$, where $f^{*}$ is the uncentered maximal function from (0.2), and $f_{r}^{*}(x) \geqslant|f(x)|$ at each Lebesgue point $x$ of $f$, hence almost everywhere. When $\psi$ is nonnegative, integrable and decreasing on $[x,+\infty)$, then

$$
\begin{equation*}
\int_{x}^{+\infty}|f(s)| \psi(s) \mathrm{d} s \leqslant\left(\int_{x}^{+\infty} \psi(s) \mathrm{d} s\right) f_{r}^{*}(x) \tag{5.6}
\end{equation*}
$$

One can get (5.6) as in (4.6), by approximating $\psi$ by a combination of functions $t_{k}^{-1} \mathbf{1}_{\left[x, x+t_{k}\right]}$. We can also define in a similar way a left maximal function $f_{\ell}^{*}$.

Lemma 5.3. - Let $\varphi$ be an integrable log-concave function on an interval $[\tau,+\infty)$, let $p$ belong to $(0,+\infty)$ and let

$$
S_{0}(\tau)=\int_{\tau}^{+\infty} \varphi(s) \mathrm{d} s, \quad S_{p}(\tau)=\int_{\tau}^{+\infty}(s-\tau)^{p} \varphi(s) \mathrm{d} s
$$

Then $S_{p}(\tau)$ is finite. Furthermore, assuming $S_{p}(\tau)>0$, we have

$$
\begin{equation*}
\varphi(\tau)^{p} \leqslant \frac{\Gamma(p+1) S_{0}(\tau)^{p+1}}{S_{p}(\tau)}, \quad \max _{s \geqslant \tau} \varphi(s)^{p} \geqslant \varphi_{r}^{*}(\tau)^{p} \geqslant \frac{S_{0}(\tau)^{p+1}}{(p+1) S_{p}(\tau)} \tag{5.7}
\end{equation*}
$$

Proof. - We have $\varphi \geqslant 0$ by definition of log-concavity. We assume $S_{p}(\tau)>0$, hence $S_{0}(\tau)>0$. We may suppose $\tau=0$ by translating and $S_{0}:=S_{0}(0)=1$ by homogeneity. We begin with the left-hand inequality in (5.7), assuming $a:=\varphi(0)>0$. Consider the log-affine probability density $\psi(s)=a \mathrm{e}^{-a s}$ on $[0,+\infty)$, chosen so that $\psi(0)=\varphi(0)$. By log-concavity, the set $I=\{\varphi \geqslant \psi\}$ is an interval, such that $0 \in I \subset[0,+\infty)$. Since $\varphi$ and $\psi$ both have integral 1 on $[0,+\infty)$, the interval $I$ is not reduced to $\{0\}$. If $I=[0,+\infty)$, the densities are equal and

$$
\begin{aligned}
S_{p}:=\int_{0}^{+\infty} s^{p} \varphi(s) \mathrm{d} s=\int_{0}^{+\infty} s^{p} \psi(s) \mathrm{d} s & =\frac{1}{a^{p}} \int_{0}^{+\infty}(a s)^{p} \mathrm{e}^{-a s} a \mathrm{~d} s \\
& =\frac{\Gamma(p+1)}{a^{p}}
\end{aligned}
$$

Otherwise, the interval $I$ is bounded, let $s_{0}:=\sup I>0$. We have $\psi \leqslant \varphi$ on $\left[0, s_{0}\right)$ and $\varphi(s)<\psi(s)$ when $s>s_{0}$, implying that $S_{p}(0)$ is finite. The antiderivative $F$ of $\varphi-\psi$ vanishing at 0 is first increasing, then decreasing on $[0,+\infty)$, and tends to 0 at infinity because $\varphi$ and $\psi$ have equal integrals. It follows that $F$ is nonnegative on $[0,+\infty)$. Recalling that $0 \leqslant \varphi(s)<\psi(s)$ at infinity, we know that $|F(s)|$ is exponentially small at infinity, and integrating by parts we obtain

$$
\int_{0}^{+\infty} s^{p}(\varphi(s)-\psi(s)) \mathrm{d} s=-p \int_{0}^{+\infty} s^{p-1} F(s) \mathrm{d} s \leqslant 0
$$

One concludes the first part by writing

$$
S_{p}=\int_{0}^{+\infty} s^{p} \varphi(s) \mathrm{d} s \leqslant \int_{0}^{+\infty} s^{p} \psi(s) \mathrm{d} s=\frac{\Gamma(p+1)}{a^{p}}
$$

For the right-hand inequality in (5.7), we let $b=\varphi_{r}^{*}(0)>0$ and consider the probability density $\psi(s)=b \mathbf{1}_{[0,1 / b]}(s)$ on $[0,+\infty)$. Let $F$ be the antiderivative of $\varphi-\psi$ vanishing at 0 . When $0<x \leqslant 1 / b$ we have by definition of $\varphi_{r}^{*}(0)$ that

$$
\frac{F(x)}{x}=\frac{1}{x} \int_{0}^{x}(\varphi(s)-\psi(s)) \mathrm{d} s=\left(\frac{1}{x} \int_{0}^{x} \varphi(s) \mathrm{d} s\right)-b \leqslant 0
$$

We see that $\psi(x)=0 \leqslant \varphi(x)$ when $x \geqslant 1 / b$. It follows that the function $F$ is $\leqslant 0$ on $[0,1 / b]$, then increasing on $[1 / b,+\infty)$, tends to 0 at infinity, thus $F$ is $\leqslant 0$ on the half-line $[0,+\infty)$. Arguing as before, we have consequently

$$
S_{p}=\int_{0}^{+\infty} s^{p} \varphi(s) \mathrm{d} s \geqslant b \int_{0}^{1 / b} s^{p} \mathrm{~d} s=\frac{1}{(p+1) b^{p}}
$$

For every $\theta \in S^{n-1}$, the function $\varphi_{\theta, C}$ associated to a symmetric convex set $C$ is even, log-concave and has integral 1 by definition. We shall thus be in a position to apply to it the following Corollary 5.4.

Corollary 5.4. - Suppose that $\varphi$ is a symmetric log-concave probability density on $\mathbb{R}$ and let $\sigma^{2}:=\int_{\mathbb{R}} s^{2} \varphi(s) \mathrm{d} s$. One has that

$$
\frac{1}{12 \sigma^{2}} \leqslant \varphi(0)^{2}=\max _{s \in \mathbb{R}} \varphi(s)^{2} \leqslant \frac{1}{2 \sigma^{2}}
$$

Proof. - Since $\varphi$ is even and log-concave, we have $\varphi(0)=\max _{s \in \mathbb{R}} \varphi(s)$. We apply Lemma 5.3 with $p=2, \tau=0$, and observe that $S_{0}(0)=1 / 2$, $S_{2}(0)=\sigma^{2} / 2$.

The preceding result is sharp, as one sees with the two examples

$$
\begin{equation*}
\varphi_{0}(s)=\frac{1}{\sqrt{2}} \mathrm{e}^{-\sqrt{2}|s|}, \quad \varphi_{1}(s)=\frac{1}{2 \sqrt{3}} \mathbf{1}_{[-\sqrt{3}, \sqrt{3}]}(s), \quad s \in \mathbb{R} \tag{5.8}
\end{equation*}
$$

The next corollary is not very sharp, but easy to deduce from Lemma 5.3. When the function $\varphi \geqslant 0$ is defined on the line and $p \in[0,+\infty)$, we set

$$
S_{p}^{+}(\tau)=\int_{\tau}^{+\infty}(s-\tau)^{p} \varphi(s) \mathrm{d} s, \quad S_{p}^{-}(\tau)=\int_{-\infty}^{\tau}|s-\tau|^{p} \varphi(s) \mathrm{d} s
$$

Corollary 5.5. - Let $\varphi$ be a centered log-concave probability density on $\mathbb{R}$ and let $\sigma^{2}:=\int_{\mathbb{R}} s^{2} \varphi(s) \mathrm{d} s$. We have that

$$
\frac{1}{24 \sigma^{2}} \leqslant \frac{\varphi_{\ell}^{*}(0)^{2}+\varphi_{r}^{*}(0)^{2}}{2} \leqslant \max _{s \in \mathbb{R}} \varphi(s)^{2} \leqslant \frac{4}{\sigma^{2}}
$$

Proof. - We begin with the rightmost inequality. Let us fix $\tau$ real. Since $\varphi$ is a centered probability density, one has that

$$
S_{2}^{+}(\tau)+S_{2}^{-}(\tau)=\int_{\mathbb{R}}(s-\tau)^{2} \varphi(s) \mathrm{d} s=\sigma^{2}+\tau^{2} \geqslant \sigma^{2}, \quad \tau \in \mathbb{R}
$$

Up to a symmetry around $\tau$, possibly replacing the function $\varphi$ by $s \mapsto$ $\varphi(2 \tau-s)$, we may assume that $S_{2}^{+}(\tau) \geqslant \sigma^{2} / 2$. We have $S_{0}^{+}(\tau)=$ $\int_{\tau}^{+\infty} \varphi(s) \mathrm{d} s \leqslant 1$ since $\varphi$ is a probability density on $\mathbb{R}$, thus by Lemma 5.3 with $p=2$ we get

$$
\varphi(\tau)^{2} \leqslant \frac{2 S_{0}^{+}(\tau)^{3}}{S_{2}^{+}(\tau)} \leqslant \frac{4}{\sigma^{2}}
$$

Since $\tau$ is arbitrary, we obtain the right-hand inequality. Let us pass to the other inequality. By Lemma 5.3 with $p=2$ on the intervals $(0,+\infty)$ and $(-\infty, 0)$, we conclude using $S_{2}^{ \pm}(0) \leqslant \sigma^{2}$ and $S_{0}^{+}(0)+S_{0}^{-}(0)=1$ that

$$
\varphi_{r}^{*}(0)^{2}+\varphi_{\ell}^{*}(0)^{2} \geqslant \frac{S_{0}^{+}(0)^{3}}{3 S_{2}^{+}(0)}+\frac{S_{0}^{-}(0)^{3}}{3 S_{2}^{-}(0)} \geqslant \frac{S_{0}^{+}(0)^{3}+S_{0}^{-}(0)^{3}}{3 \sigma^{2}} \geqslant \frac{1}{12 \sigma^{2}} .
$$

Lemma 5.6. - Let $\varphi$ be a symmetric log-concave probability density on $\mathbb{R}$, with variance $\sigma^{2}$. The function $\varphi$ decays exponentially at infinity, with a rate depending on its variance and satisfying

$$
\forall s \in \mathbb{R}, \quad \sigma \varphi(\sigma s) \leqslant \min \left(2 \mathrm{e}^{-|s| / 2}, 11 \mathrm{e}^{-|s|}\right)
$$

Proof. - Without loss of generality, we may assume that $\sigma=1$. It follows then from Corollary 5.4 that $1 /(2 \sqrt{3}) \leqslant a:=\varphi(0) \leqslant 1 / \sqrt{2}$. Consider the log-affine function $\psi_{\beta}(s)=a \mathrm{e}^{-\beta s}$ on $[0,+\infty)$, with $\beta>0$, satisfying $\psi_{\beta}(0)=$ $\varphi(0)$. If we have $\varphi\left(\tau_{0}\right) \leqslant \psi_{\beta}\left(\tau_{0}\right)$ for some $\tau_{0}>0$, it implies by log-concavity that $\varphi(s) \leqslant \psi_{\beta}(s) \leqslant \mathrm{e}^{-\beta s} / \sqrt{2}$ for $s \geqslant \tau_{0}$, and in order to obtain a bound for $\varphi$ everywhere, we can apply for the values $0 \leqslant s \leqslant \tau_{0}$ the obvious inequalities

$$
\varphi(s) \leqslant \varphi(0)=a \leqslant a \mathrm{e}^{\beta\left(\tau_{0}-s\right)} \leqslant\left(\mathrm{e}^{\beta \tau_{0}} / \sqrt{2}\right) \mathrm{e}^{-\beta s} .
$$

For any $\tau_{0}>0$, we obtain since $\varphi$ is even that

$$
\begin{equation*}
\varphi\left(\tau_{0}\right) \leqslant \psi_{\beta}\left(\tau_{0}\right) \Rightarrow \varphi(s)=\varphi(|s|) \leqslant \mathrm{e}^{\beta \tau_{0}-\ln \sqrt{2}} \mathrm{e}^{-\beta|s|}, \quad s \in \mathbb{R} \tag{5.9}
\end{equation*}
$$

On the other hand, if $\varphi(s)>\psi_{\beta}(s)$ for every $s \in(0, \tau]$, then

$$
\begin{aligned}
1 / 2 & =\int_{0}^{+\infty} s^{2} \varphi(s) \mathrm{d} s>\int_{0}^{\tau} s^{2} \psi_{\beta}(s) \mathrm{d} s=\frac{a}{\beta^{3}} \int_{0}^{\beta \tau} u^{2} \mathrm{e}^{-u} \mathrm{~d} u \\
& =\frac{a}{\beta^{3}}\left[-\mathrm{e}^{-u}\left(u^{2}+2 u+2\right)\right]_{u=0}^{\beta \tau} \geqslant \frac{1}{2 \sqrt{3} \beta^{3}}\left(2-\mathrm{e}^{-\beta \tau}\left(\beta^{2} \tau^{2}+2 \beta \tau+2\right)\right)
\end{aligned}
$$

Equivalently, when $\varphi(s)>\psi_{\beta}(s)$ for every $s \in(0, \tau]$, we get that

$$
\begin{equation*}
\mathrm{e}^{-\beta \tau}\left(\beta^{2} \tau^{2}+2 \beta \tau+2\right)>2-\sqrt{3} \beta^{3} \tag{5.10}
\end{equation*}
$$

Suppose that $\beta^{3}<2 / \sqrt{3}$. Then (5.10) cannot be true if $\tau$ is large. For every such $\beta$, there exists $\tau_{0}>0$ such that $\varphi\left(\tau_{0}\right) \leqslant \psi_{\beta}\left(\tau_{0}\right)$ and by (5.9), there is a constant $c(\beta)$ such that $\varphi(s) \leqslant c(\beta) \mathrm{e}^{-\beta|s|}$ on the line. For numerical purposes, it is more convenient to express this as follows. If $0<\sqrt{3} \beta^{3}<2$ and if

$$
\begin{equation*}
\mathrm{e}^{-x}\left(x^{2}+2 x+2\right) \leqslant 2-\sqrt{3} \beta^{3} \tag{5.11}
\end{equation*}
$$

then $x>0$, and letting $x_{0}(\beta)=x$, we know that $\varphi(s) \leqslant \psi_{\beta}(|s|) \leqslant \mathrm{e}^{-\beta|s|} / \sqrt{2}$ when $|s| \geqslant \tau_{0}(\beta):=x_{0}(\beta) / \beta$, and $\varphi(s) \leqslant c(\beta) \mathrm{e}^{-\beta|s|}$ for every $s \in \mathbb{R}$ by (5.9), with

$$
\begin{equation*}
c(\beta)=\mathrm{e}^{\beta \tau_{0}(\beta)-\ln \sqrt{2}}=\mathrm{e}^{x_{0}(\beta)-\ln \sqrt{2}} \tag{5.12}
\end{equation*}
$$

An almost optimal $x$ satisfying (5.11) can be found numerically. We have for example that $\varphi(s) \leqslant 2.218 \mathrm{e}^{-|s| / 2}$ for all $s$ when $\beta=1 / 2$, with $x_{0}(0.5)=$ 1.143. We also find $c(1)<94.295$ with a choice $x_{0}(1)=4.893$. We can then improve the first estimate given by (5.12) for $\beta=1$. When $|s| \leqslant x_{0}(1)=$ $\tau_{0}(1)$, we write

$$
\varphi(s) \leqslant 2.218 \mathrm{e}^{-|s| / 2}=2.218 \mathrm{e}^{|s| / 2} \mathrm{e}^{-|s|} \leqslant 2.218 \mathrm{e}^{\tau_{0}(1) / 2} \mathrm{e}^{-|s|}<26 \mathrm{e}^{-|s|}
$$

More generally, if we know a modified bound $c_{m}\left(\beta_{1}\right)$ such that $\varphi(s) \leqslant$ $c_{m}\left(\beta_{1}\right) \mathrm{e}^{-\beta_{1}|s|}$ for every $s$ and if $\varphi(s) \leqslant \mathrm{e}^{-\beta_{2}|s|} / \sqrt{2}$ when $|s| \geqslant \tau_{0}\left(\beta_{2}\right)$, with $\beta_{1}<\beta_{2}$, then for $|s| \leqslant \tau_{0}\left(\beta_{2}\right)$ we can write

$$
\begin{aligned}
\varphi(s) & \leqslant c_{m}\left(\beta_{1}\right) \mathrm{e}^{-\beta_{1}|s|}=c_{m}\left(\beta_{1}\right) \mathrm{e}^{\left(\beta_{2}-\beta_{1}\right)|s|} \mathrm{e}^{-\beta_{2}|s|} \\
& \leqslant c_{m}\left(\beta_{1}\right) \mathrm{e}^{\left(\beta_{2}-\beta_{1}\right) \tau_{0}\left(\beta_{2}\right)} \mathrm{e}^{-\beta_{2}|s|}
\end{aligned}
$$

so that

$$
\begin{equation*}
c_{m}\left(\beta_{2}\right) \leqslant \max \left(\mathrm{e}^{\left(\beta_{2}-\beta_{1}\right) \tau_{0}\left(\beta_{2}\right)} c_{m}\left(\beta_{1}\right), 1 / \sqrt{2}\right) \tag{5.13}
\end{equation*}
$$

The following table displays admissible values for $x_{0}(\beta), \tau_{0}(\beta)$, then the corresponding rough bound $c(\beta)$ from (5.12), and the modified bounds $c_{m}(\beta)$ obtained step by step applying (5.13), by dividing the interval $[0,1]$ in ten equal segments, beginning with $c(0)=c_{m}(0)=\varphi(0) \leqslant 1 / \sqrt{2}<0.708$. We have replaced each higher precision value of $x$ by the upper bound
$x_{0}(\beta)=\lceil 1000 . x\rceil / 1000$, and used this replacement consistently in the further calculations of $\tau_{0}(\beta), c(\beta)$ and $c_{m}(\beta)$.

| $\beta$ | $x_{0}(\beta)$ | $\tau_{0}(\beta)$ | $c(\beta)$ | $c_{m}(\beta)$ |
| :---: | :---: | :---: | :---: | ---: |
| 0.0 | 0.000 | 0.000 | 0.708 | 0.708 |
| 0.1 | 0.182 | 1.820 | 0.849 | 0.850 |
| 0.2 | 0.381 | 1.906 | 1.036 | 1.029 |
| 0.3 | 0.603 | 2.010 | 1.293 | 1.259 |
| 0.4 | 0.854 | 2.135 | 1.662 | 1.559 |
| 0.5 | 1.143 | 2.287 | 2.218 | 1.960 |
| 0.6 | 1.484 | 2.474 | 3.119 | 2.511 |
| 0.7 | 1.903 | 2.719 | 4.742 | 3.296 |
| 0.8 | 2.451 | 3.064 | 8.203 | 4.478 |
| 0.9 | 3.255 | 3.617 | 18.328 | 6.430 |
| 1.0 | 4.893 | 4.893 | 94.295 | 10.489 |

We obtain the announced bounds when $\beta=1 / 2$ and $\beta=1$. One can obviously refine the previous argument and show that

$$
\varphi(s) \leqslant c(0) \exp \left(\int_{0}^{1} \tau_{0}(\beta) \mathrm{d} \beta\right) \mathrm{e}^{-|s|} \leqslant \frac{1}{\sqrt{2}} \exp \left(\int_{0}^{1} \tau_{0}(\beta) \mathrm{d} \beta\right) \mathrm{e}^{-|s|}
$$

We may get in this way that $\varphi(s)<9 \mathrm{e}^{-|s|}$. An exact estimate could perhaps be obtained by an extreme point argument, as in [35]. Some numerical experiments suggest that for every $\beta \geqslant 0$, the maximum on $\mathbb{R}$ of $s \mapsto \mathrm{e}^{\beta|s|} \varphi(s)$, for $\varphi$ symmetric log-concave probability density with variance 1 , occurs for one of the two examples $\varphi_{0}, \varphi_{1}$ mentioned in (5.8). The example $\varphi_{0}(s)$ shows that $\mathrm{e}^{\beta|s|} \varphi(s)$ is unbounded when $\beta>\sqrt{2}$ and $\sigma=1$.

Our next estimate is so poor that it does not deserve to be given explicitly.
Corollary 5.7. - There exists a numerical value $\kappa>0$ such that for every centered log-concave probability density $\varphi$ on $\mathbb{R}$ with variance $\sigma^{2}=1$, one has

$$
\forall s \in \mathbb{R}, \quad \varphi(s) \leqslant \kappa \mathrm{e}^{-|s| / \kappa} .
$$

Proof. - Since $\varphi$ is centered, we know that $\int_{0}^{+\infty} s \varphi(s) \mathrm{d} s=$ $\int_{-\infty}^{0}|s| \varphi(s) \mathrm{d} s$, and we can thus set $S_{1}:=S_{1}^{+}(0)=S_{1}^{-}(0)$. For $p \neq 1$, let us write $S_{p}^{ \pm}$instead of $S_{p}^{ \pm}(0)$. We have that $S_{2}^{+}, S_{2}^{-} \leqslant \sigma^{2}=1$. By Corollary 5.5 and Lemma 5.3 with $p=1$, applied on the intervals $[0,+\infty)$ and $(-\infty, 0]$, we get

$$
2 \geqslant \max _{s \geqslant 0} \varphi(s) \geqslant \frac{\left(S_{0}^{+}\right)^{2}}{2 S_{1}}, \quad 2 \geqslant \max _{s \leqslant 0} \varphi(s) \geqslant \frac{\left(S_{0}^{-}\right)^{2}}{2 S_{1}} .
$$

It follows that $8 S_{1} \geqslant\left(S_{0}^{+}\right)^{2}+\left(S_{0}^{-}\right)^{2} \geqslant 1 / 2$ so $S_{1} \geqslant 1 / 16$. We also need a lower bound for $S_{0}^{ \pm}$. Let $\kappa_{1}=16$. By Cauchy-Schwarz we have

$$
\kappa_{1}^{-2} \leqslant S_{1}^{2} \leqslant S_{0}^{-} S_{2}^{-} \leqslant S_{0}^{-}, \quad \kappa_{1}^{-2} \leqslant S_{1}^{2} \leqslant S_{0}^{+} S_{2}^{+} \leqslant S_{0}^{+}
$$

hence $S_{0}^{-}, S_{0}^{+} \geqslant \kappa_{1}^{-2}$. Suppose that the maximum of $\varphi$ is reached at $s_{0} \geqslant 0$. Then $\varphi$ is non-decreasing on $\left(-\infty, s_{0}\right]$ and by Lemma 5.3 with $p=2$ we get

$$
\begin{equation*}
4 \geqslant \varphi(0)^{2}=\max _{s \leqslant 0} \varphi(s)^{2} \geqslant \frac{\left(S_{0}^{-}\right)^{3}}{3 S_{2}^{-}} \geqslant \frac{\kappa_{1}^{-6}}{3}=: \kappa_{2}^{-2} \tag{5.14}
\end{equation*}
$$

The symmetric probability density $\varphi_{1}(s)=\left(2 S_{0}^{-}\right)^{-1} \varphi(-|s|)$ on $\mathbb{R}$ is logconcave, has variance $\sigma_{1}^{2}=S_{2}^{-} / S_{0}^{-} \leqslant \kappa_{1}^{2}$. By (5.14), we have $\left(S_{0}^{-}\right)^{3} / S_{2}^{-} \leqslant 12$. By Lemma 5.6, we know that

$$
\varphi_{1}(s) \leqslant \frac{11}{\sigma_{1}} \mathrm{e}^{-|s| / \sigma_{1}}, \quad \text { and } \quad \varphi(s) \leqslant 22\left(\frac{\left(S_{0}^{-}\right)^{3}}{S_{2}^{-}}\right)^{1 / 2} \mathrm{e}^{-|s| / \sigma_{1}} \leqslant 77 \mathrm{e}^{-|s| / \kappa_{1}}
$$

for $s \leqslant 0$. Let us pass to the positive side. We let $\widetilde{\varphi}$ be equal to $\varphi\left(s_{0}\right)$ on $\left[0, s_{0}\right]$ and to $\varphi$ on $\left[s_{0},+\infty\right)$. Then $\widetilde{S}_{0}^{+} \geqslant S_{0}^{+} \geqslant \kappa_{1}^{-2}$ and since $\kappa_{2}^{-1} \leqslant \varphi(0) \leqslant$ $\varphi(x) \leqslant \varphi\left(s_{0}\right) \leqslant 2$ when $0 \leqslant x \leqslant s_{0}$, we have $\widetilde{\varphi} \leqslant 2 \kappa_{2} \varphi$ on $[0,+\infty)$. The symmetrized function $\varphi_{1}$ corresponding to $\widetilde{\varphi}$ satisfies $\sigma_{1}^{2}=\widetilde{S}_{2}^{+} / \widetilde{S}_{0}^{+} \leqslant 2 \kappa_{1}^{2} \kappa_{2}$. Also, we know that $\left(\widetilde{S}_{0}^{+}\right)^{3} / \widetilde{S}_{2}^{+} \leqslant 3 \max \widetilde{\varphi}(s)^{2} \leqslant 12$. The rest is identical to the negative case.

The next lemma is easy and classical. The (total) mass of a real valued (thus bounded) measure $\mu$ on $(\Omega, \mathcal{F})$ is defined by setting $\|\mu\|_{1}=\mu^{+}(\Omega)+$ $\mu^{-}(\Omega)=|\mu|(\Omega)$, where $\mu=\mu^{+}-\mu^{-}$is the Hahn decomposition of $\mu$ as difference of two nonnegative measures, and $|\mu|=\mu^{+}+\mu^{-}$. On the line or on $\mathbb{R}^{n}$ we have

$$
\|\mu\|_{1}=\sup \left\{\left|\int_{\mathbb{R}^{n}} \psi \mathrm{~d} \mu\right|: \psi \in \mathcal{K}\left(\mathbb{R}^{n}\right),\|\psi\|_{\infty} \leqslant 1\right\}
$$

and when $\mu$ has a density $f$, one has that $\|f(x) \mathrm{d} x\|_{1}=\|f\|_{L^{1}\left(\mathbb{R}^{n}\right)}$.
Lemma 5.8. - Let $\mu$ be a real valued measure on $\mathbb{R}$ and let $m(t)=\widehat{\mu}(t)$ be its Fourier transform. For every $t \in \mathbb{R}$ we have

$$
\begin{equation*}
|m(t)| \leqslant\|\mu\|_{1} . \tag{5.15a}
\end{equation*}
$$

If $\mathrm{d} \mu(s)=\psi(s) \mathrm{d} s$ with $\psi$ integrable, then $m=\widehat{\mu}=\widehat{\psi}$ and $|m(t)| \leqslant$ $\|\psi\|_{L^{1}\left(\mathbb{R}^{n}\right)}$.

Let us further assume that $\int_{\mathbb{R}}(1+|s|) \mathrm{d}|\mu|(s)<+\infty$. Then $m$ is $C^{1}$ on $\mathbb{R}$ and

$$
\mathrm{i} m^{\prime}(t)=2 \pi \int_{\mathbb{R}} s \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} \mu(s)
$$

so $\mathrm{i} m^{\prime}$ is the Fourier transform of the real valued measure $2 \pi s \mathrm{~d} \mu(s)$.

Let $\nu$ be a real valued measure on $\mathbb{R}$ and let $\psi(s)=\nu((-\infty, s])$, for every $s \in \mathbb{R}$. The measure $\nu$ is the derivative of $\psi$ in the sense of distributions and assuming $\psi$ integrable, we have

$$
\begin{equation*}
2 \mathrm{i} \pi t \widehat{\psi}(t)=2 \mathrm{i} \pi t \int_{\mathbb{R}} \psi(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s=\int_{\mathbb{R}} \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} \nu(s) \tag{5.15b}
\end{equation*}
$$

so $2 \mathrm{i} \pi t \widehat{\psi}(t)$ is the Fourier transform of the derivative $\nu$ of $\psi$.
Let $j, k$ be nonnegative integers. Suppose that $\psi$ is of class $C^{k-1}$ on the line, with a $k$ th derivative $\psi^{(k)}$ in the sense of distributions that is a bounded measure $\nu_{k}$ on $\mathbb{R}$, and that $\lim _{|s| \rightarrow+\infty} \psi(s)=0, \int_{\mathbb{R}}|s|^{j} \mathrm{~d}\left|\nu_{k}\right|(s)<+\infty$. Then $m$ is $C^{j}$ and

$$
\begin{equation*}
(2 \pi|t|)^{k}\left|m^{(j)}(t)\right| \leqslant(2 \pi)^{j}\left\|\left(s^{j} \psi(s)\right)^{(k)}\right\|_{1} \tag{5.15c}
\end{equation*}
$$

Consequently, for $t \neq 0$, we have that

$$
\begin{array}{r}
\left|m^{(j)}(t)\right| \leqslant \frac{(2 \pi)^{j-k}}{|t|^{k}} \sum_{i=(k-j)^{+}}^{k-1}\binom{k}{i} \frac{j!}{(j+i-k)!} \int_{\mathbb{R}}|s|^{i+j-k}\left|\psi^{(i)}(s)\right| \mathrm{d} s \\
+\frac{(2 \pi)^{j-k}}{|t|^{k}} \int_{\mathbb{R}}|s|^{j} \mathrm{~d}\left|\nu_{k}\right|(s) \tag{5.15d}
\end{array}
$$

In the line above, one can replace $\int_{\mathbb{R}}|s|^{j} \mathrm{~d}\left|\nu_{k}\right|(s)$ with $\int_{\mathbb{R}}|s|^{j}\left|\psi^{(k)}(s)\right| \mathrm{d} s$, when $\psi$ admits a derivative $\psi^{(k)}$ and $\mathrm{d} \nu_{k}(x)=\psi^{(k)}(x) \mathrm{d} x$.

Proof. - The first inequality (5.15a) is obvious. Assuming that $\int_{\mathbb{R}}|s| \mathrm{d}|\mu|(s)$ is finite, we write

$$
m(t)=\int_{\mathbb{R}} \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} \mu(s):=\int_{\mathbb{R}} \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} \mu^{+}(s)-\int_{\mathbb{R}} \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} \mu^{-}(s)
$$

and we obtain by the dominated convergence theorem that

$$
m^{\prime}(t)=-2 \mathrm{i} \pi \int_{\mathbb{R}} s \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} \mu(s) .
$$

If $\nu$ in (5.15b) has the form $\mathrm{d} \nu(x)=\psi^{\prime}(x) \mathrm{d} x$ with $\psi^{\prime}$ a true derivative, we use integration by parts, otherwise we use Fubini's theorem for $\nu^{+}$and $\nu^{-}$. We get

$$
2 \mathrm{i} \pi t \int_{\mathbb{R}} \psi(t) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s=\int_{\mathbb{R}} \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} \nu(s) .
$$

The verification of $(5.15 \mathrm{~d})$ is left to the reader. Notice that by (5.16), the hypotheses imply that $\int_{\mathbb{R}}|s|^{i+j-k}\left|\psi^{(i)}(s)\right| \mathrm{d} s<+\infty$ when $(k-j)^{+} \leqslant i<k$. Indeed, if $g^{(\ell+1)}$ is integrable on $[0,+\infty)$, then $g^{(\ell)}$ tends to a limit $L$ at infinity and if $g$ tends to 0 at infinity, it follows that $L=0$, for example by the Taylor formula.

The next lemma is straightforward.

Lemma 5.9. - Let $\nu$ be a nonnegative measure on $(0,+\infty)$ and $\alpha>0$. One has

$$
\alpha \int_{0}^{+\infty} s^{\alpha-1} \nu([s,+\infty)) \mathrm{d} s=\int_{0}^{+\infty} s^{\alpha} \mathrm{d} \nu(s)
$$

Let $F$ be a function on $(0,+\infty)$ such that $|F(s)| \leqslant \int_{s}^{+\infty} \mathrm{d} \nu(s)$ for $s>0$. One has

$$
\alpha \int_{0}^{+\infty} s^{\alpha-1}|F(s)| \mathrm{d} s \leqslant \int_{0}^{+\infty} s^{\alpha} \mathrm{d} \nu(s)
$$

Suppose that the function $g$ is differentiable on $\mathbb{R}$, with $\lim _{s \rightarrow \pm \infty} g(s)=0$ and $g^{\prime}$ integrable on the line. It follows that

$$
\begin{equation*}
\alpha \int_{\mathbb{R}}|s|^{\alpha-1}|g(s)| \mathrm{d} s \leqslant \int_{\mathbb{R}}|s|^{\alpha}\left|g^{\prime}(s)\right| \mathrm{d} s \tag{5.16}
\end{equation*}
$$

If in addition $g$ is even and non-increasing on $[0,+\infty)$, one has

$$
\int_{\mathbb{R}}|s|^{\alpha}\left|g^{\prime}(s)\right| \mathrm{d} s=\alpha \int_{\mathbb{R}}|s|^{\alpha-1} g(s) \mathrm{d} s, \quad \text { and } \int_{\mathbb{R}}\left|g^{\prime}(s)\right| \mathrm{d} s=2 g(0)
$$

Proof. - The first assertion is an immediate consequence of Fubini, because
$\alpha \int_{0}^{+\infty} s^{\alpha-1} \nu([s,+\infty)) \mathrm{d} s=\alpha \iint \mathbf{1}_{\{0<s<t\}} s^{\alpha-1} \mathrm{~d} \nu(t) \mathrm{d} s=\int_{0}^{+\infty} t^{\alpha} \mathrm{d} \nu(t)$,
with integrals finite or not. The remaining facts are left to the reader. For (5.16), use $\mathrm{d} \nu(s)=\left|g^{\prime}(s)\right| \mathrm{d} s$.

We arrive to the main result of this section.
Proposition $5.10([9, \S 4])$. - Let $K_{l c}$ be a symmetric log-concave probability density on $\mathbb{R}^{n}$, isotropic with variance $\sigma^{2}$. Let $m_{l c}$ be the Fourier transform of $K_{l c}$. For every $\xi \in \mathbb{R}^{n}$ one has that

$$
\begin{equation*}
\pi \sqrt{2} \sigma|\xi|\left|m_{l c}(\xi)\right| \leqslant 1, \quad\left|1-m_{l c}(\xi)\right| \leqslant 2 \pi \sigma|\xi|, \quad\left|\xi \cdot \nabla m_{l c}(\xi)\right| \leqslant 2 \tag{5.17.B}
\end{equation*}
$$

The middle inequality follows from the fact that for every $\theta \in S^{n-1}$, one has

$$
\left|\theta \cdot \nabla m_{l c}(t \theta)\right| \leqslant 2 \pi \sigma, \quad t \in \mathbb{R}
$$

Remark. - These inequalities are valid for $m_{C}$, when $C$ is a symmetric convex body, isotropic and normalized by variance. The case of convex bodies is the one given by Bourgain, but the proof is the same in the log-concave case.

Proof. - We have seen in (2.14) that for $\theta \in S^{n-1}$ and $t$ real, one can write

$$
m_{l c}(t \theta)=\int_{\mathbb{R}} \varphi_{\theta}(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s
$$

where $\varphi_{\theta}$ is obtained by integrating $K_{l c}$ on affine hyperplanes parallel to $\theta^{\perp}$. It is enough to prove the case $\sigma=1$. We know that $\varphi_{\theta}$ is log-concave according to Prékopa-Leindler, it is even, has integral 1 and variance 1 by hypothesis. By Lemma 5.6 , one has that $\varphi_{\theta}(s) \leqslant 2 \mathrm{e}^{-|s| / 2}$ for every $s \in \mathbb{R}$, but the desired estimates do not depend on this exponential decay, which ensures however absolute convergence for the integrals that follow. For every $t$, by ( 5.15 d ) with $j=0, k=1$ and since $\varphi_{\theta}$ is even and decreasing on $(0,+\infty)$, we have using Lemma 5.9 that

$$
\left|m_{l c}(t \theta)\right|=\left|\int_{\mathbb{R}} \varphi_{\theta}(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s\right| \leqslant \frac{1}{2 \pi|t|} \int_{\mathbb{R}}\left|\varphi_{\theta}^{\prime}(s)\right| \mathrm{d} s=\frac{\varphi_{\theta}(0)}{\pi|t|} .
$$

The function $\varphi_{\theta}$ has variance 1 by our normalization assumption, and according to Corollary 5.4 we have the upper bound $\varphi_{\theta}(0) \leqslant 1 / \sqrt{2}$. Writing $\xi=|\xi| \theta$, it follows that $\pi \sqrt{2}|\xi|\left|m_{l c}(\xi)\right| \leqslant 1$ for every $\xi$ in $\mathbb{R}^{n}$.

Notice that our writing is not correct, because $\varphi_{\theta}$ might be discontinuous at the ends of its support, so that $\varphi_{\theta}^{\prime}$ is a measure in that case, with two Dirac masses at the end points of the support. This happens for example with $\varphi_{\theta, C}$ when $C$ is polyhedral and $\theta$ orthogonal to a facet. We leave the easy changes to the reader.

Given $\theta \in S^{n-1}$, the derivative of $t \mapsto m_{l c}(t \theta)$ is expressed by

$$
\theta \cdot \nabla m_{l c}(t \theta)=\int_{\mathbb{R}}(-2 \mathrm{i} \pi s) \varphi_{\theta}(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s
$$

and

$$
\left|\theta \cdot \nabla m_{l c}(t \theta)\right| \leqslant 2 \pi \int_{\mathbb{R}}|s| \varphi_{\theta}(s) \mathrm{d} s \leqslant 2 \pi\left(\int_{\mathbb{R}} s^{2} \varphi_{\theta}(s) \mathrm{d} s\right)^{1 / 2}=2 \pi
$$

hence $\left|1-m_{l c}(\xi)\right|=\left|m_{l c}(0)-m_{l c}(|\xi| \theta)\right| \leqslant 2 \pi|\xi|$. We see also that

$$
t \theta \cdot \nabla m_{l c}(t \theta)=\int_{\mathbb{R}}(-2 \mathrm{i} \pi t) s \varphi_{\theta}(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s=-\int_{\mathbb{R}}\left(s \varphi_{\theta}(s)\right)^{\prime} \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s
$$

We estimate the two parts coming from $\left(s \varphi_{\theta}(s)\right)^{\prime}$, first

$$
\left|\int_{\mathbb{R}} \varphi_{\theta}(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s\right| \leqslant \int_{\mathbb{R}} \varphi_{\theta}(s) \mathrm{d} s=1,
$$

and as $\varphi_{\theta}$ is even and non-increasing on $[0,+\infty)$, we have by Lemma 5.9 that

$$
\left|\int_{\mathbb{R}} s \varphi_{\theta}^{\prime}(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s\right| \leqslant \int_{\mathbb{R}}\left|s \varphi_{\theta}^{\prime}(s)\right| \mathrm{d} s=\int_{\mathbb{R}} \varphi_{\theta}(s) \mathrm{d} s=1
$$

We conclude that $\left|t \theta \cdot \nabla m_{l c}(t \theta)\right| \leqslant 2$ and get $\left|\xi \cdot \nabla m_{l c}(\xi)\right| \leqslant 2$ for every $\xi$.

Lemma 5.11. - Let $K_{l c}$ be an even log-concave probability density on $\mathbb{R}^{n}$, normalized by variance, and $m_{l c}$ its Fourier transform. For every $\theta \in$ $S^{n-1}$ one has

$$
\left|\frac{\mathrm{d}^{j}}{\mathrm{~d} t^{j}} m_{l c}(t \theta)\right| \leqslant \delta_{j, c} \frac{1}{1+2 \pi|t|}, \quad j \geqslant 0, t \in \mathbb{R}
$$

where $\delta_{j, c}$ is a universal constant, estimated at (5.18).
Proof. - We know that $m_{l c}(t \theta)=\widehat{\varphi_{\theta}}(t)$. From Lemma 5.8, (5.15d) with $k=0$, it follows that

$$
\left|\frac{\mathrm{d}^{j}}{\mathrm{~d} t^{j}} m_{l c}(t \theta)\right|=\left|{\widehat{\varphi_{\theta}}}^{(j)}(t)\right| \leqslant(2 \pi)^{j} \int_{\mathbb{R}}|s|^{j} \varphi_{\theta}(s) \mathrm{d} s
$$

and with $k=1$,

$$
\left|\frac{\mathrm{d}^{j}}{\mathrm{~d} t^{j}} m_{l c}(t \theta)\right| \leqslant \frac{(2 \pi)^{j-1}}{|t|}\left(j \int_{\mathbb{R}}|s|^{j-1} \varphi_{\theta}(s) \mathrm{d} s+\int_{\mathbb{R}}|s|^{j}\left|\varphi_{\theta}^{\prime}(s)\right| \mathrm{d} s\right)
$$

The function $\varphi_{\theta}$ is a symmetric log-concave probability density on $\mathbb{R}$, with variance 1 . By Corollary 5.4, we have for $j=0$ that

$$
(1+2 \pi|t|)\left|m_{l c}(t \theta)\right| \leqslant \int_{\mathbb{R}}\left(\varphi_{\theta}(u)+\left|\varphi_{\theta}^{\prime}(u)\right|\right) \mathrm{d} u \leqslant 1+2 \varphi_{\theta}(0) \leqslant 1+\sqrt{2}
$$

For $j \geqslant 1$, we have $\int_{\mathbb{R}}|u|^{j}\left|\varphi_{\theta}{ }^{\prime}(u)\right| \mathrm{d} u=j \int_{\mathbb{R}}|u|^{j-1} \varphi_{\theta}(u) \mathrm{d} u$ by Lemma 5.9, and

$$
(1+2 \pi|t|)\left|\frac{\mathrm{d}^{j}}{\mathrm{~d} t^{j}} m_{l c}(t \theta)\right| \leqslant(2 \pi)^{j} \int_{\mathbb{R}}\left(|u|^{j}+2 j|u|^{j-1}\right) \varphi_{\theta}(u) \mathrm{d} u
$$

The function $\varphi_{\theta}$ satisfies $\int_{\mathbb{R}} s^{2} \varphi_{\theta}(s) \mathrm{d} s=1$, implying that

$$
\begin{equation*}
\delta_{0, c} \leqslant 1+\sqrt{2}<3 ; \quad \delta_{1, c} \leqslant 6 \pi ; \quad \delta_{2, c} \leqslant 20 \pi^{2} \tag{5.18a}
\end{equation*}
$$

We know by Lemma 5.6 that $\varphi_{\theta}(s) \leqslant 11 \mathrm{e}^{-|s|}$. This implies for $j>2$ that

$$
\begin{equation*}
\delta_{j, c} \leqslant 22(2 \pi)^{j} \int_{0}^{+\infty}\left(s^{j}+2 j s^{j-1}\right) \mathrm{e}^{-s} \mathrm{~d} s=66(2 \pi)^{j} \Gamma(j+1) \tag{5.18b}
\end{equation*}
$$

Remarks 5.12. - One gets $\int_{\mathbb{R}} \mid s^{j} \varphi_{\theta}(s) \mathrm{d} s \leqslant 3^{j / 2} \Gamma(j+1)$ by applying Lemma 5.3 and Corollary 5.4; Lemma 5.6 yields the bound $22 \Gamma(j+1)$, better when $j$ is large.

If the log-concave probability density $K$ on $\mathbb{R}^{n}$ is normalized by variance but is simply centered, then $\varphi_{\theta, K}$ is log-concave and centered for each $\theta$, and satisfies the exponential decay of Corollary 5.7. If $\varphi_{\theta, K}$ reaches its maximum at $s_{0}$, then

$$
\int_{\mathbb{R}}|s|^{j}\left|\varphi_{\theta, K}^{\prime}(s)\right| \mathrm{d} s \leqslant 2\left|s_{0}\right|^{j} \varphi_{\theta, K}\left(s_{0}\right)+j \int_{\mathbb{R}}|s|^{j-1} \varphi_{\theta, K}(s) \mathrm{d} s
$$

admits a universal bound $\kappa_{j}$. Lemma 5.11 remains valid in this extended case, with other constants $\left(\delta_{j}\right)_{j \geqslant 0}$ for which we do not have satisfactory explicit expressions. Fradelizi [34, Theorem 5] extended the $L^{p}\left(\mathbb{R}^{n}\right)$ result of Theorem 6.2 (Bourgain, Carbery) to centered bodies $C$ in $\mathbb{R}^{n}$, not necessarily symmetric (unluckily, the word "centered" was forgotten in the statement given in [34]).

If $C$ is an arbitrary convex body, then $\mathrm{M}_{C}$ is bounded on $L^{p}\left(\mathbb{R}^{n}\right), p \in$ $(1,+\infty]$, but for each fixed $n \geqslant 1$ and $p<+\infty$, there is no uniform bound for the family of arbitrary convex bodies in $\mathbb{R}^{n}$ (if $n=1$, examine $\mathrm{M}_{C} f$ when $C=[1,1+\varepsilon], f=\mathbf{1}_{C}$ and $\left.\varepsilon \rightarrow 0\right)$. In a somewhat related direction, it is known that the $L^{p}\left(\mathbb{R}^{n}\right)$ norm of the uncentered operator in (0.2) is $\geqslant C_{p}^{n}$ for some $C_{p}>1$, when $1<p<+\infty$ [40].
Corollary 5.13. - Let $K_{l c}$ be a symmetric log-concave probability density on $\mathbb{R}^{n}$, isotropic with variance $\sigma^{2}$, and let $m_{l c}$ be its Fourier transform. For every $\xi \in \mathbb{R}^{n}$ and $j \geqslant 0$ one has that

$$
\begin{equation*}
\left|\frac{\mathrm{d}^{j}}{\mathrm{~d} t^{j}} m_{l c}(t \xi)\right| \leqslant \delta_{j, c} \frac{|\sigma \xi|^{j}}{1+2 \pi|t \sigma \xi|}, \quad t \in \mathbb{R} \tag{5.19}
\end{equation*}
$$

where $\delta_{j, c}$ is the universal constant of Lemma 5.11.
Proof. - The result is obvious when $\xi=0$, otherwise we apply Lemma 5.11 with $\theta=|\xi|^{-1} \xi$ to the normalized Fourier transform $N(\xi)=$ $m_{l c}(\xi / \sigma)$, obtaining thus

$$
\frac{\mathrm{d}^{j}}{\mathrm{~d} t^{j}} m_{l c}(t \xi)=\frac{\mathrm{d}^{j}}{\mathrm{~d} t^{j}} N(t|\sigma \xi| \theta)=\left.|\sigma \xi|^{j} \frac{\mathrm{~d}^{j}}{\mathrm{~d} u^{j}} N(u \theta)\right|_{u=t|\sigma \xi|} \leqslant \delta_{j, c} \frac{|\sigma \xi|^{j}}{1+2 \pi|t \sigma \xi|}
$$

### 5.3. Fourier analysis in $L^{2}\left(\mathbb{R}^{n}\right)$

Lemma 5.14 (Bourgain [9]). - Let $K$ be a kernel in $L^{1}\left(\mathbb{R}^{n}\right)$ and assume that its Fourier transform $m$ is $C^{1}$ outside the origin. For every $j \in \mathbb{Z}$, define

$$
\alpha_{j}(m)=\sup _{2^{j-1} \leqslant|\xi| \leqslant 2^{j+1}}|m(\xi)| \quad \text { and } \quad \beta_{j}(m)=\sup _{2^{j-1} \leqslant|\xi| \leqslant 2^{j+1}}|\xi \cdot \nabla m(\xi)|
$$

If $\Gamma_{B}(K):=\sum_{j \in \mathbb{Z}} \sqrt{\alpha_{j}(m)} \sqrt{\alpha_{j}(m)+\beta_{j}(m)}<+\infty$, then the maximal operator $M_{K}$ associated to $K$ is bounded on $L^{2}\left(\mathbb{R}^{n}\right)$. More precisely, one has that

$$
\left\|M_{K} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}=\left\|\sup _{t>0}\left|K_{(t)} * f\right|\right\|_{L^{2}\left(\mathbb{R}^{n}\right)} \leqslant 2 \Gamma_{B}(K)\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)}, \quad f \in L^{2}\left(\mathbb{R}^{n}\right)
$$

We shall simply write $\alpha_{j}=\alpha_{j}(m)$ and $\beta_{j}=\beta_{j}(m)$ in the rest of the section.

Remark. - Clearly, we have that

$$
\sum_{j \in \mathbb{Z}} \sqrt{\alpha_{j}} \sqrt{\alpha_{j}+\beta_{j}} \leqslant \sum_{j \in \mathbb{Z}} \alpha_{j}+\sum_{j \in \mathbb{Z}} \sqrt{\alpha_{j} \beta_{j}},
$$

and each of the two terms in the right-hand side is less than the left-hand side. Bourgain employs both expressions as definitions of $\Gamma_{B}(K)$, one in [9] and the other in [10] or in [13]. The convergence of the series of $\alpha_{j} \mathrm{~S}$ when $j$ tends to $-\infty$ implies that $m(\xi)$ tends to 0 when $\xi$ tends to 0 , thus $m(0)=0$, which means that the integral of $K$ on $\mathbb{R}^{n}$ is equal to 0 . This lemma will not be applied to $K_{C}$ or $K_{l c}$, but typically, to the difference of two kernels with equal integrals.

Proof. - We shall give a proof less rough than Bourgain's, relying on the tools introduced in Section 4. We consider a $C^{\infty}$ function $\eta$ on $\mathbb{R}$ such that

$$
\eta(t)=1 \text { if } t \leqslant 1, \quad \eta(t)=0 \text { if } t \geqslant 2, \quad \text { and } 0 \leqslant \eta \leqslant 1 .
$$

Next, we set $\rho(t)=\eta(t)-\eta(2 t)$ for $t \in \mathbb{R}$. We see that $\rho$ vanishes outside $[1 / 2,2]$. Also, $\rho(t)=1-\eta(2 t)$ on $[1 / 2,1]$ and $\rho(t)=\eta(t)$ on $[1,2]$, so that $0 \leqslant \rho(t) \leqslant 1$ and

$$
d_{0}:=\sup _{t \in \mathbb{R}}\left|t \rho^{\prime}(t)\right|=\sup _{t \in \mathbb{R}}\left|t \eta^{\prime}(t)\right|=\sup _{t \in[1,2]} t\left|\eta^{\prime}(t)\right| .
$$

Let $\varepsilon>0$ be given. One can make sure that $d_{0}<(1+\varepsilon) / \ln 2$, choosing for $\eta$ a $C^{\infty}$ approximation of the function $\eta_{0}$ defined on $[0,2]$ by $\eta_{0}(t)=$ $\min \left(1,1-\log _{2} t\right)$, for which $t\left|\eta_{0}^{\prime}(t)\right|=1 / \ln (2)$ when $t \in[1,2]$.

For every $j \in \mathbb{Z}$ and $\xi \in \mathbb{R}^{n}$, let $\varphi_{j}(\xi)=\rho\left(2^{-j}|\xi|\right)$ and consider the annulus

$$
C_{j}=\left\{\xi \in \mathbb{R}^{n}: 2^{j-1} \leqslant|\xi| \leqslant 2^{j+1}\right\} \subset \mathbb{R}^{n} .
$$

From the properties of $\rho$, we have that $0 \leqslant \varphi_{j} \leqslant 1, \varphi_{j}$ vanishes outside $C_{j}$, and

$$
\sum_{j \in \mathbb{Z}} \varphi_{j}(\xi)=\sum_{j \in \mathbb{Z}}\left(\eta\left(2^{-j}|\xi|\right)-\eta\left(2^{-j+1}|\xi|\right)\right)=1
$$

for every $\xi \neq 0$, because $\eta\left(2^{-j}|\xi|\right)=0$ when $j \leqslant \log _{2}(|\xi|)-1$ and $\eta\left(2^{-j}|\xi|\right)=$ 1 when $j \geqslant \log _{2}(|\xi|)$. We introduce for every $j \in \mathbb{Z}$ a multiplier $m_{j}$ defined by

$$
m_{j}(\xi)=\varphi_{j}(\xi) m(\xi), \quad \xi \in \mathbb{R}^{n}
$$

and we let $K_{j}=m_{j}^{\vee}=\varphi_{j}^{\vee} * K$. One has $\sum_{j \in \mathbb{Z}} K_{j}=K$, which allows us to write for $f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$ and every $x \in \mathbb{R}^{n}$ the upper bound

$$
\left(M_{K} f\right)(x)=\sup _{t>0}\left|\left(K_{(t)} * f\right)(x)\right| \leqslant \sup _{t>0} \sum_{j \in \mathbb{Z}}\left|\left[\left(K_{j}\right)_{(t)} * f\right](x)\right| \leqslant \sum_{j \in \mathbb{Z}}\left(M_{K_{j}} f\right)(x) .
$$

By Lemma 4.7 with $r=4$, one has

$$
\begin{equation*}
\left\|M_{K_{j}} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}^{2} \leqslant 2 \ln 4\left\|m_{j}\right\|_{L^{\infty}\left(\mathbb{R}^{n}\right)}\left\|m_{j}^{*}\right\|_{L^{\infty}\left(\mathbb{R}^{n}\right)}\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)}^{2} . \tag{5.20}
\end{equation*}
$$

We see that $\left\|m_{j}\right\|_{\infty} \leqslant \alpha_{j}$, since $\left|m_{j}\right| \leqslant|m|$ and since $m_{j}$ is supported in the annulus $C_{j}$. On the other hand, $m_{j}^{*}(\xi)=\xi \cdot \nabla m_{j}(\xi)$ and we have

$$
\nabla m_{j}(\xi)=\varphi_{j}(\xi) \nabla m(\xi)+m(\xi) \nabla \varphi_{j}(\xi)
$$

As $\varphi_{j}$ is supported in $C_{j}$, we get $\left|\varphi_{j}(\xi) \xi \cdot \nabla m(\xi)\right| \leqslant \beta_{j}<(1+\varepsilon) \beta_{j} / \ln 2$, and

$$
\left|m(\xi) \xi \cdot \nabla \varphi_{j}(\xi)\right| \leqslant \alpha_{j}\left|\xi \cdot 2^{-j} \rho^{\prime}\left(2^{-j}|\xi|\right) \frac{\xi}{|\xi|}\right| \leqslant \alpha_{j} d_{0}<(1+\varepsilon) \alpha_{j} / \ln 2 .
$$

It follows that $\left\|m_{j}^{*}\right\|_{\infty} \leqslant(1+\varepsilon)\left(\alpha_{j}+\beta_{j}\right) / \ln 2$. By (5.20) we get

$$
\left\|M_{K_{j}} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)} \leqslant 2 \sqrt{1+\varepsilon} \sqrt{\alpha_{j}} \sqrt{\alpha_{j}+\beta_{j}}\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)}
$$

After summation in $j \in \mathbb{Z}$ and letting $\varepsilon \rightarrow 0$, we conclude that

$$
\left\|M_{K} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)} \leqslant 2 \Gamma_{B}(K)\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)}
$$

We pass from $f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$ to $f \in L^{2}\left(\mathbb{R}^{n}\right)$ as explained in Section 3.3.

### 5.3.1. Conclusion of Bourgain's argument

End of the proof of Theorem 5.1. - We begin with a version of the proof that illustrates well the fact that Lemma 5.14 is a comparison lemma: in vague terms, if we know that the conclusion of Theorem 5.1 is true for one family of convex sets, then it is true for all convex sets.

We rely here on Stein's Theorem 4.1 for the Euclidean ball $B$, asserting that the maximal operator $\mathrm{M}_{B}$ is bounded on $L^{p}\left(\mathbb{R}^{n}\right)$ for every $p$ in $(1,+\infty]$, with a bound independent of the dimension $n$. In this paragraph, we only use the $L^{2}$ case of this result. Let us call $B=B_{n, V}$ the Euclidean ball in $\mathbb{R}^{n}$, centered at 0 and normalized by variance, which has radius $\sqrt{n+2}$ by (5.4). Let $m_{B}$ denote the Fourier transform of $K_{B}$. Consider also a symmetric log-concave probability density $K_{l c}$ on $\mathbb{R}^{n}$, isotropic and normalized by variance. The two functions $m_{l c}$ and $m_{B}$ satisfy the estimates (5.17.B) of Proposition 5.10. We apply Lemma 5.14 to the difference kernel $K=K_{l c}-K_{B}$. According to (5.17.B), for every $\xi \in \mathbb{R}^{n}$, the Fourier transform $m=m_{l c}-m_{B}$ satisfies

$$
\begin{aligned}
|\xi||m(\xi)| \leqslant \sqrt{2} / \pi,|m(\xi)| \leqslant\left|1-m_{l c}(\xi)\right|+\left|1-m_{B}(\xi)\right| \leqslant 4 \pi|\xi| \\
|\xi \cdot \nabla m(\xi)| \leqslant 4
\end{aligned}
$$

We deduce that $\beta_{j}=\sup _{2^{j-1} \leqslant|\xi| \leqslant 2^{j+1}}|\xi \cdot \nabla m(\xi)| \leqslant 4$ for $j \in \mathbb{Z}$. For $j<0$ one has

$$
\alpha_{j}=\sup _{2^{j-1} \leqslant|\xi| \leqslant 2^{j+1}}|m(\xi)| \leqslant 4 \pi 2^{j+1}=4 \pi 2^{-|j|+1} \leqslant 32.2^{-|j|}
$$

and for $j \geqslant 0$, we have $\alpha_{j} \leqslant \sqrt{2} \pi^{-1} 2^{-j+1} \leqslant 2^{-j}$. It follows that the two series $\sum_{j \in \mathbb{Z}} \alpha_{j}$ and $\sum_{j \in \mathbb{Z}} \sqrt{\alpha_{j} \beta_{j}}$ converge, and

$$
\sum_{j \in \mathbb{Z}} \alpha_{j} \leqslant 32+2, \quad \sum_{j \in \mathbb{Z}} \sqrt{\alpha_{j} \beta_{j}} \leqslant 20+10 \sqrt{2}
$$

thus the maximal operator $f \mapsto \sup _{t>0}\left|K_{(t)} * f\right|$ is bounded on $L^{2}\left(\mathbb{R}^{n}\right)$ by a constant independent of the dimension, say, less than $2 \Gamma_{B}(K)<2(54+$ $10 \sqrt{2})<137$. Finally, for $f \geqslant 0$, we write

$$
\begin{aligned}
\mathrm{M}_{K_{l c}} f & =\sup _{t>0}\left|\left(K_{l c}\right)_{(t)} * f\right| \\
& \leqslant \sup _{t>0}\left|\left(K_{B}\right)_{(t)} * f\right|+\sup _{t>0}\left|\left(K_{l c}-K_{B}\right)_{(t)} * f\right|=\mathrm{M}_{B} f+M_{K} f
\end{aligned}
$$

and we can estimate $\mathrm{M}_{K_{l c}}$ by the sum of two operators that are bounded on $L^{2}\left(\mathbb{R}^{n}\right)$ by constants independent of the dimension $n$.

The proof actually given by Bourgain [9] bypasses the $L^{2}$ result of Stein on Euclidean balls. The kernel $K$ is now given as $K=K_{l c}-P$, where $P$ is the Poisson kernel $P=P_{1}$ from (1.32) for the value $t=1$ of the parameter. We know by $\left(1.31 . P^{*}\right)$ that the maximal operator $f \mapsto \sup _{t>0}\left|P_{t} f\right|$ associated to the Poisson kernel acts boundedly on $L^{p}\left(\mathbb{R}^{n}\right), 1<p \leqslant+\infty$, with a bound $\leqslant 2$ when $p=2$, thus independent of the dimension $n$. Now, everything is said: we replace the multiplier $m_{B}$ by $\widehat{P}$ and it suffices to see that $\widehat{P}$ also satisfies good estimates similar to (5.17.B). But $\widehat{P}(\xi)=\mathrm{e}^{-2 \pi|\xi|}$ clearly satisfies the even better estimates

$$
\begin{gather*}
|\xi||\widehat{P}(\xi)|=|\xi| \mathrm{e}^{-2 \pi|\xi|} \leqslant(2 \pi \mathrm{e})^{-1}  \tag{5.21a}\\
|1-\widehat{P}(\xi)| \leqslant 2 \pi|\xi|, \quad|\xi \cdot \nabla \widehat{P}(\xi)|=2 \pi|\xi| \mathrm{e}^{-2 \pi|\xi|} \leqslant \mathrm{e}^{-1} \tag{5.21b}
\end{gather*}
$$

where we made use of the inequality $x \mathrm{e}^{-x} \leqslant \mathrm{e}^{-1}$, true for every $x \geqslant 0$. This ends the second proof of Theorem 5.1, with different constants whose exact values are rather irrelevant. However, we found here an explicit bound $\kappa_{2}<2+137<140$, explicit but definitely not sharp.

## 6. The $L^{p}$ results of Bourgain and Carbery

One gives again a symmetric convex body $C$ in $\mathbb{R}^{n}$, and $\mu_{C}$ denotes the uniform probability measure on $C$. Beside the maximal function $\mathrm{M}_{C} f$ from (0.3.M $)$, for every function $f \in L_{\text {loc }}^{1}\left(\mathbb{R}^{n}\right)$ and every $x \in \mathbb{R}^{n}$ we set

$$
\left(\mathrm{M}_{C}^{(d)} f\right)(x)=\sup _{j \in \mathbb{Z}} \frac{1}{\left|2^{j} C\right|} \int_{x+2^{j} C}|f(y)| \mathrm{d} y=\sup _{j \in \mathbb{Z}} \int_{\mathbb{R}^{n}}\left|f\left(x+2^{j} v\right)\right| \mathrm{d} \mu_{C}(v)
$$

One can call $\mathrm{M}_{C}^{(d)} f$ the "dyadic" maximal function associated to the convex set $C$. Obviously, $\mathrm{M}_{C}^{(d)} \leqslant \mathrm{M}_{C}$. More generally, we associate to every kernel $K$ integrable on $\mathbb{R}^{n}$ the dyadic maximal function

$$
\left(M_{K}^{(d)} f\right)(x)=\sup _{j \in \mathbb{Z}}\left|\int_{\mathbb{R}^{n}} f\left(x+2^{j} v\right) K(v) \mathrm{d} v\right|, \quad x \in \mathbb{R}^{n}
$$

In 1986, Bourgain and Carbery have obtained identical results for $L^{p}\left(\mathbb{R}^{n}\right)$. Somewhat surprisingly, the cases $\mathrm{M}_{C}^{(d)}$ and $\mathrm{M}_{C}$ are different, the boundedness of $\mathrm{M}_{C}$ on $L^{p}\left(\mathbb{R}^{n}\right)$ being obtained only when $p>3 / 2$, as opposed to $p>1$ for $\mathrm{M}_{C}^{(d)}$.

Theorem 6.1 (Bourgain [10], Carbery [21]). - For every p in $(1,+\infty]$, there exists a constant $\kappa^{(d)}(p)$ such that for every integer $n \geqslant 1$ and every symmetric convex body $C \subset \mathbb{R}^{n}$, one has

$$
\forall f \in L^{p}\left(\mathbb{R}^{n}\right),\left\|\mathrm{M}_{C}^{(d)} f\right\|_{L^{p}\left(\mathbb{R}^{n}\right)} \leqslant \kappa^{(d)}(p)\|f\|_{L^{p}\left(\mathbb{R}^{n}\right)}
$$

Theorem 6.2 (Bourgain [10], Carbery [21]). - For every $p$ in $(3 / 2,+\infty]$, there exists a constant $\kappa(p)$ such that for every integer $n \geqslant 1$ and for every symmetric convex set $C \subset \mathbb{R}^{n}$, one has that

$$
\forall f \in L^{p}\left(\mathbb{R}^{n}\right),\left\|\mathrm{M}_{C} f\right\|_{L^{p}\left(\mathbb{R}^{n}\right)} \leqslant \kappa(p)\|f\|_{L^{p}\left(\mathbb{R}^{n}\right)}
$$

We recalled in the Introduction that the maximal theorem of strong type is not true for $p=1$, even with a constant depending on $n$, and even for the smaller function $\mathrm{M}_{C}^{(d)} f$, since $\mathrm{M}_{C} f \leqslant 2^{n} \mathrm{M}_{C}^{(d)} f$. Note that Theorems 6.1 and 6.2 are obvious for $L^{\infty}\left(\mathbb{R}^{n}\right)$, with $\kappa^{(d)}(\infty)=\kappa(\infty)=1$. By Bourgain [9], we have the result in $L^{2}\left(\mathbb{R}^{n}\right)$, so we obtain it for $p \in[2,+\infty]$ by interpolation. Consequently, our work will be limited to values of $p$ in the interval $(1,2]$. We shall follow Carbery's approach to both theorems. This approach has been applied later in the Detlef Müller article [59] (see Section 7), on which relies Bourgain's recent article [13] devoted to the maximal function associated to high dimensional cubes (see Section 8).

The proof will use the inequalities (5.17.B) and (5.19), which are also true for log-concave densities, and by simply following the proofs of Bourgain or Carbery, we can extend the results to the log-concave setting. As suggested in [10], one can actually take one more step, forget convexity and exploit only the inequalities on the Fourier transform given by Lemma 5.11. In this more general framework, we consider a probability density $K_{g}$ on $\mathbb{R}^{n}$, or merely a kernel $K_{g}$ integrable on $\mathbb{R}^{n}$ and having a Fourier transform $m_{g}$ which satisfies the following: there exist $\delta_{0, g}, \delta_{1, g}>0$ such that for every
$\theta \in S^{n-1}$, we have

$$
\begin{equation*}
\left|m_{g}(t \theta)\right| \leqslant \frac{\delta_{0, g}}{1+|t|}, \quad\left|\frac{\mathrm{d}}{\mathrm{~d} t} m_{g}(t \theta)\right|=\left|\theta \cdot \nabla m_{g}(t \theta)\right| \leqslant \frac{\delta_{1, g}}{1+|t|}, \quad t \in \mathbb{R} \tag{6.1.H}
\end{equation*}
$$

The form of the $\delta_{0, g}$-bound of $m_{g}$ has been chosen for the sake of uniformity, but when $K_{g}$ is a probability density, we know of course that $\left\|m_{g}\right\|_{L^{\infty}\left(\mathbb{R}^{n}\right)}=1$ and in particular we have $\delta_{0, g} \geqslant 1$ in that case.

Proposition 6.3. - Theorems 6.1 and 6.2 are also valid for any symmetric log-concave probability density $K_{l c}$ on $\mathbb{R}^{n}$, namely

$$
\begin{array}{lr}
\left\|\mathrm{M}_{K_{l c}}^{(d)} f\right\|_{L^{p}\left(\mathbb{R}^{n}\right)} \leqslant \kappa^{(d)}(p)\|f\|_{L^{p}\left(\mathbb{R}^{n}\right)}, & 1<p \leqslant+\infty \\
\left\|\mathrm{M}_{K_{l c}} f\right\|_{L^{p}\left(\mathbb{R}^{n}\right)} \leqslant \kappa(p)\|f\|_{L^{p}\left(\mathbb{R}^{n}\right)}, & 3 / 2<p \leqslant+\infty
\end{array}
$$

If a probability density $K_{g}$ satisfies $(6.1 . \mathbf{H})$, then for $3 / 2<p \leqslant 2$ we have

$$
\left\|\mathrm{M}_{K_{g}}\right\|_{p \rightarrow p} \leqslant \kappa_{p}\left(\delta_{0, g}+\delta_{1, g}\right)^{2-2 / p}
$$

and this result extends to every $p \in(1,2]$ in the case of the dyadic operator $\mathrm{M}_{K_{g}}^{(d)}$.

All these results are obvious when $p=+\infty$, and easy when $p>2$ by interpolation $\left(L^{2}, L^{\infty}\right)$ after the case $p=2$ is obtained. When $p \leqslant 2$, the log-concave statements follow from the "general" one. Indeed, for the study of maximal functions, we may assume that the convex set $C$ or the symmetric log-concave probability density $K_{l c}$ is isotropic and normalized by variance. Then, by (5.17.B) or by Lemma 5.11, $m_{C}$ or $m_{l c}$ satisfy (6.1.H) with universal constants $\delta_{0, c}$ and $\delta_{1, c}$.

### 6.1. A priori estimate and interpolation

Suppose that a family $\left(T_{j}\right)_{j \in \mathbb{Z}}$ of operators on $L^{p}(X, \Sigma, \mu)$ is given, for a set of values of $p$ and on a certain measure space $(X, \Sigma, \mu)$ (further down, it will be $\mathbb{R}^{n}$, equipped with the Lebesgue measure). These operators can be linear operators, or nonlinear operators of the form

$$
T_{j} f=\sup _{v \in V}\left|T_{j, v} f\right|
$$

where each $T_{j, v}$ is linear and where $v$ runs over a certain set $V$ of indices. We want to study the maximal function

$$
T^{*} f=\sup _{j \in \mathbb{Z}}\left|T_{j} f\right|=\sup _{j \in \mathbb{Z}, v \in V}\left|T_{j, v} f\right|
$$

We also consider later a kernel $K$ integrable on $\mathbb{R}^{n}$. In the application to the geometrical problem, this kernel will be (as in Section 5.3.1) the difference $K=K_{1}-K_{2}$ of two kernels, where $K_{1}$ is the uniform probability density on an isotropic convex set $C$ or a probability density $K_{g}$ satisfying (6.1.H), and $K_{2}$ is a kernel for which the dimensionless maximal inequality is already known. We have to deal with two cases. In the first one, $T_{j}$ will be the convolution with the dilate $K_{\left(2^{j}\right)}$ from (2.7) of $K$, and the maximal function $T^{*} f=M_{K}^{(d)} f$ will then permit us to relate the dyadic maximal function $\mathrm{M}_{C}^{(d)} f$ to a maximal function whose bounded character on $L^{p}\left(\mathbb{R}^{n}\right)$ is already known. In the second one, the operator $T_{j, v}$ will be the convolution with $K_{\left(v 2^{j}\right)}$ with $v \in[1,2]=V$, in which case

$$
\begin{equation*}
T_{j} f=\sup _{2^{j} \leqslant t \leqslant 2^{j+1}}\left|K_{(t)} * f\right| \tag{6.2}
\end{equation*}
$$

and $T^{*} f=M_{K} f$ allows us to study the "global" maximal function $\mathrm{M}_{C} f$ or $\mathrm{M}_{K_{g}} f$.

We assume that linear operators $\left(Q_{j}\right)_{j \in \mathbb{Z}}$ such that $\sum_{j \in \mathbb{Z}} Q_{j}=$ Id are given. In the applications to come, these operators will be those of Equation (2.6), in the Section 2.1 on Littlewood-Paley functions.

Definition 6.4 (Carbery [21]). - Given families $\left(T_{j}\right)_{j \in \mathbb{Z}}$ and $\left(Q_{j}\right)_{j \in \mathbb{Z}}$ as above, we say that $T^{*}$ is weakly bounded on $L^{p}(X, \Sigma, \mu)$ if there exists a constant $A$ such that

$$
\forall f \in L^{p}(X, \Sigma, \mu), \quad \forall k \in \mathbb{Z}, \quad\left\|\sup _{j \in \mathbb{Z}}\left|T_{j} Q_{j+k} f\right|\right\|_{L^{p}(\mu)} \leqslant A\|f\|_{L^{p}(\mu)} \cdot\left(\mathbf{W}_{p}\right)
$$

We say that $T^{*}$ is strongly bounded on $L^{p}(X, \Sigma, \mu)$ if there exists a real nonnegative sequence $\left(a_{k}\right)_{k \in \mathbb{Z}}$, satisfying $\sum_{k \in \mathbb{Z}} a_{k}^{r}<+\infty$ for every $r>0$, and such that

$$
\begin{equation*}
\forall f \in L^{p}(X, \Sigma, \mu), \quad \forall k \in \mathbb{Z}, \quad\left\|\sup _{j \in \mathbb{Z}}\left|T_{j} Q_{j+k} f\right|\right\|_{L^{p}(\mu)} \leqslant a_{k}\|f\|_{L^{p}(\mu)} \tag{p}
\end{equation*}
$$

By $T_{j} Q_{j+k} f$, we mean of course $T_{j}\left(Q_{j+k} f\right)$.
Remarks 6.5. - In this generality, the supremum for $v \in V$ in $T_{j} f=$ $\sup _{v \in V}\left|T_{j, v} f\right|$ must be understood as essential supremum, as explained in Section 3.3. In our cases of application, the function $v \mapsto T_{j, v}(x), x \in X$, will be a continuous function on an interval $V$ of the line, in which case the pointwise supremum coincides with the supremum on any countable dense subset of $V$.

It is evident that $\left(\mathbf{S}_{p}\right)$ implies $\left(\mathbf{W}_{p}\right)$, and $\left(\mathbf{S}_{p}\right)$ implies that $T^{*}$ is bounded, because

$$
\left|T_{j, v} f\right|=\left|\sum_{k \in \mathbb{Z}} T_{j, v} Q_{j+k} f\right| \leqslant \sum_{k \in \mathbb{Z}}\left|T_{j, v} Q_{j+k} f\right| \leqslant \sum_{k \in \mathbb{Z}}\left|T_{j} Q_{j+k} f\right|
$$

thus

$$
\left|T_{j} f\right|=\sup _{v \in V}\left|T_{j, v} f\right| \leqslant \sum_{k \in \mathbb{Z}}\left|T_{j} Q_{j+k} f\right|, \quad \text { then } \quad T^{*} f \leqslant \sum_{k \in \mathbb{Z}} \sup _{j \in \mathbb{Z}}\left|T_{j} Q_{j+k} f\right|
$$

and

$$
\begin{equation*}
\left\|T^{*} f\right\|_{L^{p}(\mu)} \leqslant \sum_{k \in \mathbb{Z}}\left\|\sup _{j \in \mathbb{Z}}\left|T_{j} Q_{j+k} f\right|\right\|_{L^{p}(\mu)} \leqslant\left(\sum_{k \in \mathbb{Z}} a_{k}\right)\|f\|_{L^{p}(\mu)} \tag{6.3}
\end{equation*}
$$

If one has $\left(\mathbf{W}_{p_{0}}\right)$ and $\left(\mathbf{S}_{p_{1}}\right)$ and if $1 / p=(1-\theta) / p_{0}+\theta / p_{1}$, with $0<\theta \leqslant 1$, then as in (3.26) we obtain by interpolation

$$
\forall f \in L^{p}(\mu), \forall k \in \mathbb{Z}, \quad\left\|\sup _{j \in \mathbb{Z}}\left|T_{j} Q_{j+k} f\right|\right\|_{L^{p}(\mu)} \leqslant A^{1-\theta} a_{k}^{\theta}\|f\|_{L^{p}(\mu)}
$$

and $\sum_{k \in \mathbb{Z}} A^{(1-\theta) r} a_{k}^{\theta r}<+\infty$ for every $r>0$, so $\left(\mathbf{S}_{p}\right)$ is satisfied. In order to obtain this, we apply the complex interpolation of linear operators between spaces $L^{p}\left(\ell^{q}\right)$ [7, Chap. 5, Th. 5.1.2]. Here, the range space is of the form $L^{p}\left(\mu, \ell^{\infty}(\mathbb{Z})\right)$, a case covered by complex interpolation. Indeed, in the simpler case where the $T_{j}$ s are linear, we obtain the result by considering for each $k \in \mathbb{Z}$ the linear operator

$$
f \mapsto\left(T_{j} Q_{j+k} f\right)_{j \in \mathbb{Z}} \in L^{p}\left(X, \Sigma, \mu, \ell^{\infty}(\mathbb{Z})\right), \quad f \in L^{p}(\mu)
$$

If $V$ has more than one element, the range space will be $L^{p}\left(\mu, \ell^{\infty}(\mathbb{Z} \times V)\right)$. The nonlinear operator $f \mapsto \sup _{j \in \mathbb{Z}}\left|T_{j} Q_{j+k} f\right|$ belongs to the class of linearizable operators considered in [36].

We now describe the assumptions that will be made in the main result of this section. First of all, we assume that there exist constants $C_{p}, 1<p \leqslant 2$, such that

$$
\begin{equation*}
\forall p \in(1,2], \quad \forall f \in L^{p}(\mu), \quad\left\|\left(\sum_{j \in \mathbb{Z}}\left|Q_{j} f\right|^{2}\right)^{1 / 2}\right\|_{L^{p}(\mu)} \leqslant C_{p}\|f\|_{L^{p}(\mu)} \tag{0}
\end{equation*}
$$

If the $\left(Q_{j}\right)_{j \in \mathbb{Z}}$ are those of (2.4), then we can take $C_{p}=\mathrm{q}_{p}$ which behaves as $1 /(p-1)$ when $p \rightarrow 1$, according to (2.5).

We assume that $T_{j, v}=U_{j, v}-S_{j, v}$, where $U_{j, v}$ and $S_{j, v}$ are positive linear operators, and we assume for $S^{*}$, defined by $S^{*} f=\sup _{j \in \mathbb{Z}, v \in V}\left|S_{j, v} f\right|$, that there exist $p_{\text {min }}$ in the open interval $(1,2)$ and constants $C_{p}^{\prime}, p_{\min }<p \leqslant 2$, such that

$$
\begin{equation*}
\forall p \in\left(p_{\min }, 2\right], \quad\left\|S^{*}\right\|_{p} \leqslant C_{p}^{\prime} \tag{1}
\end{equation*}
$$

where $\|R\|_{p}$ is a shorter notation for the norm $\|R\|_{p \rightarrow p}$ of an operator $R$. The condition " $U_{j, v}$ positive" will be the only reason for requiring that the kernel $K_{g}$ in Proposition 6.3 be a probability density rather than an arbitrary integrable kernel. The $U_{j, v} \mathrm{~s}$ will correspond to the kernel $K_{g}$ under study, while the $S_{j, v} \mathrm{~s}$ will often refer to Poisson kernels for which the maximal function estimates in $L^{p}\left(\mathbb{R}^{n}\right)$ are already known by (1.31.P $P^{*}$ ).

We assume that for every $p \in\left(p_{\text {min }}, 2\right]$, there exists a constant $C_{p}^{\prime \prime}$ such that

$$
\begin{equation*}
\forall j \in \mathbb{Z}, \quad\left\|T_{j}\right\|_{p} \leqslant C_{p}^{\prime \prime} \tag{2}
\end{equation*}
$$

We shall assume that $T^{*}$ satisfies $\left(\mathbf{S}_{2}\right)$, hence we have that

$$
\begin{equation*}
\forall f \in L^{2}(\mu), \forall k \in \mathbb{Z},\left\|\sup _{j \in \mathbb{Z}}\left|T_{j} Q_{j+k} f\right|\right\|_{L^{2}(\mu)} \leqslant a_{k}\|f\|_{L^{2}(\mu)} \tag{3}
\end{equation*}
$$

where $\sum_{k \in \mathbb{Z}} a_{k}^{r}<+\infty$ for every $r>0$.
Proposition 6.6 (Carbery [21]). - Under the assumptions $\left(A_{\mathbf{0}}\right),\left(A_{\mathbf{1}}\right)$, $\left(A_{\mathbf{2}}\right)$ and $\left(A_{\mathbf{3}}\right)$, the maximal operator $T^{*}$ is bounded on $L^{p}(X, \Sigma, \mu)$ for every real number $p$ such that $p_{\min }<p \leqslant 2$. For every $p_{0}$ such that $p_{\min }<p_{0}<$ $p \leqslant 2$, we have

$$
\begin{gather*}
\left\|T^{*}\right\|_{p} \leqslant\left(C_{r_{0}}\right)^{2 \gamma / p_{0}}\left(C_{p_{0}}^{\prime \prime}\right)^{\gamma}\left(\sum_{k \in \mathbb{Z}} a_{k}^{(1-\gamma) p / 2}\right)^{2 / p}+2 C_{p}^{\prime}  \tag{6.4}\\
\text { with } r_{0}=2 p /\left(p+2-p_{0}\right) \in\left(p_{0}, p\right) \text { and } \gamma=[1 / p-1 / 2] /\left[1 / p_{0}-1 / 2\right] .
\end{gather*}
$$

Our main interest in applications will be the maximal operator $U^{*}$, which is also bounded on $L^{p}(X, \Sigma, \mu)$ since $S^{*}$ is bounded on $L^{p}(X, \Sigma, \mu)$ according to $\left(A_{1}\right)$.

Proof. - Under the assumption $\left(A_{\mathbf{3}}\right)$, one already knows by (6.3) that $T^{*}$ is bounded on $L^{2}(X, \Sigma, \mu)$. We fix $p_{1}=p$ such that $p_{\min }<p_{1}<2$ and we try to prove that $T^{*}$ is bounded on $L^{p_{1}}(X, \Sigma, \mu)$. For doing this, it is enough to show that for every finite subfamily $\left(T_{j}\right)_{j \in J}$ of $\left(T_{j}\right)_{j \in \mathbb{Z}}$, the corresponding maximal operator

$$
f \mapsto \max _{j \in J}\left|T_{j} f\right|
$$

is $L^{p_{1}}$-bounded by a constant independent of the chosen finite subset $J \subset \mathbb{Z}$.
We thus consider a family $\left(T_{j}\right)$ that has only a finite number of nonzero terms, implying that $\left\|T^{*}\right\|_{p_{1}}<+\infty$ by Property $\left(A_{\mathbf{2}}\right)$. We choose $p_{0}$ arbitrary such that $p_{\min }<p_{0}<p_{1}$, and we introduce $r_{0}$ such that $p_{\min }<p_{0}<$ $r_{0}<p_{1}<r_{1}:=2$, defined in this way: if $\theta \in(0,1)$ is such that

$$
\begin{equation*}
\frac{1}{2}=\frac{1-\theta}{p_{0}}+\frac{\theta}{\infty} \tag{6.5a}
\end{equation*}
$$

that is to say, if $\theta=1-p_{0} / 2$, then we set

$$
\begin{equation*}
\frac{1}{r_{0}}=\frac{1-\theta}{p_{0}}+\frac{\theta}{p_{1}} \quad\left(=\frac{1}{2}+\frac{1}{p_{1}}-\frac{p_{0}}{2 p_{1}}, r_{0}=\frac{2 p_{1}}{p_{1}+2-p_{0}}\right) . \tag{6.5b}
\end{equation*}
$$

Here is the plan: by a first interpolation between $p_{0}$ and $p_{1}$, we will show that $T^{*}$ satisfies $\left(\mathbf{W}_{r_{0}}\right)$ with a constant bounded by a function of $\left\|T^{*}\right\|_{p_{1}}$. Next, we will interpolate between $\left(\mathbf{W}_{r_{0}}\right)$ and $\left(\mathbf{S}_{r_{1}}\right)=\left(\mathbf{S}_{2}\right)$ and obtain $\left(\mathbf{S}_{p_{1}}\right)$, giving a new bound for the norm $\left\|T^{*}\right\|_{p_{1}}$, whose particular form

$$
\left\|T^{*}\right\|_{p_{1}} \leqslant A\left(\left\|T^{*}\right\|_{p_{1}}+B\right)^{\beta}, \quad \text { for some } \beta \in(0,1)
$$

implies that $\left\|T^{*}\right\|_{p_{1}}$ is bounded by a constant independent of the chosen finite subfamily. This will complete the proof.

For $1 \leqslant r, s \leqslant+\infty$, let $\kappa(r, s)$ be the smallest constant such that

$$
\left\|\left(\sum_{j \in \mathbb{Z}}\left|T_{j} g_{j}\right|^{s}\right)^{1 / s}\right\|_{L^{r}} \leqslant \kappa(r, s)\left\|\left(\sum_{j \in \mathbb{Z}}\left|g_{j}\right|^{s}\right)^{1 / s}\right\|_{L^{r}}
$$

for every sequence $\left(g_{j}\right)_{j \in \mathbb{Z}}$ in $L^{r}(X, \Sigma, \mu)$.

- One sees that $\kappa\left(p_{0}, p_{0}\right) \leqslant C_{p_{0}}^{\prime \prime}$, by $\left(A_{\mathbf{2}}\right)$ and the simple sum-integral inversion

$$
\begin{aligned}
\left\|\left(\sum_{j \in \mathbb{Z}}\left|T_{j} g_{j}\right|^{p_{0}}\right)^{1 / p_{0}}\right\|_{L^{p_{0}}}^{p_{0}} & =\sum_{j \in \mathbb{Z}}\left\|T_{j} g_{j}\right\|_{L^{p_{0}}}^{p_{0}} \\
& \leqslant\left(C_{p_{0}}^{\prime \prime}\right)^{p_{0}}\left\|\left(\sum_{j \in \mathbb{Z}}\left|g_{j}\right|^{p_{0}}\right)^{1 / p_{0}}\right\|_{L^{p_{0}}}^{p_{0}} .
\end{aligned}
$$

- One has also $\kappa\left(p_{1},+\infty\right) \leqslant\left\|T^{*}\right\|_{p_{1}}+2 C_{p_{1}}^{\prime}$. Indeed, when $\left(W_{j}\right)_{j \in \mathbb{Z}}$ is a family of positive operators and $g=\sup _{j \in \mathbb{Z}}\left|g_{j}\right|$, one has

$$
\left|W_{j} g_{j}\right| \leqslant W_{j}\left|g_{j}\right| \leqslant W_{j} g, \quad \sup _{j \in \mathbb{Z}}\left|W_{j} g_{j}\right| \leqslant \sup _{j \in \mathbb{Z}} W_{j} g
$$

Because $S_{j, v}$ is positive, we have $\sup _{j \in \mathbb{Z}}\left|S_{j, v} g_{j}\right| \leqslant \sup _{j \in \mathbb{Z}} S_{j, v} g$ for every $v \in V$, and letting $S_{j} g_{j}=\sup _{v \in V}\left|S_{j, v} g_{j}\right|$ we see according to $\left(A_{1}\right)$ that

$$
\sup _{j \in \mathbb{Z}} S_{j} g_{j} \leqslant S^{*} g, \quad\left\|\sup _{j \in \mathbb{Z}} S_{j} g_{j}\right\|_{L^{p_{1}}} \leqslant\left\|S^{*} g\right\|_{L^{p_{1}}} \leqslant C_{p_{1}}^{\prime}\left\|\sup _{j \in \mathbb{Z}}\left|g_{j}\right|\right\|_{L^{p_{1}}}
$$

Since $U_{j, v}=T_{j, v}+S_{j, v}$ is positive, we obtain also for $U_{j} f=$ $\sup _{v \in V}\left|U_{j, v} f\right|$ that

$$
\begin{aligned}
\left\|\sup _{j \in \mathbb{Z}}\left|U_{j} g_{j}\right|\right\|_{L^{p_{1}}} \leqslant\left\|U^{*} g\right\|_{L^{p_{1}}} & \leqslant\left\|T^{*} g\right\|_{L^{p_{1}}}+\left\|S^{*} g\right\|_{L^{p_{1}}} \\
& \leqslant\left(\left\|T^{*}\right\|_{p_{1}}+C_{p_{1}}^{\prime}\right)\|g\|_{L^{p_{1}}}
\end{aligned}
$$

and finally $\left\|\sup _{j \in \mathbb{Z}}\left|T_{j} g_{j}\right|\right\|_{L^{p_{1}}} \leqslant\left(\left\|T^{*}\right\|_{p_{1}}+2 C_{p_{1}}^{\prime}\right)\left\|\sup _{j \in \mathbb{Z}}\left|g_{j}\right|\right\|_{L^{p_{1}}}$, which proves the inequality $\kappa\left(p_{1},+\infty\right) \leqslant\left\|T^{*}\right\|_{p_{1}}+2 C_{p_{1}}^{\prime}$.

We apply complex interpolation between spaces $L^{p}\left(\ell^{q}\right)$ [7, Chap. 5, Th. 5.1.2], namely between the spaces $L^{p_{0}}\left(\ell^{p_{0}}\right)$ and $L^{p_{1}}\left(\ell^{\infty}\right)$, which gives the space $L^{r_{0}}\left(\ell^{2}\right)$ for the chosen value $\theta$ of the interpolation parameter, by (6.5a) and (6.5b). We already explained that the case where $T_{j}$ is not linear can also be covered by complex interpolation. It follows from (3.26) that

$$
\kappa\left(r_{0}, 2\right) \leqslant \kappa\left(p_{0}, p_{0}\right)^{1-\theta} \kappa\left(p_{1},+\infty\right)^{\theta} \leqslant\left(C_{p_{0}}^{\prime \prime}\right)^{1-\theta}\left(\left\|T^{*}\right\|_{p_{1}}+2 C_{p_{1}}^{\prime}\right)^{\theta}
$$

Remark (in passing). - It is exactly in this manner that Stein [73, Chap. VI, Th. 8, p. 103] shows the inequality (6.6) on the square function $\left(\sum_{n}\left|E_{n} f_{n}\right|^{2}\right)^{1 / 2}$ of a sequence $\left(E_{n}\right)$ of conditional expectations with respect to an increasing sequence of $\sigma$-fields, stating that

$$
\begin{equation*}
\left\|\left(\sum_{n}\left|E_{n} f_{n}\right|^{2}\right)^{1 / 2}\right\|_{q} \leqslant \kappa_{q}\left\|\left(\sum_{n}\left|f_{n}\right|^{2}\right)^{1 / 2}\right\|_{q}, \quad 1<q<+\infty \tag{6.6}
\end{equation*}
$$

When $1<q<2$, the proof applies inversion for a pair ( $q_{0}, q_{0}$ ), and Doob's maximal theorem for a pair $\left(q_{1},+\infty\right)$ with $q_{0}<q<q_{1}$ and $q\left(q_{1}-q_{0}\right)=$ $2\left(q_{1}-q\right)$.
Thus, with $g_{j}=Q_{j+k} f$ for a fixed $k \in \mathbb{Z}$, one has

$$
\begin{aligned}
\left\|\sup _{j \in \mathbb{Z}}\left|T_{j} Q_{j+k} f\right|\right\|_{L^{r_{0}}} & \leqslant\left\|\left(\sum_{j \in \mathbb{Z}}\left|T_{j} Q_{j+k} f\right|^{2}\right)^{1 / 2}\right\|_{L^{r_{0}}} \\
& \leqslant \kappa\left(r_{0}, 2\right)\left\|\left(\sum_{j \in \mathbb{Z}}\left|Q_{j+k} f\right|^{2}\right)^{1 / 2}\right\|_{L^{r_{0}}} \\
& =\kappa\left(r_{0}, 2\right)\left\|\left(\sum_{j \in \mathbb{Z}}\left|Q_{j} f\right|^{2}\right)^{1 / 2}\right\|_{L^{r_{0}}} \leqslant C_{r_{0}} \kappa\left(r_{0}, 2\right)\|f\|_{L^{r_{0}}} .
\end{aligned}
$$

We have proved the property $\left(\mathbf{W}_{r_{0}}\right)$, since we got that

$$
\forall f \in L^{r_{0}}, \forall k \in \mathbb{Z}, \quad\left\|\sup _{j \in \mathbb{Z}}\left|T_{j} Q_{j+k} f\right|\right\|_{L^{r_{0}}} \leqslant C_{r_{0}} \kappa\left(r_{0}, 2\right)\|f\|_{L^{r_{0}}}
$$

If for a certain $\rho \in(0,1)$, we write

$$
\frac{1}{p_{1}}=\frac{1-\rho}{r_{0}}+\frac{\rho}{r_{1}}=\frac{1-\rho}{r_{0}}+\frac{\rho}{2} \quad\left(\rho=\frac{p_{1}-p_{0}}{2-p_{0}}\right)
$$

we get $\left(\mathbf{S}_{p_{1}}\right)$ by interpolating between $\left(\mathbf{W}_{r_{0}}\right)$ and $\left(\mathbf{S}_{2}\right)=\left(\mathbf{S}_{r_{1}}\right)$, obtaining thus

$$
\forall k \in \mathbb{Z}, \quad\left\|\sup _{j \in \mathbb{Z}}\left|T_{j} Q_{j+k} f\right|\right\|_{L^{p_{1}}} \leqslant\left(C_{r_{0}} \kappa\left(r_{0}, 2\right)\right)^{1-\rho} a_{k}^{\rho}\|f\|_{L^{p_{1}}}
$$

By (6.3), it follows that

$$
\left\|\sup _{j \in \mathbb{Z}}\left|T_{j} f\right|\right\|_{L^{p_{1}}} \leqslant\left(C_{r_{0}} \kappa\left(r_{0}, 2\right)\right)^{1-\rho}\left(\sum_{k \in \mathbb{Z}} a_{k}^{\rho}\right)\|f\|_{L^{p_{1}}}
$$

One has finally an implicit inequality about $\left\|T^{*}\right\|_{p_{1}}$, namely

$$
\begin{aligned}
\left\|T^{*}\right\|_{p_{1}} & \leqslant\left[C_{r_{0}} \kappa\left(r_{0}, 2\right)\right]^{1-\rho}\left(\sum_{k \in \mathbb{Z}} a_{k}^{\rho}\right) \\
& \leqslant\left[C_{r_{0}}\left(C_{p_{0}}^{\prime \prime}\right)^{1-\theta}\left(\left\|T^{*}\right\|_{p_{1}}+2 C_{p_{1}}^{\prime}\right)^{\theta}\right]^{1-\rho}\left(\sum_{k \in \mathbb{Z}} a_{k}^{\rho}\right) \\
& =\left(C_{r_{0}}\left(C_{p_{0}}^{\prime \prime}\right)^{1-\theta}\right)^{1-\rho}\left(\sum_{k \in \mathbb{Z}} a_{k}^{\rho}\right)\left(\left\|T^{*}\right\|_{p_{1}}+2 C_{p_{1}}^{\prime}\right)^{\theta(1-\rho)},
\end{aligned}
$$

implying that $\left\|T^{*}\right\|_{p_{1}}$ is bounded by a constant depending only upon $C_{r_{0}}$, $C_{p_{1}}^{\prime}, C_{p_{0}}^{\prime \prime}$ and the $a_{k} \mathrm{~s}$. Indeed, suppose that $C \geqslant 0$ satisfies $C \leqslant A(C+B)^{\beta}$, where $A, B>0$ and $0<\beta<1$. We write

$$
C \leqslant\left(A^{1 /(1-\beta)}\right)^{1-\beta}(C+B)^{\beta} \leqslant(1-\beta) A^{1 /(1-\beta)}+\beta(C+B)
$$

yielding

$$
C \leqslant A^{1 /(1-\beta)}+\frac{\beta}{1-\beta} B
$$

This bound is essentially correct when $B$ is small, and we shall use it below with $A=\left(C_{r_{0}}\left(C_{p_{0}}^{\prime \prime}\right)^{1-\theta}\right)^{1-\rho}\left(\sum_{k \in \mathbb{Z}} a_{k}^{\rho}\right), B=2 C_{p_{1}}^{\prime}$ and $\beta=\theta(1-\rho)$.

However, when $B \geqslant A^{1 /(1-\beta)}$, a better bound $(1-\beta)^{-1} A B^{\beta}$ is available.
In this case, $A \leqslant B^{1-\beta}$, thus $C \leqslant B^{1-\beta}(C+B)^{\beta} \leqslant B+\beta C$, hence

$$
C \leqslant A\left(\frac{B}{1-\beta}+B\right)^{\beta}=\left(\frac{2-\beta}{1-\beta}\right)^{\beta} A B^{\beta} \leqslant \frac{1}{1-\beta} A B^{\beta},
$$

because $(2-\beta)^{\beta}(1-\beta)^{1-\beta} \leqslant \beta(2-\beta)+(1-\beta)^{2}=1$.
Recall that $\rho=\left(p_{1}-p_{0}\right) /\left(2-p_{0}\right)$, so $\beta=\theta(1-\rho)=1-p_{1} / 2<1$. We find an explicit bound for $\left\|T^{*}\right\|_{p_{1}}$, independent of the finite subfamily $\left(T_{j}\right)_{j \in J}$ of $\left(T_{j}\right)$ s that was chosen at the beginning, of the form

$$
\begin{aligned}
\left\|T^{*}\right\|_{p_{1}} & \leqslant\left(C_{r_{0}}\left(C_{p_{0}}^{\prime \prime}\right)^{1-\theta}\right)^{2(1-\rho) / p_{1}}\left(\sum_{k \in \mathbb{Z}} a_{k}^{\rho}\right)^{2 / p_{1}}+\frac{2-p_{1}}{p_{1}} 2 C_{p_{1}}^{\prime} \\
& \leqslant\left(C_{r_{0}}\right)^{2 \gamma / p_{0}}\left(C_{p_{0}}^{\prime \prime}\right)^{\gamma}\left(\sum_{k \in \mathbb{Z}} a_{k}^{\rho}\right)^{2 / p_{1}}+2 C_{p_{1}}^{\prime}
\end{aligned}
$$

with $\gamma=\left[1 / p_{1}-1 / 2\right] /\left[1 / p_{0}-1 / 2\right]$. Observe that $\rho=\left[p_{1} /\left(2 p_{0}\right)-1 / 2\right] /\left[1 / p_{0}-\right.$ $1 / 2]=(1-\gamma) p_{1} / 2$. We get in particular a bound of $C_{p_{1}}^{\prime \prime}$ by a power $\gamma<1$ of $C_{p_{0}}^{\prime \prime}$. There is no miracle: this power $\gamma$ is the one corresponding to interpolation between $C_{p_{0}}^{\prime \prime}$ and the value $C_{2}^{\prime \prime}$ hidden in the assumption $\left(A_{3}\right)$.

### 6.2. Fractional derivatives

If a function $h$ is given in the Schwartz space $\mathcal{S}(\mathbb{R})$, one can express it as Fourier transform of another function $k \in \mathcal{S}(\mathbb{R})$ and write

$$
\forall t \in \mathbb{R}, \quad h(t)=\int_{\mathbb{R}} k(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s
$$

One has then an expression for the derivatives of $h$ by means of (unbounded) multipliers. For every integer $j \geqslant 1$ and every $t \in \mathbb{R}$, one sees that

$$
(-1)^{j} h^{(j)}(t)=\int_{\mathbb{R}}(2 \mathrm{i} \pi s)^{j} k(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s
$$

It is tempting to extend the notion of derivative, from the integer case $j \in \mathbb{N}$ to every complex value $z$ such that $\operatorname{Re} z>-1$, by setting

$$
\begin{equation*}
\forall t \in \mathbb{R}, \quad\left(D^{z} h\right)(t)=\int_{\mathbb{R}}(2 \mathrm{i} \pi s)^{z} k(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s \tag{6.7}
\end{equation*}
$$

Note that $D^{1} h=-h^{\prime}$ with this definition. We define complex powers by

$$
(2 \mathrm{i} \pi s)^{z}=\mathrm{e}^{z \ln (2 \mathrm{i} \pi s)}=\mathrm{e}^{z(\ln (2 \pi|s|)+\mathrm{i} \operatorname{Arg}(2 \mathrm{i} \pi s))}=|2 \pi s|^{z} \mathrm{e}^{\mathrm{i} \pi z \operatorname{sign}(s) / 2}
$$

and we have that $(\lambda \mathrm{i} s)^{z}=\lambda^{z}(\mathrm{i} s)^{z}$ when $\lambda>0$. If we dilate the function $h$ to $h_{[\lambda]}$, with $\lambda>0$ as in (2.7), we know that $h_{[\lambda]}=\mathcal{F}\left(k_{(\lambda)}\right)$, therefore

$$
\begin{aligned}
\left(D^{z} h_{[\lambda]}\right)(t) & =\int_{\mathbb{R}}(2 \mathrm{i} \pi s)^{z} \lambda^{-1} k\left(\lambda^{-1} s\right) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s \\
& =\lambda^{z} \int_{\mathbb{R}}(2 \mathrm{i} \pi u)^{z} k(u) \mathrm{e}^{-2 \mathrm{i} \pi u \lambda t} \mathrm{~d} u
\end{aligned}
$$

This means that

$$
\begin{equation*}
D^{z}\left(h_{[\lambda]}\right)=\lambda^{z}\left(D^{z} h\right)_{[\lambda]}, \quad \text { or } \quad D_{t}^{z} h(\lambda t)=\lambda^{z}\left(D^{z} h\right)(\lambda t) \tag{6.8}
\end{equation*}
$$

where we use the notation $D_{t}^{z} h(\lambda t)$ when the function of $t$ does not have an explicit name, as in $t \mapsto h(\lambda t)$. For a specific value, we shall write for example $\left.D_{t}^{z} h(\lambda t)\right|_{t=1}$.

If we would like to extend $D^{z}$ to $h=\mathbf{1}$, we might consider the function $\mathbf{1}$ as the limit of $h_{[\lambda]}$ when $h(0)=1$ and $\lambda \searrow 0$. Then (6.8) suggests that $D^{z} \mathbf{1}=0$ when $\operatorname{Re} z>0$, and that $D^{z} \mathbf{1}$ is undefined if $\operatorname{Re} z<0$.

When $z$ is not a nonnegative integer, the operator $D^{z}$ is not local. We will see later however that $\left(D^{z} h\right)\left(t_{0}\right)$ depends only on the values of $h$ on $\left[t_{0},+\infty\right)$. This could be checked right now by arguments involving holomorphic functions.

When $-1<\operatorname{Re} z<0$, the differentiation $D^{z}$ is in fact a fractional integration. We shall see below that $\left(D^{z} h\right)(t)=\left(I^{-z} h\right)(t)$, where $I^{w}$ is given for $\operatorname{Re} w>0$ by

$$
\begin{equation*}
\left(I^{w} h\right)(t)=\frac{1}{\Gamma(w)} \int_{t}^{+\infty}(u-t)^{w-1} h(u) \mathrm{d} u \tag{6.9}
\end{equation*}
$$

The next lemma provides the tool that relates the definitions (6.7) and (6.9).
Lemma 6.7. - Let $\zeta$ be a complex number such that $\operatorname{Re} \zeta<0$ and let $\varepsilon>0$. The inverse Fourier transform of the function $t \mapsto$ $\Gamma(-\zeta)^{-1} \mathbf{1}_{(-\infty, 0)}(t)(-t)^{-\zeta-1} \mathrm{e}^{\varepsilon t}$ is equal to $s \mapsto(\varepsilon+2 \mathrm{i} \pi s)^{\zeta}$, namely

$$
\frac{1}{\Gamma(-\zeta)} \int_{\mathbb{R}} \mathbf{1}_{(-\infty, 0)}(t)(-t)^{-\zeta-1} \mathrm{e}^{\varepsilon t} \mathrm{e}^{2 \mathrm{i} \pi s t} \mathrm{~d} t=(\varepsilon+2 \mathrm{i} \pi s)^{\zeta}, \quad s \in \mathbb{R}
$$

Proof. - By a contour integral of $(-z)^{-\zeta-1} \mathrm{e}^{z}$, running along the negative real half-line and along the half-line $H_{s}=\{(\varepsilon+2 \mathrm{i} \pi s) t \in \mathbb{C}: t<0\}$, we obtain

$$
\Gamma(-\zeta)=\int_{-\infty}^{0}(-t)^{-\zeta-1} \mathrm{e}^{t} \mathrm{~d} t=(\varepsilon+2 \mathrm{i} \pi s)^{-\zeta} \int_{-\infty}^{0}(-t)^{-\zeta-1} \mathrm{e}^{(\varepsilon+2 \mathrm{i} \pi s) t} \mathrm{~d} t
$$

giving the announced result.
Integrating (6.9) by parts, we see that

$$
\left(I^{w} h\right)(t)=-\frac{1}{\Gamma(w+1)} \int_{t}^{+\infty}(u-t)^{w} h^{\prime}(u) \mathrm{d} u
$$

This new formula makes sense for $\operatorname{Re} w>-1$ and could be used for defining the fractional derivative $D^{z}$ if $z=-w$ and $\operatorname{Re} w \in(-1,0)$, by setting for $t$ real

$$
\begin{equation*}
\left(D^{z} h\right)(t)=-\frac{1}{\Gamma(1-z)} \int_{t}^{+\infty}(u-t)^{-z} h^{\prime}(u) \mathrm{d} u \tag{6.10}
\end{equation*}
$$

This is proved in Lemma 6.8. It is coherent with the fact that $D^{\alpha}$, for $0<$ $\alpha<1$, can be considered as the antiderivative of order $1-\alpha$ of the derivative $D^{1} h=-h^{\prime}$,

$$
D^{\alpha} h=D^{\alpha-1} D^{1} h=-D^{\alpha-1} h^{\prime}=-I^{1-\alpha} h^{\prime}
$$

The operation $D^{z}$ is not symmetric on $\mathbb{R}$; this is obvious from the formulas for $I^{w}$. The choice that was done of $(2 \mathrm{i} \pi s)^{z}$ instead of $(-2 \mathrm{i} \pi s)^{z}$ in (6.7) induces the direction in which the fractional antiderivative is computed. This direction, to $+\infty$, is well adapted to the "radial" Carbery's method introduced in [20].

Lemma 6.8. - Let $\alpha \in(0,1), t_{0} \in \mathbb{R}$ be given and let $k$ be a function on $\mathbb{R}$ such that $\left(1+|s|^{\alpha}\right) k(s)$ is integrable on the real line. Assume that $h=\widehat{k}$
is Lipschitz with $\left|h^{\prime}(t)\right| \leqslant \kappa_{1}(1+|t|)^{-1}$ for almost every $t \geqslant t_{0}$. Then, for every $t>t_{0}$ and $z$ such that $\operatorname{Re} z=\alpha$, we have

$$
-\frac{1}{\Gamma(1-z)} \int_{t}^{+\infty}(u-t)^{-z} h^{\prime}(u) \mathrm{d} u=\int_{\mathbb{R}}(2 \mathrm{i} \pi s)^{z} k(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s
$$

Proof. - Let $\eta$ be a nonnegative $C^{\infty}$ function on $\mathbb{R}$, with integral 1 and with compact support in $[-1,1]$. Consider $\varepsilon \in(0,1)$ and

$$
k_{\varepsilon}(s)=k(s)\left(\eta^{\vee}\right)_{[\varepsilon]}(s)=k(s) \eta^{\vee}(\varepsilon s), \quad s \in \mathbb{R}
$$

Then $\eta^{\vee} \in \mathcal{S}(\mathbb{R}), s k_{\varepsilon}(s)$ is integrable and $h_{\varepsilon}:=\widehat{k_{\varepsilon}}=h * \eta_{(\varepsilon)}$ is $C^{1}$. We can write

$$
\begin{equation*}
-h_{\varepsilon}^{\prime}(t)=\int_{\mathbb{R}} 2 \mathrm{i} \pi s k_{\varepsilon}(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s, \quad t \in \mathbb{R} \tag{6.11}
\end{equation*}
$$

Since $h$ is Lipschitz, we also know that $h_{\varepsilon}^{\prime}=h^{\prime} * \eta_{(\varepsilon)}$. Fix $t \geqslant t_{0}+\varepsilon$. When $|\tau| \leqslant 1$ and $u \geqslant t$, we have $u-\varepsilon \tau \geqslant t_{0}, 1+|u| \leqslant 1+\varepsilon|\tau|+|u-\varepsilon \tau| \leqslant 2+2|u-\varepsilon \tau|$, so

$$
\begin{equation*}
\left|h_{\varepsilon}^{\prime}(u)\right|=\left|\int_{-1}^{1} h^{\prime}(u-\varepsilon \tau) \eta(\tau) \mathrm{d} \tau\right| \leqslant \kappa_{1} \int_{-1}^{1} \frac{\eta(\tau)}{1+|u-\varepsilon \tau|} \mathrm{d} \tau \leqslant \frac{2 \kappa_{1}}{1+|u|} \tag{6.12}
\end{equation*}
$$

Applying (6.11) and $\left|(u-t)^{-z}\right|=(u-t)^{-\alpha}$, Fubini's theorem and the inverse Fourier transform of $v \mapsto\left[(-v)_{+}\right]^{-z} \mathrm{e}^{\varepsilon v}$ given by Lemma 6.7 with $\zeta=z-1$, we get

$$
\begin{aligned}
& -\frac{1}{\Gamma(1-z)} \int_{t}^{+\infty}(u-t)^{-z} \mathrm{e}^{\varepsilon(t-u)} h_{\varepsilon}^{\prime}(u) \mathrm{d} u \\
& =\frac{1}{\Gamma(1-z)} \int_{t}^{+\infty}(u-t)^{-z} \mathrm{e}^{\varepsilon(t-u)}\left(\int_{\mathbb{R}} 2 \mathrm{i} \pi s k_{\varepsilon}(s) \mathrm{e}^{-2 \mathrm{i} \pi s u} \mathrm{~d} s\right) \mathrm{d} u \\
& =\frac{1}{\Gamma(1-z)} \iint \mathbf{1}_{\{t-u<0\}}(u-t)^{-z} \mathrm{e}^{\varepsilon(t-u)} 2 \mathrm{i} \pi s k_{\varepsilon}(s) \mathrm{e}^{2 \mathrm{i} \pi s(t-u)} \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s \mathrm{~d} u \\
& =\int_{\mathbb{R}}(\varepsilon+2 \mathrm{i} \pi s)^{z-1}(2 \mathrm{i} \pi s) k_{\varepsilon}(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s
\end{aligned}
$$

Letting $\varepsilon$ tend to 0 , by a double application of Lebesgue's dominated convergence, using (6.12) and since $h_{\varepsilon}^{\prime}(u) \rightarrow h^{\prime}(u)$ at every Lebesgue point $u$ of $h^{\prime}$, we obtain

$$
-\frac{1}{\Gamma(1-z)} \int_{t}^{+\infty}(u-t)^{-z} h^{\prime}(u) \mathrm{d} u=\int_{\mathbb{R}}(2 \mathrm{i} \pi s)^{z} k(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s
$$

It is quite comforting to have two possible ways of defining $D^{z} h$. However, we will have to handle cases where the Fourier transform $h(t)$ is well controlled, but where the estimates on $k(s)$ are not so good. We shall therefore concentrate on the integral definition (6.10) of $D^{z} h$. We have to check
that the properties obtained with the first definition remain true when only the second applies.

When $\alpha \in(0,1)$ tends to 1 , one has $\Gamma(1-\alpha) \simeq(1-\alpha)^{-1}$ and for $\varepsilon>0$ we get

$$
\begin{aligned}
& -\frac{1}{\Gamma(1-\alpha)} \int_{t+\varepsilon}^{+\infty}(u-t)^{-\alpha} h^{\prime}(u) \mathrm{d} u \rightarrow 0, \\
& \\
& \quad(1-\alpha) \int_{t}^{t+\varepsilon}(u-t)^{-\alpha} \mathrm{d} u=\varepsilon^{1-\alpha} \rightarrow 1 .
\end{aligned}
$$

We recover the fact that $\left(D^{1} h\right)(t)=-h^{\prime}(t)$, already known by Fourier.
Let us mention the case of $h(t)=\mathrm{e}^{-\lambda|t|}$, the Fourier transform of a Cauchy kernel. When $t \geqslant 0$ and $0<\operatorname{Re} z<1$, we have

$$
\begin{align*}
D_{t}^{z} \mathrm{e}^{-\lambda|t|}= & \frac{1}{\Gamma(1-z)} \mathrm{e}^{-\lambda t} \int_{t}^{+\infty}(u-t)^{-z} \lambda \mathrm{e}^{-\lambda(u-t)} \mathrm{d} u  \tag{6.13}\\
& =\lambda^{z} \mathrm{e}^{-\lambda|t|}
\end{align*}
$$

The dilation relation (6.8) follows from a simple change of variable similar to the one in the line above, and is left to the reader.

We have introduced in (5.5) the right maximal function $h_{r}^{*}$ of $h$. Notice that for $h$ Lipschitz on $\left(t_{0},+\infty\right)$ and for every $t \geqslant t_{0}, \delta>0$, we have

$$
\begin{equation*}
|h(t+\delta)| \leqslant|h(t)|+\int_{t}^{t+\delta}\left|h^{\prime}(u)\right| \mathrm{d} u \leqslant h_{r}^{*}(t)+\delta\left(h^{\prime}\right)_{r}^{*}(t) \tag{6.14}
\end{equation*}
$$

Lemma 6.9. - Let $h$ be Lipschitz on $\left(t_{0},+\infty\right), \alpha \in(0,1)$ and $h(t)=$ $o\left(t^{\alpha}\right)$ at $+\infty$. Let $h_{0}=h_{r}^{*}$ be the right maximal function of $h$ and $h_{1}=\left(h^{\prime}\right)_{r}^{*}$ that of $h^{\prime}$. Then

$$
\left|\left(D^{\alpha} h\right)(t)\right| \leqslant 6 h_{0}(t)^{1-\alpha} h_{1}(t)^{\alpha}, \quad t \geqslant t_{0} .
$$

If $w$ is complex and $\operatorname{Re} w=\alpha$, then for every $t \geqslant t_{0}$ we have

$$
\left|\left(D^{w} h\right)(t)\right| \leqslant \frac{2}{\alpha(1-\alpha)} \frac{(1+|w|)^{1-\alpha}}{|\Gamma(1-w)|} h_{0}(t)^{1-\alpha} h_{1}(t)^{\alpha} .
$$

Proof. - For $t \geqslant t_{0}$ and $\delta>0$, we express $E_{\alpha}:=-\Gamma(1-\alpha)\left(D^{\alpha} h\right)(t)$ as

$$
\int_{t}^{t+\delta}(u-t)^{-\alpha} h^{\prime}(u) \mathrm{d} u+\int_{t+\delta}^{+\infty}(u-t)^{-\alpha} h^{\prime}(u) \mathrm{d} u
$$

Applying (5.6) and integration by parts, we bound each of the two pieces

$$
\begin{aligned}
\left|E_{\alpha}\right| \leqslant & \frac{\delta^{1-\alpha}}{1-\alpha} h_{1}(t)+\left|\left[(u-t)^{-\alpha} h(u)\right]_{t+\delta}^{+\infty}\right|+\alpha\left|\int_{t+\delta}^{+\infty}(u-t)^{-\alpha-1} h(u) \mathrm{d} u\right| \\
\leqslant & \frac{\delta^{1-\alpha}}{1-\alpha} h_{1}(t)+\delta^{-\alpha}|h(t+\delta)| \\
& \quad+\alpha\left(\int_{t}^{t+\delta} \delta^{-\alpha-1}|h(u)| \mathrm{d} u+\int_{t+\delta}^{+\infty}(u-t)^{-\alpha-1}|h(u)| \mathrm{d} u\right)
\end{aligned}
$$

By (6.14), by (5.6) for the non-decreasing function $\psi$ defined by $\psi(u)=$ $\delta^{-\alpha-1}$ when $u \in[t, t+\delta]$ and $\psi(u)=(u-t)^{-\alpha-1}$ for $u \geqslant t+\delta$, we obtain

$$
\begin{aligned}
\left|E_{\alpha}\right| & \leqslant \frac{\delta^{1-\alpha}}{1-\alpha} h_{1}(t)+\delta^{-\alpha}\left(h_{0}(t)+\delta h_{1}(t)\right)+(1+\alpha) \delta^{-\alpha} h_{0}(t) \\
& =\frac{2-\alpha}{1-\alpha} \delta^{1-\alpha} h_{1}(t)+(2+\alpha) \delta^{-\alpha} h_{0}(t) \\
& \leqslant \frac{2}{1-\alpha} \delta^{1-\alpha} h_{1}(t)+3 \delta^{-\alpha} h_{0}(t) .
\end{aligned}
$$

We choose $\delta=\delta_{0}=h_{0}(t) / h_{1}(t)$ and get that

$$
\left|E_{\alpha}\right| \leqslant\left(\frac{2}{1-\alpha}+3\right) h_{0}(t)^{1-\alpha} h_{1}(t)^{\alpha} .
$$

Recalling $\Gamma(1-\alpha) \geqslant 1$ and the minimal value $\Gamma\left(x_{\Gamma}\right)>0.88$ in (3.7) we have

$$
\begin{aligned}
\left|\left(D^{\alpha} h\right)(t)\right| & \leqslant\left(\frac{2}{\Gamma(2-\alpha)}+\frac{3}{\Gamma(1-\alpha)}\right) h_{0}(t)^{1-\alpha} h_{1}(t)^{\alpha} \\
& \leqslant\left(\frac{2}{\Gamma\left(x_{\Gamma}\right)}+3\right) h_{0}(t)^{1-\alpha} h_{1}(t)^{\alpha} \leqslant 6 h_{0}(t)^{1-\alpha} h_{1}(t)^{\alpha}
\end{aligned}
$$

When $w$ is complex and $\operatorname{Re} w=\alpha$, we use $\left|(u-t)^{-w}\right|=(u-t)^{-\alpha}$, the same integration by parts, $\left|(u-t)^{-w-1}\right|=(u-t)^{-\alpha-1}$ and we get

$$
\begin{aligned}
\left|E_{w}\right| & \leqslant \frac{\delta^{1-\alpha}}{1-\alpha} h_{1}(t)+\delta^{-\alpha}\left(h_{0}(t)+\delta h_{1}(t)\right)+|w|\left(1+\frac{1}{\alpha}\right) \delta^{-\alpha} h_{0}(t) \\
& \leqslant \frac{2 \delta^{1-\alpha}}{1-\alpha} h_{1}(t)+\frac{2}{\alpha}(1+|w|) \delta^{-\alpha} h_{0}(t)
\end{aligned}
$$

Choosing $\delta=(1+|w|) h_{0}(t) / h_{1}(t)$ we obtain the announced result.
In what follows, we shall consider the following assumptions on a function $h$ :

$$
\left\{\begin{array}{l}
h \text { is Lipschitz on }\left[t_{0},+\infty\right)  \tag{6.15}\\
|h(t)| \leqslant \kappa_{0}(1+|t|)^{-1} \text { for } t \geqslant t_{0} \\
\left|h^{\prime}(t)\right| \leqslant \kappa_{1}(1+|t|)^{-1} \text { for almost every } t \geqslant t_{0}
\end{array}\right.
$$

Corollary 6.10. - Suppose that the function $h$ defined on $\left(t_{0},+\infty\right)$, $t_{0} \geqslant 0$, satisfies (6.15). Then for every $\alpha \in(0,1)$, we have

$$
\left|\left(D^{\alpha} h\right)(t)\right| \leqslant 6 \frac{\kappa_{0}^{1-\alpha} \kappa_{1}^{\alpha}}{1+|t|}, \quad t \geqslant t_{0}
$$

Proof. - The two upper bounds in (6.15) are decreasing functions of $t \in\left[t_{0},+\infty\right)$, hence they also bound $h_{r}^{*}$ or $\left(h^{\prime}\right)_{r}^{*}$. We conclude by applying Lemma 6.9.

Assuming that $h$ has enough derivatives and continuing integrations by parts, starting from (6.10), we get successive formulas for $D^{z} h$ for each integer $j>0$, which make sense when $\operatorname{Re} z<j$. Let $z=j-1+w$, with $j \geqslant 1$ and $\operatorname{Re} w \in(0,1)$. We obtain that

$$
\left(D^{z} h\right)(t)=(-1)^{j-1}\left(D^{w} h^{(j-1)}\right)(t)=\frac{(-1)^{j}}{\Gamma(1-w)} \int_{t}^{+\infty}(u-t)^{-w} h^{(j)}(u) \mathrm{d} u,
$$

and for every $z \in \mathbb{C}$ such that $\operatorname{Re} z<j$, we have

$$
\begin{equation*}
\left(D^{z} h\right)(t)=\frac{(-1)^{j}}{\Gamma(j-z)} \int_{t}^{+\infty}(u-t)^{-z+j-1} h^{(j)}(u) \mathrm{d} u \tag{6.16}
\end{equation*}
$$

By gluing the successive definitions, we define entire functions of $z$ for every fixed $t$ and $h \in \mathcal{S}(\mathbb{R})$. By the principle of analytic continuation, we conclude that the integral formula for $D^{z} h$ coincides when $\operatorname{Re} z>-1$ with the one obtained by Fourier transform (a fact that we have checked in Lemma 6.8 when $0<\operatorname{Re} z<1$ ).

Lemma 6.11. - Let $\alpha$ be in $(0,1)$. Suppose that the function $h$ satisfies the assumptions (6.15) on $\left[t_{0},+\infty\right), t_{0} \geqslant 0$, and define $D^{\alpha} h$ by (6.10). We have that

$$
\left(I^{\alpha} D^{\alpha} h\right)(t)=t, \quad t \geqslant t_{0}
$$

Proof. - We first assume in addition that

$$
\int_{t_{0}}^{+\infty}\left|h^{\prime}(u)\right| \mathrm{d} u<+\infty, \quad \text { thus } \quad h(t)=-\int_{t}^{+\infty} h^{\prime}(u) \mathrm{d} u
$$

for every $t \geqslant t_{0}$ since $h$ is Lipschitz. For $u \geqslant t_{0}$, accepting possibly infinite integrals of nonnegative measurable functions, set

$$
G(u)=\frac{1}{\Gamma(1-\alpha)} \int_{u}^{+\infty}(v-u)^{-\alpha}\left|h^{\prime}(v)\right| \mathrm{d} v
$$

When $h$ is decreasing on $\left(t_{0},+\infty\right)$, the function $G$ is equal to $D^{\alpha} h$, and $\left|D^{\alpha} h\right| \leqslant G$ in general. Then, consider $F$, equal to $I^{\alpha} G$ in good cases, defined
for $t \geqslant t_{0}$ by

$$
\begin{aligned}
F(t) & :=\frac{1}{\Gamma(\alpha)} \int_{t}^{+\infty}(u-t)^{\alpha-1} G(u) \mathrm{d} u \\
& =\frac{1}{\Gamma(\alpha) \Gamma(1-\alpha)} \int_{t}^{+\infty}(u-t)^{\alpha-1} \int_{u}^{+\infty}(v-u)^{-\alpha}\left|h^{\prime}(v)\right| \mathrm{d} v \mathrm{~d} u \\
& =\frac{1}{\Gamma(\alpha) \Gamma(1-\alpha)} \int_{t}^{+\infty}\left(\int \mathbf{1}_{t \leqslant u \leqslant v}(u-t)^{\alpha-1}(v-u)^{-\alpha} \mathrm{d} u\right)\left|h^{\prime}(v)\right| \mathrm{d} v
\end{aligned}
$$

Setting $u=t+y(v-t)$, one gets with $\gamma_{\alpha}=\Gamma(\alpha) \Gamma(1-\alpha)$ that

$$
\begin{aligned}
F(t) & =\gamma_{\alpha}^{-1}\left(\int_{0}^{1} y^{\alpha-1}(1-y)^{-\alpha} \mathrm{d} y\right) \int_{t}^{+\infty}\left|h^{\prime}(v)\right| \mathrm{d} v \\
& =\int_{t}^{+\infty}\left|h^{\prime}(v)\right| \mathrm{d} v<+\infty
\end{aligned}
$$

The last equality can be deduced from (6.13) by applying the preceding computation to $h(v)=\mathrm{e}^{-\left|v-t_{0}\right|}$, or one can check directly that $\gamma_{\alpha}=\int_{0}^{1} y^{\alpha-1}(1-$ $y)^{-\alpha} \mathrm{d} y$. From the Fubini theorem and the same calculation without absolute values, it follows that if $\int_{t_{0}}^{+\infty}\left|h^{\prime}(u)\right| \mathrm{d} u<+\infty$, then for every $t \geqslant t_{0}$ we have

$$
\left(I^{\alpha} D^{\alpha} h\right)(t)=-\int_{t}^{+\infty} h^{\prime}(u) \mathrm{d} u=h(t)
$$

Under (6.15), we introduce $h_{\varepsilon}(t)=\mathrm{e}^{-\varepsilon\left|t-t_{0}\right|} h(t)$ with $\varepsilon>0$, for which we use the preceding case and convergence when $\varepsilon \rightarrow 0$. When $\varepsilon \in(0,1)$ and $t>t_{0}$ we have

$$
\left|h_{\varepsilon}(t)\right| \leqslant|h(t)| \leqslant \frac{\kappa_{0}}{1+|t|}, \quad\left|h_{\varepsilon}^{\prime}(t)\right| \leqslant\left(\varepsilon|h(t)|+\left|h^{\prime}(t)\right|\right) \leqslant \frac{\kappa_{0}+\kappa_{1}}{1+|t|}
$$

By Corollary 6.10 , we have $\left|D^{\alpha} h_{\varepsilon}\right| \leqslant \kappa(1+|t|)^{-1}$, and we can apply twice dominated convergence when $\varepsilon \rightarrow 0$ in

$$
\int_{t}^{+\infty}(u-t)^{\alpha-1}\left(\int_{u}^{+\infty}(v-u)^{-\alpha} h_{\varepsilon}^{\prime}(v) \mathrm{d} v\right) \mathrm{d} u=h_{\varepsilon}(t)
$$

Assuming (6.15) and $\operatorname{Re} z>0$, we have

$$
\begin{equation*}
D_{t}^{z}(t h(t))=t\left(D^{z} h\right)(t)-z\left(D^{z-1} h\right)(t) . \tag{6.17}
\end{equation*}
$$

This is obtained when $0<\operatorname{Re} z<1$ with an integration by parts, writing

$$
\begin{aligned}
\Gamma(1-z)\left(-D_{t}^{z}(t h(t)\right. & )+t\left(D^{z} h\right)(t)\right) \\
= & \int_{t}^{+\infty}(u-t)^{-z}\left((u-t) h^{\prime}(u)+h(u)\right) \mathrm{d} u \\
= & \int_{t}^{+\infty}(u-t)^{-z+1} h^{\prime}(u) \mathrm{d} u+\int_{t}^{+\infty}(u-t)^{-z} h(u) \mathrm{d} u \\
= & z \int_{t}^{+\infty}(u-t)^{-z} h(u) \mathrm{d} u=z \Gamma(1-z)\left(D^{z-1} h\right)(t)
\end{aligned}
$$

### 6.2.1. Multipliers associated to fractional derivatives

If $K$ is a kernel integrable on $\mathbb{R}^{n}$, we know by (2.15) that its Fourier transform $m$ is expressed for $\xi \neq 0$ as

$$
m(u \xi)=\int_{\mathbb{R}} \varphi_{\theta}(s) \mathrm{e}^{-2 \mathrm{i} \pi s u|\xi|} \mathrm{d} s=\int_{\mathbb{R}} \frac{1}{|\xi|} \varphi_{\theta}\left(\frac{v}{|\xi|}\right) \mathrm{e}^{-2 \mathrm{i} \pi v u} \mathrm{~d} v, \quad u \in \mathbb{R}
$$

where $\theta=|\xi|^{-1} \xi$ and where the function $\varphi_{\theta}$ is defined on $\mathbb{R}$ by (2.14). Letting $\alpha>0$ and assuming that $x \mapsto|x|^{\alpha} K(x)$ is integrable on $\mathbb{R}^{n}$, this yields

$$
\begin{aligned}
D_{u}^{\alpha} m(u \xi) & =\int_{\mathbb{R}}(2 \mathrm{i} \pi v)^{\alpha} \frac{1}{|\xi|} \varphi_{\theta}\left(\frac{v}{|\xi|}\right) \mathrm{e}^{-2 \mathrm{i} \pi v u} \mathrm{~d} v \\
& =\int_{\mathbb{R}}(2 \mathrm{i} \pi s|\xi|)^{\alpha} \varphi_{\theta}(s) \mathrm{e}^{-2 \mathrm{i} \pi s|\xi| u} \mathrm{~d} s \\
& =\int_{\mathbb{R}^{n}}(2 \mathrm{i} \pi x \cdot \xi)^{\alpha} K(x) \mathrm{e}^{-2 \mathrm{i} \pi u x \cdot \xi} \mathrm{~d} x
\end{aligned}
$$

which is naturally extended by 0 when $\xi=0$. We set in what follows

$$
\begin{align*}
(\xi \cdot \nabla)^{\alpha} m(\xi):= & \left.D_{u}^{\alpha} m(u \xi)\right|_{u=1} \\
& =\int_{\mathbb{R}^{n}}(2 \mathrm{i} \pi x \cdot \xi)^{\alpha} K(x) \mathrm{e}^{-2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} x \\
& =\int_{\mathbb{R}}(2 \mathrm{i} \pi s|\xi|)^{\alpha} \varphi_{\theta}(s) \mathrm{e}^{-2 \mathrm{i} \pi s|\xi|} \mathrm{d} s
\end{align*}
$$

When $\alpha=1$ and $\xi \neq 0$, the quantity $(\xi \cdot \nabla)^{1} m(\xi)$ is equal to $-\xi \cdot \nabla m(\xi)$, which is the product by $-|\xi|$ of the usual directional derivative of the function $m$ in the direction of the norm-one vector $\theta=|\xi|^{-1} \xi$. When $0<\alpha<1$, under the assumptions (6.15), we can give according to Lemma 6.8 the integral formula

$$
\begin{equation*}
(\xi \cdot \nabla)^{\alpha} m(\xi)=-\frac{1}{\Gamma(1-\alpha)} \int_{1}^{+\infty}(u-1)^{-\alpha} \frac{\mathrm{d}}{\mathrm{~d} u} m(u \xi) \mathrm{d} u \tag{6.19}
\end{equation*}
$$

We shall use the integral formula (6.19) when $m(\xi)$ is Lipschitz outside the origin and when for every $u_{0}>0$ and $u \geqslant u_{0}$, we have for every $\theta \in S^{n-1}$ that

$$
|m(u \theta)|+\left|\frac{\mathrm{d}}{\mathrm{~d} u} m(u \theta)\right| \leqslant \frac{\kappa\left(\theta, u_{0}\right)}{1+|u|}
$$

If $K$ is an isotropic log-concave probability density with variance $\sigma^{2}$, we know by Corollary 5.13 that $|(\mathrm{d} / \mathrm{d} u) m(u \xi)| \leqslant \delta_{1, c}|\sigma \xi| /(1+2 \pi|u \sigma \xi|) \leqslant$ $\delta_{1, c} /(2 \pi|u|)$, thus

$$
\begin{equation*}
\left|(\xi \cdot \nabla)^{\alpha} m(\xi)\right| \leqslant \frac{\delta_{1, c}}{2 \pi|\Gamma(1-\alpha)|} \int_{1}^{+\infty}(u-1)^{-\alpha} u^{-1} \mathrm{~d} u=\kappa_{\alpha} \delta_{1, c} \tag{6.20}
\end{equation*}
$$

and the bounded function $\xi \mapsto(\xi \cdot \nabla)^{\alpha} m(\xi)$ defines an $L^{2}$ multiplier. We reach of course the same conclusion under (6.1.H) for a "general" kernel $K_{g}$.

We have seen in (2.10) that the multiplier norm of $m(\xi)$ on $L^{p}\left(\mathbb{R}^{n}\right)$ is the same as that of the dilate $m(\lambda \xi)$, for every $\lambda>0$. It is thus natural to look for a norm invariant by dilation, if we want a norm capable to control the action on $L^{p}$ of a multiplier. Since we shall work radially with Carbery's approach, we begin with a smooth function $h$ compactly supported in $(0,+\infty)$, and when $\alpha \in(0,1)$ we set with Carbery [21]

$$
\begin{equation*}
\|h\|_{L_{\alpha}^{2}}:=\left(\int_{0}^{+\infty}\left|t^{\alpha+1} D_{t}^{\alpha}\left(\frac{h(t)}{t}\right)\right|^{2} \frac{\mathrm{~d} t}{t}\right)^{1 / 2} \tag{6.21}
\end{equation*}
$$

One verifies that this norm is invariant by dilation. By (6.8), we have

$$
\begin{equation*}
t^{\alpha+1} D_{t}^{\alpha}\left(\frac{h_{[\lambda]}(t)}{t}\right)=t^{\alpha+1} \lambda D_{t}^{\alpha}\left(\frac{h(\lambda t)}{\lambda t}\right)=\left.(\lambda t)^{\alpha+1} D_{v}^{\alpha}\left(\frac{h(v)}{v}\right)\right|_{v=\lambda t} \tag{6.22}
\end{equation*}
$$

and the change of variable $u=\lambda t$ in (6.21) completes the proof. Let $h$ be Lipschitz on $\left(t_{0},+\infty\right)$ for all $t_{0}>0$. Applying (6.17) to $\widetilde{h}(t)=h(t) / t$, we get for all $t>0$

$$
\begin{align*}
D_{t}^{\alpha}\left(\frac{h(t)}{t}\right) & =\frac{\alpha}{t} D_{t}^{\alpha-1}\left(\frac{h(t)}{t}\right)+\frac{1}{t}\left(D^{\alpha} h\right)(t) \\
& =\frac{1}{t} D_{t}^{\alpha-1}\left(\frac{\alpha h(t)}{t}-h^{\prime}(t)\right) . \tag{6.23}
\end{align*}
$$

Remark 6.12. - When $1 / 2<\alpha<1$, the $L_{\alpha}^{2}$ norm dominates the $L^{\infty}(0,+\infty)$ norm of the function $h$. For a justification, let us assume in addition that $h$ is bounded and Lipschitz on each interval $(t,+\infty)$ with $t>0$. Then $H: u \mapsto h(u) / u$ satisfies (6.15) on $(t,+\infty)$ and we can apply

Lemma 6.11, giving $I^{\alpha} D^{\alpha} H=H$, thus

$$
\begin{aligned}
\frac{h(t)}{t} & =\frac{1}{\Gamma(\alpha)} \int_{t}^{+\infty}(u-t)^{\alpha-1} D_{u}^{\alpha}\left(\frac{h(u)}{u}\right) \mathrm{d} u \\
& =\frac{1}{t \Gamma(\alpha)} \int_{t}^{+\infty}(t / u)(1-t / u)^{\alpha-1}\left[u^{\alpha+1} D_{u}^{\alpha}\left(\frac{h(u)}{u}\right)\right] \frac{\mathrm{d} u}{u}
\end{aligned}
$$

Applying Cauchy-Schwarz, $\Gamma(\alpha)>1$ for $\alpha \in(0,1)$ and letting $y=t / u$, we get

$$
\begin{aligned}
h(t)^{2} & \leqslant\left(\int_{t}^{+\infty}(t / u)^{2}(1-t / u)^{2 \alpha-2} \frac{\mathrm{~d} u}{u}\right)\left(\int_{t}^{+\infty}\left[u^{\alpha+1} D_{u}^{\alpha}\left(\frac{h(u)}{u}\right)\right]^{2} \frac{\mathrm{~d} u}{u}\right) \\
& \leqslant\left(\int_{0}^{1} y(1-y)^{2 \alpha-2} \mathrm{~d} y\right)\|h\|_{L_{\alpha}^{2}}^{2} \leqslant \frac{1}{2 \alpha-1}\|h\|_{L_{\alpha}^{2}}^{2}
\end{aligned}
$$

The latter calculation is the basis for the $L^{2}$ part of Carbery's Proposition 6.14.

Remark 6.13. - Using the second expression in (6.23), we see that $\|h\|_{L_{\alpha}^{2}}^{2}$ is the integral on $(0,+\infty)$, and with respect to $(\mathrm{d} t) / t$, of the square of the modulus of

$$
\begin{aligned}
t^{\alpha} D_{t}^{\alpha-1}\left(\frac{\alpha h(t)}{t}-h^{\prime}(t)\right) & =\frac{1}{\Gamma(1-\alpha)} \int_{t}^{+\infty}(u / t-1)^{-\alpha}\left(\alpha h(u)-u h^{\prime}(u)\right) \frac{\mathrm{d} u}{u} \\
& =\frac{1}{\Gamma(1-\alpha)} \int_{1}^{+\infty}(v-1)^{-\alpha}\left(\alpha h(t v)-t v h^{\prime}(t v)\right) \frac{\mathrm{d} v}{v}
\end{aligned}
$$

In most cases, this expression tends to $\kappa h(0)$ when $t \rightarrow 0$, with $\kappa>0$, and then we have that $\|h\|_{L_{\alpha}^{2}}$ is finite only if $h(0)=0$, as for Bourgain's criterion $\Gamma_{B}(K)$.

We do not see an easy way to compare the $L_{\alpha}^{2}$ norm and the quantity appearing in the $\Gamma_{B}$ criterion. However, in the very special case where $H(t)=$ $h(t) / t$ is $\geqslant 0$, convex and decreasing on $(0,+\infty)$, the function $\left|H^{\prime}\right|=-H^{\prime}$ is decreasing and it follows from Lemma 6.9 that $\left(D^{1 / 2} H\right)(t)$ is bounded by $\kappa \sqrt{\left|H(t) H^{\prime}(t)\right|}$, hence

$$
\begin{aligned}
\|h\|_{L_{1 / 2}^{2}}^{2} & \leqslant \kappa \int_{0}^{+\infty} t^{3} \frac{|h(t)|}{t}\left(\frac{\left|h^{\prime}(t)\right|}{t}+\frac{|h(t)|}{t^{2}}\right) \frac{\mathrm{d} t}{t} \\
& \leqslant \kappa \int_{0}^{+\infty}\left(|h(t)|\left|t h^{\prime}(t)\right|+|h(t)|^{2}\right) \frac{\mathrm{d} t}{t} \\
& \leqslant \kappa^{\prime} \sum_{j \in \mathbb{Z}}\left(\alpha_{j}(h) \beta_{j}(h)+\alpha_{j}(h)^{2}\right)
\end{aligned}
$$

We obtain then (in this very special situation) that $\|h\|_{L_{1 / 2}^{2}} \leqslant \kappa \Gamma_{B}\left(h^{\vee}\right)$.

### 6.3. Fourier criteria for bounding the maximal function

In the next proposition due to Carbery, we impose conditions that fit into our presentation but are certainly too restrictive.

Proposition 6.14 (Carbery [21]). - Let $K$ be a kernel integrable on $\mathbb{R}^{n}$ with integral equal to 0 , let $m$ be the Fourier transform of $K$. Assume that $m_{\theta}:=u \mapsto m(u \theta)$ is differentiable on $(0,+\infty)$ for every $\theta \in S^{n-1}$, and that $m_{\theta}^{\prime}(u), u>0$, is bounded by a constant independent of $\theta$.
(1) If there exists $\alpha \in(1 / 2,1)$ such that

$$
\begin{equation*}
C_{\alpha}(m):=\sup _{\theta \in S^{n-1}}\|t \mapsto m(t \theta)\|_{L_{\alpha}^{2}}<+\infty \tag{6.24}
\end{equation*}
$$

then for every function $f \in L^{2}\left(\mathbb{R}^{n}\right)$ one has
$\left\|M_{K} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}=\left\|\sup _{t>0}\left|K_{(t)} * f\right|\right\|_{L^{2}\left(\mathbb{R}^{n}\right)} \leqslant \frac{1}{\sqrt{2 \alpha-1}} C_{\alpha}(m)\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)}$.
(2) Suppose that $p<+\infty$ and $1 / p<\alpha<1$. If the multiplier $(\xi \cdot \nabla)^{\alpha} m(\xi)$ from $\left(6.18 . \nabla^{\alpha}\right)$ is bounded on $L^{p}\left(\mathbb{R}^{n}\right)$, then for every $f$ in $L^{p}\left(\mathbb{R}^{n}\right)$ one has that

$$
\begin{array}{rl}
\| \sup _{1 \leqslant t \leqslant 2} \mid K_{(t)} * & f \mid \|_{L^{p}\left(\mathbb{R}^{n}\right)} \\
& \leqslant \kappa_{\alpha, p}\left(2\|m\|_{p \rightarrow p}+\left\|(\xi \cdot \nabla)^{\alpha} m(\xi)\right\|_{p \rightarrow p}\right)\|f\|_{L^{p}\left(\mathbb{R}^{n}\right)}  \tag{6.25}\\
\text { with } \kappa_{\alpha, p} & \leqslant(2 \alpha)^{1-1 / p}(p-1)^{1-2 / p}(\alpha-1 / p)^{1 / p-1}
\end{array}
$$

When $1<p \leqslant 2$, one has the simpler larger bound $\kappa_{\alpha, p} \leqslant \sqrt{2}(\alpha-$ $1 / p)^{-1 / p}$. Indeed, for $0<\alpha<1$, we have that $2^{1 / 2-1 / p} \alpha^{1-1 / p}(p-1)^{1-2 / p}(\alpha-$ $1 / p)^{2 / p-1}$ is less than $([\alpha-1 / p] /[\sqrt{2 \alpha}(p-1)])^{2 / p-1}$. When $1<p \leqslant 2$, this expression increases with $\alpha \in(1 / p, 1]$, and for $\alpha=1$, one has $(1-1 / p) /(\sqrt{2}(p-1))=1 /(\sqrt{2} p) \leqslant 1$.

Observe that if we set $\xi=|\xi| \theta$ for some nonzero vector $\xi \in \mathbb{R}^{n}$, we have

$$
\|t \mapsto m(t \xi)\|_{L_{\alpha}^{2}}=\|t \mapsto m(t \theta)\|_{L_{\alpha}^{2}}
$$

according to the invariance by dilation (6.22) of the norm $L_{\alpha}^{2}$. So the supremum in (1) is also the supremum on $\xi \in \mathbb{R}^{n}$. We shall need the following Lemma, slightly more general than the conclusion (1) in Proposition 6.14.

Lemma 6.15. - Let $\left(K_{t}\right)_{t>0}$ be a family of integrable kernels on $\mathbb{R}^{n}$, and denote by $\xi \mapsto m(\xi, t)$ the Fourier transform of $K_{t}$. Assume that for
every $u_{0}>0$, there exist $N$ and $\kappa\left(u_{0}\right)$ satisfying the following: for every $\xi$ in $\mathbb{R}^{n}$, the function $g_{\xi}: u \mapsto m(\xi, u) / u$, for $u \in\left[u_{0},+\infty\right)$, is Lipschitz and

$$
\begin{equation*}
\left|g_{\xi}(u)\right|+\left|g_{\xi}^{\prime}(u)\right| \leqslant \kappa\left(u_{0}\right) \frac{(1+|\xi|)^{N}}{1+|u|}, \quad \xi \in \mathbb{R}^{n}, u \geqslant u_{0} \tag{6.26}
\end{equation*}
$$

If there is $\alpha \in(1 / 2,1)$ such that $c_{\alpha}:=\sup _{\xi \in \mathbb{R}^{n}}\|t \mapsto m(\xi, t)\|_{L_{\alpha}^{2}}<+\infty$, then

$$
\forall f \in \mathcal{S}\left(\mathbb{R}^{n}\right), \quad\left\|\sup _{t>0}\left|K_{t} * f\right|\right\|_{L^{2}\left(\mathbb{R}^{n}\right)} \leqslant \frac{1}{\sqrt{2 \alpha-1}} c_{\alpha}\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)}
$$

Proof. - By the assumptions, the function $g_{\xi}$ satisfies (6.15). As in Remark 6.12 , we obtain by Lemma 6.11 for all $\xi \in \mathbb{R}^{n}$ and $t>0$ that

$$
\frac{m(\xi, t)}{t}=\frac{1}{\Gamma(\alpha)} \int_{t}^{+\infty}(u-t)^{\alpha-1} D_{u}^{\alpha}\left(\frac{m(\xi, u)}{u}\right) \mathrm{d} u
$$

For $f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$, according to (6.26) and Corollary 6.10, we can use Fubini and get

$$
\begin{aligned}
& \left(K_{t} * f\right)(x)=\int_{\mathbb{R}^{n}} m(\xi, t) \widehat{f}(\xi) \mathrm{e}^{2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} \xi \\
& =\frac{1}{\Gamma(\alpha)} \int_{t}^{+\infty} t(u-t)^{\alpha-1} \int_{\mathbb{R}^{n}} D_{u}^{\alpha}\left(\frac{m(\xi, u)}{u}\right) \widehat{f}(\xi) \mathrm{e}^{2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} \xi \mathrm{~d} u \\
& =\frac{1}{\Gamma(\alpha)} \int_{t}^{+\infty}(t / u)(1-t / u)^{\alpha-1}\left(\int_{\mathbb{R}^{n}} u^{\alpha+1} D_{u}^{\alpha}\left(\frac{m(\xi, u)}{u}\right) \widehat{f}(\xi) \mathrm{e}^{2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} \xi\right) \frac{\mathrm{d} u}{u}
\end{aligned}
$$

For $u>0$ and $x \in \mathbb{R}^{n}$, let us set

$$
\left(P_{u}^{\alpha} f\right)(x)=\int_{\mathbb{R}^{n}} u^{\alpha+1} D_{u}^{\alpha}\left(\frac{m(\xi, u)}{u}\right) \widehat{f}(\xi) \mathrm{e}^{2 \mathrm{i} \pi x \cdot \xi} \mathrm{~d} \xi
$$

This operator $P_{u}^{\alpha}$ is associated to the multiplier

$$
p_{u}^{\alpha}(\xi)=\left.u^{\alpha+1} D_{v}^{\alpha}\left(\frac{m(\xi, v)}{v}\right)\right|_{v=u}, \quad \xi \in \mathbb{R}^{n}
$$

One can rewrite

$$
\begin{equation*}
\left(K_{t} * f\right)(x)=\frac{1}{\Gamma(\alpha)} \int_{t}^{+\infty}(t / u)(1-t / u)^{\alpha-1}\left(P_{u}^{\alpha} f\right)(x) \frac{\mathrm{d} u}{u} \tag{6.27}
\end{equation*}
$$

By Cauchy-Schwarz and since $\Gamma(\alpha)>1$ when $\alpha \in(0,1)$, we get

$$
\left|\left(K_{t} * f\right)(x)\right|^{2} \leqslant\left(\int_{t}^{+\infty}(t / u)^{2}(1-t / u)^{2(\alpha-1)} \frac{\mathrm{d} u}{u}\right)\left(\int_{0}^{+\infty}\left|\left(P_{u}^{\alpha} f\right)(x)\right|^{2} \frac{\mathrm{~d} u}{u}\right)
$$

For $\alpha>1 / 2$, one has $2(\alpha-1)>-1$ and letting $y=t / u$, one sees that

$$
\int_{t}^{+\infty}(t / u)^{2}(1-t / u)^{2(\alpha-1)} \frac{\mathrm{d} u}{u}=\int_{0}^{1} y(1-y)^{2(\alpha-1)} \mathrm{d} y<\frac{1}{2 \alpha-1}
$$

We have obtained for $\left|\left(K_{t} * f\right)(x)\right|^{2}$ a bound independent of $t$, hence

$$
\sup _{t>0}\left|\left(K_{t} * f\right)(x)\right|^{2} \leqslant \kappa_{\alpha}^{2}\left(\int_{0}^{+\infty}\left|\left(P_{u}^{\alpha} f\right)(x)\right|^{2} \frac{\mathrm{~d} u}{u}\right)
$$

with $\kappa_{\alpha}^{-2}=2 \alpha-1$. By Fubini and Parseval, we have

$$
\begin{aligned}
\left\|\sup _{t>0}\left|K_{t} * f\right|\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}^{2} & \leqslant \kappa_{\alpha}^{2} \int_{\mathbb{R}^{n}}\left(\int_{0}^{+\infty}\left|\left(P_{u}^{\alpha} f\right)(x)\right|^{2} \frac{\mathrm{~d} u}{u}\right) \mathrm{d} x \\
& =\kappa_{\alpha}^{2} \int_{0}^{+\infty}\left\|P_{u}^{\alpha} f\right\|_{L^{2}\left(\mathbb{R}^{n}\right)}^{2} \frac{\mathrm{~d} u}{u} \\
& =\kappa_{\alpha}^{2} \int_{\mathbb{R}^{n}} \int_{0}^{+\infty}\left|u^{\alpha+1} D_{u}^{\alpha}\left(\frac{m(\xi, u)}{u}\right) \widehat{f}(\xi)\right|^{2} \frac{\mathrm{~d} u}{u} \mathrm{~d} \xi \\
& \leqslant \kappa_{\alpha}^{2} \int_{\mathbb{R}^{n}} c_{\alpha}^{2}|\widehat{f}(\xi)|^{2} \mathrm{~d} \xi=\kappa_{\alpha}^{2} c_{\alpha}^{2}\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)}^{2} .
\end{aligned}
$$

Remark 6.16. - If $|a(t)| \leqslant c\left(t_{0}\right)$ when $t \geqslant t_{0}>0$ and if $b(t)=a(t) / t$, then we have $(1+t)|b(t)|=\left(t^{-1}+1\right)|a(t)| \leqslant c\left(t_{0}\right)\left(1+t_{0}^{-1}\right)$ when $t \geqslant t_{0}$. If we add that $\left|a^{\prime}(t)\right| \leqslant c\left(t_{0}\right)$ for $t \geqslant t_{0}$, we have also $(1+t)\left|a^{\prime}(t) / t\right| \leqslant c\left(t_{0}\right)\left(1+t_{0}^{-1}\right)$, $t \geqslant t_{0}$, and

$$
\left|b^{\prime}(t)\right| \leqslant\left(\frac{\left|a^{\prime}(t)\right|}{t}+\frac{|b(t)|}{t}\right) \leqslant \frac{c\left(t_{0}\right)\left(1+t_{0}^{-1}\right)^{2}}{1+t}, \quad t \geqslant t_{0}>0 .
$$

If we know that for every $u_{0}>0$, there is $c\left(u_{0}\right)$ such that

$$
|m(\xi, u)|+\left|\frac{\mathrm{d}}{\mathrm{~d} u} m(\xi, u)\right| \leqslant c\left(u_{0}\right)(1+|\xi|)^{N}, \quad \xi \in \mathbb{R}^{n}, u \geqslant u_{0}
$$

it follows that $(6.26)$ is true, with $\kappa\left(u_{0}\right) \leqslant 2 c\left(u_{0}\right)\left(1+t_{0}^{-1}\right)^{2}$.
Proof of Proposition 6.14. - We apply Lemma 6.15 to the family $K_{t}=$ $K_{(t)}$ of dilates of $K, t>0$. Under the assumptions of Proposition 6.14, we first have that $|m(t \xi)|+|(\mathrm{d} / \mathrm{d} t) m(t \xi)| \leqslant \kappa(1+|\xi|)$. Remark 6.16 implies then that the family of functions $g_{\xi}: t \mapsto m(t \xi) / t$ satisfies (6.26). We thus obtain by Lemma 6.15 the $L^{2}$-maximal inequality when $f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$, and we may extend it to all functions in $L^{2}\left(\mathbb{R}^{n}\right)$ by the density of $\mathcal{S}\left(\mathbb{R}^{n}\right)$ in $L^{2}\left(\mathbb{R}^{n}\right)$, as explained in Section 3.3.

Let us pass to the proof of (2), the $L^{p}$ case. We use the notation of the proof of Lemma 6.15 , adapted to $m(\xi, t)=m(t \xi)$. Denote by $q$ the conjugate exponent of $p$, and observe that $q-2>-1$ because $p<+\infty$. When $\alpha \in(1 / p, 1)$ and $t \geqslant 1$, by applying Hölder to (6.27) and since $\alpha-1>-1 / q$,
$\Gamma(\alpha)>1$, we obtain

$$
\begin{aligned}
& \left|\left(K_{(t)} * f\right)(x)\right| \\
& \leqslant \Gamma(\alpha)^{-1}\left(\int_{t}^{+\infty}(t / u)^{q}(1-t / u)^{q(\alpha-1)} \mathrm{d} u\right)^{1 / q}\left(\int_{t}^{+\infty}\left|\left(P_{u}^{\alpha} f\right)(x)\right|^{p} \frac{\mathrm{~d} u}{u^{p}}\right)^{1 / p} \\
& \leqslant t^{1 / q}\left(\int_{t}^{+\infty}(t / u)^{q}(1-t / u)^{q(\alpha-1)} \frac{\mathrm{d} u}{t}\right)^{1 / q}\left(\int_{t}^{+\infty}\left|\left(P_{u}^{\alpha} f\right)(x)\right|^{p} \frac{\mathrm{~d} u}{u^{p}}\right)^{1 / p} \\
& \leqslant t^{1 / q}\left(\int_{1}^{+\infty} v^{-q \alpha}(v-1)^{q(\alpha-1)} \mathrm{d} v\right)^{1 / q}\left(\int_{1}^{+\infty}\left|\left(P_{u}^{\alpha} f\right)(x)\right|^{p} \frac{\mathrm{~d} u}{u^{p}}\right)^{1 / p} \\
& \leqslant\left(\int_{1}^{2}(v-1)^{q \alpha-q} \mathrm{~d} v+\int_{2}^{+\infty}(v-1)^{-q} \mathrm{~d} v\right)^{1 / q} \\
& \times t^{1 / q}\left(\int_{1}^{+\infty}\left|\left(P_{u}^{\alpha} f\right)(x)\right|^{p} \frac{\mathrm{~d} u}{u^{p}}\right)^{1 / p}
\end{aligned}
$$

With $c_{\alpha, p}^{q}=1 /(q \alpha-q+1)+1 /(q-1)=\alpha(p-1) /(\alpha-1 / p)$, it follows that

$$
\begin{aligned}
\left\|\sup _{1 \leqslant t \leqslant 2}\left|K_{(t)} * f\right|\right\|_{L^{p}\left(\mathbb{R}^{n}\right)}^{p} & \leqslant c_{\alpha, p}^{p} 2^{p / q} \int_{1}^{+\infty}\left(\int_{\mathbb{R}^{n}}\left|\left(P_{u}^{\alpha} f\right)(x)\right|^{p} \mathrm{~d} x\right) \frac{\mathrm{d} u}{u^{p}} \\
& =c_{\alpha, p}^{p} 2^{p / q} \int_{1}^{+\infty}\left\|P_{u}^{\alpha} f\right\|_{p}^{p} \frac{\mathrm{~d} u}{u^{p}} \\
& \leqslant \frac{c_{\alpha, p}^{p} 2^{p / q}}{p-1} \sup _{u \geqslant 1}\left\|P_{u}^{\alpha} f\right\|_{p}^{p}
\end{aligned}
$$

and we shall see that $\left\|P_{u}^{\alpha}\right\|_{p \rightarrow p} \leqslant 2\|m\|_{p \rightarrow p}+\left\|(\xi \cdot \nabla)^{\alpha} m(\xi)\right\|_{p \rightarrow p}$. The multipliers $p_{u}^{\alpha}$ are dilates of one another, indeed, for every $\lambda>0$, we have by (6.22) that

$$
\begin{aligned}
p_{u}^{\alpha}(\lambda \xi) & =\left.u^{\alpha+1} D_{v}^{\alpha}\left(\frac{m(v \lambda \xi)}{v}\right)\right|_{v=u} \\
& =\left.u^{\alpha+1} \lambda^{\alpha+1} D_{v}^{\alpha}\left(\frac{m(v \xi)}{v}\right)\right|_{v=\lambda u}=p_{\lambda u}^{\alpha}(\xi)
\end{aligned}
$$

It suffices therefore to consider $p_{1}^{\alpha}$. According to (6.23), one has

$$
p_{1}^{\alpha}(\xi)=\left.D_{t}^{\alpha}\left(\frac{m(t \xi)}{t}\right)\right|_{t=1}=\left.\alpha D_{t}^{\alpha-1}\left(\frac{m(t \xi)}{t}\right)\right|_{t=1}+\left.D_{t}^{\alpha} m(t \xi)\right|_{t=1}
$$

The multiplier $\left.D_{t}^{\alpha} m(t \xi)\right|_{t=1}$ is precisely equal to $(\xi \cdot \nabla)^{\alpha} m(\xi)$. The other term, since $\alpha-1<0$, can be written by (6.9) as

$$
U(\xi)=\left.\alpha D_{t}^{\alpha-1}\left(\frac{m(t \xi)}{t}\right)\right|_{t=1}=\frac{\alpha}{\Gamma(1-\alpha)} \int_{1}^{+\infty}(v-1)^{-\alpha}\left(\frac{m(v \xi)}{v}\right) \mathrm{d} v
$$

By Lemma 2.1, we have $\|U\|_{p \rightarrow p} \leqslant 2\|m\|_{p \rightarrow p}$ because

$$
\frac{\alpha}{\Gamma(1-\alpha)} \int_{1}^{+\infty}(v-1)^{-\alpha} \frac{\mathrm{d} v}{v} \leqslant \frac{\alpha}{\Gamma(1-\alpha)}\left(\frac{1}{1-\alpha}+\frac{1}{\alpha}\right)=\frac{1}{\Gamma(2-\alpha)} \leqslant 2
$$

cutting $\int_{1}^{+\infty}$ at $v=2$, and using (3.7).

### 6.4. Proofs of Theorems 6.1 and 6.2, and Proposition 6.3

We need only show Proposition 6.3 , and we can limit ourselves to $1<$ $p \leqslant 2$. As in Bourgain's proof of the $L^{2}$ theorem for $\mathrm{M}_{C}$ at the end of Section 5.3.1, the kernel $K$ to which we shall apply Proposition 6.6 is given by $K=K_{g}-P$, where $P$ is the Poisson kernel $P_{1}$ from (1.32), and $K_{g}$ is a probability density on $\mathbb{R}^{n}$ satisfying (6.1.H) with two constants $\delta_{0, g} \geqslant 1$ and $\delta_{1, g}$ controlling the Fourier transform $m_{g}$ and its gradient. We know by $\left(1.31 . P^{*}\right)$ that the maximal operator associated to the Poisson kernel acts on $L^{r}\left(\mathbb{R}^{n}\right), 1<r \leqslant+\infty$, with constants independent of the dimension $n$. Letting $B$ denote the Euclidean ball normalized by variance in $\mathbb{R}^{n}$, we could replace $P$ by $K_{B}$ and invoke Stein's Theorem 4.1 instead.

We shall apply Proposition 6.6 in the two cases corresponding to Theorems 6.1 and 6.2, in order to show that the maximal function (or the dyadic maximal function) associated to the kernel $K$ is bounded on $L^{p}$ for $p>3 / 2$ (or for $p>1$ ). We shall get by difference that the maximal function for $K_{g}$ (or $K_{l c}, K_{C}$ ) is bounded as well. In the "dyadic" case of Theorem 6.1, the operator $T_{j}$, for $j \in \mathbb{Z}$, is the convolution with the dilate $K_{\left(2^{j}\right)}$ of $K$. For Theorem 6.2, $T_{j, v}$ is the convolution with $K_{\left(v 2^{j}\right)}, 1 \leqslant v \leqslant 2$, and $T_{j}$ is given by

$$
T_{j} f=\sup _{1 \leqslant v \leqslant 2}\left|T_{j, v} f\right|=\sup _{2^{j} \leqslant t \leqslant 2^{j+1}}\left|K_{(t)} * f\right| .
$$

One has to check that the assumptions of Proposition 6.6, namely, $\left(A_{\mathbf{0}}\right)$, $\left(A_{\mathbf{1}}\right),\left(A_{\mathbf{2}}\right)$ and $\left(A_{\mathbf{3}}\right)$, are satisfied in these two cases. If the $\left(Q_{j}\right)$ are those of Littlewood-Paley, defined by

$$
\widehat{Q}_{j}(\xi)=\mathrm{e}^{-2 \pi 2^{j}|\xi|}-\mathrm{e}^{-2 \pi 2^{j+1}|\xi|}, \quad \xi \in \mathbb{R}^{n}
$$

then the assumption $\left(A_{\mathbf{0}}\right)$ is satisfied according to (2.4), with $C_{p}=\mathrm{q}_{p}$.
For $\left(A_{\mathbf{1}}\right)$, we write $T_{j, v}=U_{j, v}-S_{j, v}$, where the $U_{j, v}=\left(K_{g}\right)_{\left(v 2^{j}\right)}$ are related to $K_{g}$ and the $S_{j, v}=P_{\left(v 2^{j}\right)}$ to the Poisson kernel. The operators $U_{j, v}$ and $S_{j, v}$ are positive, as convolutions with probability densities. As mentioned before, this is the only place where we need $K_{g}$ to be a probability density rather than a general integrable kernel. We know by (1.31. $P^{*}$ ) that the maximal operator $S^{*}$ associated to the Poisson kernel is bounded on
$L^{p}\left(\mathbb{R}^{n}\right), 1<p<+\infty$, by a constant $C_{p}^{\prime}$ independent of the dimension $n$. Consequently, the property $\left(A_{1}\right)$ is satisfied.

Let us consider $\left(A_{\mathbf{2}}\right)$. The first case is when $T_{j}=K_{\left(2^{j}\right)}$ and in this case, according to (2.13), the operator $T_{j}$ is bounded on all the spaces $L^{p}\left(\mathbb{R}^{n}\right)$, $1 \leqslant p \leqslant+\infty$, by the $L^{1}$ norm of $K$ and we get that

$$
\begin{equation*}
\left\|T_{j}\right\|_{p \rightarrow p} \leqslant\|K\|_{L^{1}\left(\mathbb{R}^{n}\right)} \leqslant 2 \tag{6.28}
\end{equation*}
$$

In the second case, we have to use the part (2) of Proposition 6.14. This will be discussed below.

Finally, we must show $\left(A_{\mathbf{3}}\right)$, i.e., prove that $T^{*}$ satisfies the property $\left(\mathbf{S}_{2}\right)$. For $k$ fixed in $\mathbb{Z}$, we shall bound the maximal operator of the kernel $N_{k}:=K * Q_{k}$ using the conclusion (1) of Proposition 6.14. We show in Section 6.5 that for every value $\alpha \in(1 / 2,1)$, the norm $b_{k}:=C_{\alpha}\left(\widehat{N_{k}}\right)$ decays exponentially with $|k|$, with constants depending on $\alpha$ and (linearly) on $\delta_{0, g}+\delta_{1, g}$. In the "dyadic case", the bound obtained in this way by (1) for the maximal function of $N_{k}$ implies that

$$
\begin{aligned}
\left\|\sup _{j \in \mathbb{Z}}\left|T_{j} Q_{j+k} f\right|\right\|_{2} & =\left\|\sup _{j \in \mathbb{Z}}\left|K_{\left(2^{j}\right)} *\left(P_{\left(2^{j+k}\right)}-P_{\left(2^{j+1+k}\right)}\right) * f\right|\right\|_{2} \\
& =\left\|\sup _{j \in \mathbb{Z}}\left|\left(N_{k}\right)_{\left(2^{j}\right)} * f\right|\right\|_{2} \leqslant\left\|\sup _{t>0}\left|\left(N_{k}\right)_{(t)} * f\right|\right\|_{2} \leqslant \kappa_{\alpha} b_{k}
\end{aligned}
$$

which proves $\left(\mathbf{S}_{2}\right)$ in this case. The case of the global maximal function requires a small adaptation, Carbery says: "This is not exactly what being strongly bounded on $L^{2}$ means, but a slight modification of this argument will give precisely what we require". Indeed, there is now a gap between what we get from Proposition 6.14 and the assumption we need for applying Proposition 6.6. We shall discuss it in the subsection 6.4.1 and resolve this "gap question" in the subsection 6.5.1. We obtain at last by Lemma 6.19 and by Lemma 6.15 that there exist universal coefficients $\left(a_{k}\right)_{k \in \mathbb{Z}}$ such that $\sum_{k \in \mathbb{Z}} a_{k}^{s}<+\infty$ for every $s>0$, and such that

$$
\begin{equation*}
\left\|\sup _{j \in \mathbb{Z}}\left|T_{j} Q_{j+k} f\right|\right\|_{2} \leqslant\left(\delta_{0, g}+\delta_{1, g}\right) a_{k}, \quad k \in \mathbb{Z} \tag{6.29}
\end{equation*}
$$

For $\left(A_{\mathbf{2}}\right)$ in the "global" case, we study the operators $\left(W_{t}\right)_{t>0}$ defined by

$$
W_{t} f=\sup _{t \leqslant u \leqslant 2 t}\left|K_{(u)} * f\right|, \quad t>0
$$

and we want to prove $\left(A_{\mathbf{2}}\right)$ for the family of $T_{j}=W_{2^{j}}$ from (6.2), with $j \in \mathbb{Z}$. Using the invariance by dilation (2.11) of multiplier norms, we see that the operators $W_{t}$ have the same norm when $t$ varies, hence we need to find a bound for $T_{0}=W_{1}$ only. For this, we want to apply the conclusion (2) of Proposition 6.14 , so we must show that the multipliers $m$ and $(\xi \cdot \nabla)^{\alpha} m(\xi)$
are bounded on $L^{p}\left(\mathbb{R}^{n}\right)$ for some $\alpha \in(1 / p, 1)$. For $m$ it is clear by the elementary fact (2.13).

For $(\xi \cdot \nabla)^{\alpha} m(\xi)$ we shall use complex interpolation between $(\xi \cdot \nabla)^{0} m(\xi)=$ $m(\xi)$ that acts on $L^{1}\left(\mathbb{R}^{n}\right)$, and $(\xi \cdot \nabla) m(\xi)$ that acts on $L^{2}\left(\mathbb{R}^{n}\right)$ since it is a bounded function on $\mathbb{R}^{n}$ by (6.20) and (6.1.H). We get by interpolation that the multiplier $(\xi \cdot \nabla)^{\alpha} m(\xi)$ is bounded on $L^{p}\left(\mathbb{R}^{n}\right)$, with $p$ given by

$$
\frac{1}{p}=\frac{1-\alpha}{1}+\frac{\alpha}{2}=1-\frac{\alpha}{2}
$$

and we need $1-\alpha / 2=1 / p<\alpha$ for applying (2), thus $1<3 \alpha / 2=3-3 / p$. We must therefore have $p>3 / 2$ in order to conclude. We see that the reason for the restriction on the values of $p$ in Theorem 6.2 is to be found precisely here.

This sketch is not fully accurate. For being able to interpolate, one must control in $L^{2}$ the values $\alpha=1+\mathrm{i} \tau$, for every $\tau$ real, which causes no difficulty, but also the values $\alpha=0+\mathrm{i} \tau$ in $L^{1}$, and this is more technical. The precise work, involving a slight modification of the strategy described here, is done in Section 7.3 when we are well embedded by Müller [59] in the mood for interpolation. For every $p \in(3 / 2,2]$, we shall then obtain for some $\alpha>1 / p$, function of $p$, a bound of the form $\left\|(\xi \cdot \nabla)^{\alpha} m(\xi)\right\|_{p \rightarrow p} \leqslant \kappa_{p}\left(\delta_{0, g}+\delta_{1, g}\right)^{2-2 / p}$. By Proposition 6.14, we deduce

$$
\left\|T_{0} f\right\|_{L^{p}\left(\mathbb{R}^{n}\right)}=\left\|\sup _{1 \leqslant t \leqslant 2}\left|K_{(t)} * f\right|\right\|_{L^{p}\left(\mathbb{R}^{n}\right)} \leqslant \kappa_{p}^{\prime}\left(\delta_{0, g}+\delta_{1, g}\right)^{2-2 / p}\|f\|_{L^{p}\left(\mathbb{R}^{n}\right)}
$$

for every function $f \in L^{p}\left(\mathbb{R}^{n}\right)$. We get $\left(A_{\mathbf{2}}\right)$ with $p_{\min }=3 / 2$, since

$$
\begin{equation*}
\left\|T_{j}\right\|_{p \rightarrow p}=\left\|T_{0}\right\|_{p \rightarrow p} \leqslant \kappa_{p}^{\prime}\left(\delta_{0, g}+\delta_{1, g}\right)^{2-2 / p}, \quad j \in \mathbb{Z}, 3 / 2<p \leqslant 2 \tag{6.30}
\end{equation*}
$$

Applying Proposition 6.6, we finish the proof of Proposition 6.3. For $p \in(3 / 2,2]$, we shall bound $T^{*}=M_{K}$ in $L^{p}\left(\mathbb{R}^{n}\right)$, thus also $\mathrm{M}_{K_{g}}$. We choose a value $p_{0}$, function of $p$, such that $3 / 2<p_{0}<p$, and we let $\delta=\delta_{0, g}+\delta_{1, g}$. We have by (6.30) that $C_{p_{0}}^{\prime \prime} \leqslant \kappa_{p}^{\prime \prime} \delta^{2-2 / p_{0}}$. Then, applying (6.4), (6.29), (6.30) and $\delta_{0, g} \geqslant 1$, we obtain

$$
\begin{aligned}
\left\|\mathrm{M}_{K_{g}}\right\|_{p \rightarrow p} \leqslant\left\|T^{*}\right\|_{p \rightarrow p}+\kappa_{p, 0} & \leqslant \kappa_{p, 1}\left(C_{p_{0}}^{\prime \prime}\right)^{\gamma}\left(\sum_{k \in \mathbb{Z}}\left(\delta a_{k}\right)^{(1-\gamma) p / 2}\right)^{2 / p}+\kappa_{p, 2} \\
& \leqslant \kappa_{p} \delta^{2-2 / p}
\end{aligned}
$$

as announced, observing that $1-\gamma=\left[1 / p_{0}-1 / p\right] /\left[1 / p_{0}-1 / 2\right]$ is the interpolation parameter for $L^{p}$ and the pair $\left(L^{p_{0}}, L^{2}\right)$, and that the powers of $\delta$ under the exponents $\gamma$ and $1-\gamma$ are of the form $2-2 / r, r=p_{0}$ or 2 . In the dyadic case, we may replace (6.30) by (6.28) and obtain the result for $\mathrm{M}_{K_{g}}^{(d)}$ when $p \in(1,2]$.

Remark 6.17. - Bringing back the question to the Poisson kernel leads to some complications, because the function $\varphi_{\theta}(s)$ associated to the Poisson kernel, i.e., the Cauchy kernel (1.33.C), does not have decay properties as good as that of the function $\varphi_{\theta, C}$ of a convex set. This approach however does not depend on the $L^{p}$ result of Stein for the Euclidean ball.

Why not employ the Gaussian semi-group instead? In some non Euclidean situations, like Heisenberg groups or Grushin operators for instance, and especially for the weak type $(1,1)$ property of associated maximal functions, the Poisson kernel is preferable. Indeed, some asymptotic estimates, uniform in the dimension, are required on the kernel and are easier to obtain for the Poisson kernel. But in the Euclidean case, we cannot see a compelling obstacle to the use of the Gaussian kernel. We would get an excellent decay, both in the space variable and in the Fourier variable. We have chosen to stick to the original proofs, but we urge the reader to rewrite them with Gaussian kernels instead. We shall see in Section 8 that Bourgain uses Gaussian kernels.

### 6.4.1. Where is the gap?

As was said above, we will arrive for $N_{k}=K * Q_{k}$ at

$$
C_{\alpha}\left(N_{k}\right):=\sup _{\theta \in S^{n-1}}\left\|t \mapsto N_{k}(t \theta)\right\|_{L_{\alpha}^{2}} \leqslant \kappa_{\alpha} 2^{-\gamma|k|}, \quad k \in \mathbb{Z}
$$

for some $\gamma>0$. This implies by Proposition 6.14(1), that

$$
\left\|\sup _{t>0}\left|\left(N_{k}\right)_{(t)} * f\right|\right\|_{2} \leqslant \kappa_{\alpha} 2^{-\gamma|k|}
$$

Translating the definition of $N_{k}$ gives

$$
\left\|\sup _{t>0}\left|\left(K * Q_{k}\right)_{(t)} * f\right|\right\|_{2} \leqslant \kappa_{\alpha} 2^{-\gamma|k|}
$$

where $K=K_{g}-P$, or

$$
\| \sup _{v \in[1,2]} \sup _{j \in \mathbb{Z}} \mid\left(K_{\left(v 2^{j}\right)} *\left(P_{\left(v 2^{j+k}\right)}-P_{\left(v 2^{j+k+1}\right)}\right) * f \mid \|_{2} \leqslant \kappa_{\alpha} 2^{-\gamma|k|}\right.
$$

This must be compared to bounding the expression

$$
\| \sup _{v \in[1,2]} \sup _{j \in \mathbb{Z}} \mid\left(K_{\left(v 2^{j}\right)} *\left(P_{\left(2^{j+k}\right)}-P_{\left(2^{j+k+1}\right)}\right) * f \mid \|_{2}\right.
$$

which is what we are waiting for, in the definition of Property $\left(\mathbf{S}_{2}\right)$ for the family of operators $\left(T_{j, v}\right), j \in \mathbb{Z}, v \in[1,2]$.

### 6.5. A proof for the property $\left(\mathrm{S}_{2}\right)$

In what follows, $m=m_{g}-\widehat{P}$ is the Fourier transform of the kernel $K=K_{g}-P$ that appears in the proof of Proposition 6.3, where $K_{g}$ is a probability density on $\mathbb{R}^{n}$ satisfying (6.1.H). We have

$$
\widehat{P}(\xi)=\mathrm{e}^{-2 \pi|\xi|} \text { and we let } \rho(\xi)=\widehat{P}(\xi)-\widehat{P}(2 \xi), \quad \xi \in \mathbb{R}^{n}
$$

For every $k \in \mathbb{Z}$, every $\xi \in \mathbb{R}^{n}$ and $u>0$, we set

$$
\begin{gathered}
m_{k}(\xi)=\widehat{N_{k}}(\xi)=m(\xi)\left(\mathrm{e}^{-2^{k+1} \pi|\xi|}-\mathrm{e}^{-2^{k+2} \pi|\xi|}\right)=m(\xi) \rho\left(2^{k} \xi\right) \\
h_{k}^{\xi}(u)=\frac{m_{k}(u \xi)}{u}
\end{gathered}
$$

One must show that for any given $\alpha \in(1 / 2,1)$, the quantity

$$
C_{\alpha}\left(m_{k}\right)^{2}=\sup _{\theta \in S^{n-1}}\left\|u \mapsto m_{k}(u \theta)\right\|_{L_{\alpha}^{2}}^{2}=\sup _{\theta \in S^{n-1}} \int_{0}^{+\infty}\left(u^{\alpha+1}\left(D^{\alpha} h_{k}^{\theta}\right)(u)\right)^{2} \frac{\mathrm{~d} u}{u}
$$

introduced in (6.24) decays exponentially to 0 when $|k|$ tends to infinity. We fix therefore $\theta \in S^{n-1}$ and for $u \in \mathbb{R}$, we set

$$
\phi(u)=m(u \theta), \quad \chi(u)=\mathrm{e}^{-2 \pi|u|}-\mathrm{e}^{-4 \pi|u|}=\widehat{P}(u \theta)-\widehat{P}(2 u \theta)=\rho(u \theta) .
$$

Let $\delta=\delta_{0, g}+\delta_{1, g} \geqslant 1$, where $\delta_{0, g}, \delta_{1, g}$ are the constants in (6.1.H). We know that

$$
\begin{gathered}
|u|\left|m_{g}(u \theta)\right| \leqslant \delta_{0, g} \leqslant \delta,\left|\theta \cdot \nabla m_{g}(u \theta)\right| \leqslant \delta_{1, g} \leqslant \delta, \\
\left|u \theta \cdot \nabla m_{g}(u \theta)\right| \leqslant \delta, u \in \mathbb{R} .
\end{gathered}
$$

On the other hand, the derivative with respect to $u>0$ of $\widehat{P}(u \theta)=\mathrm{e}^{-2 \pi|u|}$ is bounded by $2 \pi$, and according to (5.21a), (5.21b), we have

$$
|u \widehat{P}(u \theta)| \leqslant(2 \pi \mathrm{e})^{-1}<1 \leqslant \delta, \quad\left|u \frac{\mathrm{~d}}{\mathrm{~d} u} \widehat{P}(u \theta)\right| \leqslant \mathrm{e}^{-1}<\delta
$$

For $\phi(u)=m(u \theta)=m_{g}(u \theta)-\widehat{P}(u \theta)$ we get $\left|\phi^{\prime}(u)\right| \leqslant \delta+2 \pi$. Using again $\delta>1$, we simplify this bound as $\left|\phi^{\prime}(u)\right|<8 \delta$. It follows first that $|\phi(u)| \leqslant$ $8 \delta|u|$, and

$$
\begin{equation*}
|\phi(u)| \leqslant 8 \delta\left(|u| \wedge|u|^{-1}\right), \quad\left|\phi^{\prime}(u)\right| \leqslant 8 \delta\left(1 \wedge|u|^{-1}\right) \tag{6.31a}
\end{equation*}
$$

For $\chi(u)$, we see when $u>0$ that $0 \leqslant \chi(u) \leqslant \mathrm{e}^{-2 \pi u}$ and

$$
-2 \pi \mathrm{e}^{-2 \pi u} \leqslant \chi^{\prime}(u)=-2 \pi \mathrm{e}^{-2 \pi u}+4 \pi \mathrm{e}^{-4 \pi u} \leqslant 2 \pi \mathrm{e}^{-2 \pi u}
$$

implying that $\left|\chi^{\prime}(u)\right| \leqslant 2 \pi$ for $u \neq 0$ and

$$
\begin{equation*}
|\chi(u)| \leqslant(2 \pi|u|) \wedge|2 \pi \mathrm{e} u|^{-1}, \quad\left|\chi^{\prime}(u)\right| \leqslant(2 \pi) \wedge|\mathrm{e} u|^{-1} \tag{6.31b}
\end{equation*}
$$

We obtain a symmetric treatment of the two functions $\chi$ and $\phi_{\delta}:=\delta^{-1} \phi$ since, up to some universal multiple $\kappa$ (we express this by the sign $\lesssim$ ), we have

$$
\begin{equation*}
\left|\phi_{\delta}(u)\right|,|\chi(u)| \lesssim|u| \wedge|u|^{-1}, \quad\left|\phi_{\delta}^{\prime}(u)\right|,\left|\chi^{\prime}(u)\right| \lesssim 1 \wedge|u|^{-1} \tag{6.32}
\end{equation*}
$$

We set $p_{k}(u)=m_{k}(u \theta)=\phi(u) \chi\left(2^{k} u\right), h_{k}(u)=p_{k}(u) / u$ and we want to estimate $\left\|p_{k}\right\|_{L_{\alpha}^{2}}$ for every $k \in \mathbb{Z}$. Notice that

$$
p_{-k}\left(2^{k} v\right)=\chi(v) \phi\left(2^{k} v\right) .
$$

The $L_{\alpha}^{2}$ norm is invariant by dilation and the assumptions on $\phi_{\delta}$ and $\chi$ are identical, we may therefore restrict the verification to the case $k \geqslant 0$. Let us fix an integer $k \geqslant 0$. We have the following table, divided into the three regions where the chosen bounds (6.32) for the functions $h_{k}$ and $h_{k}^{\prime}$ keep the same analytical expression, namely, the intervals $\left(0,2^{-k}\right),\left(2^{-k}, 1\right)$ and $(1,+\infty)$. We consider that $h_{k}^{\prime}$ is the derivative of the product of $u^{-1} \phi(u)$ and $\chi\left(2^{k} u\right)$, we bound therefore $\left|h_{k}^{\prime}\right|$ by the sum of $\left|\left(u^{-1} \phi(u)\right)^{\prime}\right|\left|\chi\left(2^{k} u\right)\right|$ and $\left|u^{-1} \phi(u)\right| 2^{k}\left|\chi^{\prime}\left(2^{k} u\right)\right|$.

| $u:$ | 0 | $2^{-k}$ |  | 1 |
| :---: | :---: | :---: | :---: | :---: |
| $u^{-1}\left\|\phi_{\delta}(u)\right\|$ | $\lesssim$ | 1 | 1 | $u^{-2}$ |
| $\left\|\chi\left(2^{k} u\right)\right\|$ | $\lesssim$ | $2^{k} u$ | $2^{-k} u^{-1}$ | $2^{-k} u^{-1}$ |
| $u^{-1}\left\|\phi_{\delta}^{\prime}(u)\right\|+u^{-2}\left\|\phi_{\delta}(u)\right\|$ | $\lesssim$ | $u^{-1}+u^{-1}$ | $u^{-1}+u^{-1}$ | $u^{-2}+u^{-3} \lesssim u^{-2}$ |
| $2^{k}\left\|\chi^{\prime}\left(2^{k} u\right)\right\|$ | $\lesssim$ | $2^{k}$ | $u^{-1}$ | $u^{-1}$ |
| $\delta^{-1}\left\|h_{k}(u)\right\|$ | $\lesssim$ | $2^{k} u \leqslant 2^{-k} u^{-1}$ | $2^{-k} u^{-1}$ | $2^{-k} u^{-3}$ |
| $\delta^{-1}\left\|h_{k}^{\prime}(u)\right\|$ | $\lesssim$ | $2^{k}+2^{k} \lesssim u^{-1}$ | $2^{-k} u^{-2}+u^{-1} \lesssim u^{-1}$ | $2^{-k} u^{-3}+u^{-3} \lesssim u^{-3}$ |

We see that $\delta^{-1}\left|h_{k}^{\prime}(u)\right| \lesssim H_{1}(u):=u^{-1} \wedge u^{-3}$. This function $H_{1}$ is nonincreasing on $(0,+\infty)$ and independent of $k$, and $\delta^{-1}\left|h_{k}(u)\right| \lesssim H_{0, k}(u)=$ $2^{-k} H_{1}(u)$. It follows from Lemma 6.9 that for $t>0$, we have

$$
\delta^{-1}\left|\left(D^{\alpha} h_{k}\right)(t)\right| \lesssim H_{0, k}(t)^{1-\alpha} H_{1}(t)^{\alpha} \lesssim 2^{-(1-\alpha) k} H_{1}(t)
$$

and the conclusion is reached since we obtain then

$$
\begin{aligned}
\left\|\phi \chi_{\left[2^{k}\right]}\right\|_{L_{\alpha}^{2}}^{2} & =\left\|p_{k}\right\|_{L_{\alpha}^{2}}^{2}=\int_{0}^{+\infty}\left|t^{\alpha+1}\left(D^{\alpha} h_{k}\right)(t)\right|^{2} \frac{\mathrm{~d} t}{t} \\
& \lesssim \delta^{2} 2^{-2(1-\alpha) k}\left(\int_{0}^{1}\left(t^{\alpha+1} t^{-1}\right)^{2} \frac{\mathrm{~d} t}{t}+\int_{1}^{\infty}\left(t^{\alpha+1} t^{-3}\right)^{2} \frac{\mathrm{~d} t}{t}\right)
\end{aligned}
$$

and

$$
\int_{0}^{1} t^{2 \alpha-1} \mathrm{~d} t+\int_{1}^{\infty} t^{2 \alpha-5} \mathrm{~d} t=\frac{1}{2 \alpha}+\frac{1}{4-2 \alpha}=\frac{1}{\alpha(2-\alpha)}<\frac{1}{\alpha}<+\infty
$$

thus $\left\|p_{k}\right\|_{L_{\alpha}^{2}} \lesssim \delta \alpha^{-1 / 2} 2^{-(1-\alpha) k}$ when $k \geqslant 0$, and $\left\|p_{k}\right\|_{L_{\alpha}^{2}} \leqslant \kappa \alpha^{-1 / 2} \delta 2^{-(1-\alpha)|k|}$ when $k \in \mathbb{Z}$. This implies by Proposition 6.14 (1) that

$$
\begin{equation*}
\left\|\sup _{r>0}\left(\left[\left(m_{k}\right)_{[r]} \widehat{f}\right]^{\vee}\right)\right\|_{L^{2}\left(\mathbb{R}^{n}\right)} \leqslant \kappa \delta 2^{-(1-\alpha)|k|}\|f\|_{L^{2}\left(\mathbb{R}^{n}\right)} \tag{6.33}
\end{equation*}
$$

for every $\alpha \in(1 / 2,1)$, giving the property $\left(\mathbf{S}_{2}\right)$ (see Definition 6.4) in the dyadic case.

It would be just as simple to work with the $\Gamma_{B}(K)$ criterion of Bourgain given in Section 5.3. We prove a general Lemma that will be invoked again in Section 8 for the cube problem.

Lemma 6.18. - Suppose that two integrable kernels $K_{1}$ and $K_{2}$ on $\mathbb{R}^{n}$ satisfy, for a certain $\kappa$ and every $\theta \in S^{n-1}$, that $\left|\widehat{K}_{j}(u \theta)\right| \leqslant \kappa\left(|u| \wedge|u|^{-1}\right), \quad\left|\theta \cdot \nabla \widehat{K}_{j}(u \theta)\right| \leqslant \kappa\left(1 \wedge|u|^{-1}\right), j=1,2, u \in \mathbb{R}$. It follows that $\Gamma_{B}\left(K_{1} *\left(K_{2}\right)_{\left(2^{k}\right)}\right) \leqslant C(\kappa) 2^{-|k| / 2}$ for $k \in \mathbb{Z}$.

Proof. - We fix $\theta \in S^{n-1}$, and in order to remind us about the preceding case, we let $m$ be the Fourier transform of $K_{1}$ and $\rho$ that of $K_{2}$. We will modify the table above, in order to emphasize now $\phi(u):=m(u \theta)$ and $u \theta \cdot \nabla m(u \theta)=u \phi^{\prime}(u)$ that appear in the components $\alpha_{j}(m)$ and $\beta_{j}(m)$ of $\Gamma_{B}(K)$, and we proceed similarly for $\chi(u):=\rho(u \theta)$.

Let $m_{k}$ be the Fourier transform of the kernel $K_{1} *\left(K_{2}\right)_{\left(2^{k}\right)}$. We have that $m_{k}(u \theta)=m(u \theta) \rho\left(2^{k} u \theta\right)$ and we may again restrict ourselves to $k \geqslant 0$, since a dilation by $2^{i}$ on a multiplier $g(\xi)$ produces a shift of $i$ places on the indices $j$ of the sequences $\left(\alpha_{j}(g)\right)_{j \in \mathbb{Z}},\left(\beta_{j}(g)\right)_{j \in \mathbb{Z}}$, leaving $\sum_{j \in \mathbb{Z}}$ unchanged. The bounds below do not depend on $\theta \in S^{n-1}$, so we will be able to estimate

$$
A_{k}(u):=\sup _{\theta \in S^{n-1}}\left|m_{k}(u \theta)\right| \quad \text { and } \quad B_{k}(u):=\sup _{\theta \in S^{n-1}}\left|u \theta \cdot \nabla m_{k}(u \theta)\right| .
$$

Note that $B_{k}(u)$ is controlled by $\phi(u) 2^{k} u \chi^{\prime}\left(2^{k} u\right)$ and $u \phi^{\prime}(u) \chi\left(2^{k} u\right)$. We have $\alpha_{j}\left(m_{k}\right) \sim A_{k}\left(2^{j}\right), \beta_{j}\left(m_{k}\right) \sim B_{k}\left(2^{j}\right)$, for every $j \in \mathbb{Z}$. The new table is divided into the same three regions as before.

| $u:$ | 0 | $2^{-k}$ |  | 1 |
| :---: | :---: | :---: | :---: | :---: |
| $\|\phi(u)\|$ | $\lesssim$ | $u$ | $u$ | $u^{-1}$ |
| $\left\|\chi\left(2^{k} u\right)\right\|$ | $\lesssim$ | $2^{k} u$ | $2^{-k} u^{-1}$ | $2^{-k} u^{-1}$ |
| $u\left\|\phi^{\prime}(u)\right\|$ | $\lesssim$ | $u$ | $u$ | 1 |
| $2^{k} u\left\|\chi^{\prime}\left(2^{k} u\right)\right\|$ | $\lesssim$ | $2^{k} u$ | 1 | 1 |
| $A_{k}(u)$ | $\lesssim$ | $2^{k} u^{2}$ | $2^{-k}$ | $2^{-k} u^{-2}$ |
| $B_{k}(u)$ | $\lesssim$ | $2^{k} u^{2}+2^{k} u^{2}$ | $u+2^{-k} \lesssim u$ | $u^{-1}+2^{-k} u^{-1} \lesssim u^{-1}$ |
| $\sqrt{A_{k}(u) B_{k}(u)}$ | $\lesssim$ | $2^{k} u^{2}$ | $2^{-k / 2} u^{1 / 2}$ | $2^{-k / 2} u^{-3 / 2}$ |

It follows that for every $j \in \mathbb{Z}$, we have

$$
\alpha_{j}\left(m_{k}\right) \lesssim \begin{cases}2^{k+2 j} & \text { if } j \leqslant-k, \\ 2^{-k} & \text { if }-k \leqslant j \leqslant 0, \quad \text { so } \sum_{j \in \mathbb{Z}} \alpha_{j}\left(m_{k}\right) \lesssim(k+1) 2^{-k} \\ 2^{-k-2 j} & \text { if } 0 \leqslant j,\end{cases}
$$

and

$$
\begin{aligned}
& \sqrt{\alpha_{j}\left(m_{k}\right) \beta_{j}\left(m_{k}\right)} \lesssim \begin{cases}2^{k+2 j} & \text { if } j \leqslant-k \\
2^{-k / 2+j / 2} & \text { if }-k \leqslant j \leqslant 0 \\
2^{-k / 2-3 j / 2} & \text { if } 0 \leqslant j\end{cases} \\
& \text { so } \sum_{j \in \mathbb{Z}} \sqrt{\alpha_{j}\left(m_{k}\right) \beta_{j}\left(m_{k}\right)} \lesssim 2^{-k / 2}
\end{aligned}
$$

Taking the supremum, we obtain $\Gamma_{B}\left(K_{1} *\left(K_{2}\right)_{\left(2^{k}\right)}\right) \leqslant C(\kappa) 2^{-|k| / 2}$, for $k \in \mathbb{Z}$.

Coming back to Carbery's situation, we obtain in this way by Lemma 5.14 that

$$
\left\|m_{k}\right\|_{2 \rightarrow 2} \leqslant \kappa \delta 2^{-|k| / 2}, \quad k \in \mathbb{Z}
$$

slightly better than what we got with $C_{\alpha}\left(m_{k}\right)$. Indeed, we must choose $\alpha>1 / 2$ with Carbery, and we have obtained for $C_{\alpha}\left(m_{k}\right)$ a bound of or$\operatorname{der} \delta 2^{-(1-\alpha)|k|}$.

### 6.5.1. A solution to the gap question

The gap question has been exposed in Section 6.4.1. Instead of the function studied precedently, equal to

$$
\widehat{N_{k}}(\xi): t \mapsto m_{[t]}(\xi)\left(\widehat{P}_{\left[t 2^{k}\right]}-\widehat{P}_{\left[t 2^{k+1}\right]}\right)(\xi), \quad t>0, \xi \in \mathbb{R}^{n}
$$

we need to study the family of multipliers defined by

$$
\widehat{n}_{k}(\xi, t)=m_{[t]}(\xi)\left(\widehat{P}_{\left[2^{j+k}\right]}-\widehat{P}_{\left[2^{j+k+1}\right]}\right)(\xi), \quad j \in \mathbb{Z} \text { and } 2^{j} \leqslant t \leqslant 2^{j+1}
$$

which are the Fourier transforms of the kernels $K_{(t)} *\left(P_{\left(2^{j(t)+k}\right)}-P_{\left(2^{j(t)+k+1)}\right.}\right)$ with $j(t)=\left\lfloor\log _{2} t\right\rfloor$. They do not fit into the setting of Proposition 6.14, but can be treated using Lemma 6.15. We do the following: for every $j \in \mathbb{Z}$, let $x_{j}=2^{j}+2^{j-1}$ be the midpoint of the interval $I_{j}=\left[2^{j}, 2^{j+1}\right]$. Let the "new" function be

$$
t \mapsto m_{\left[2^{j}+2\left(t-2^{j}\right)\right]}(\xi)\left(\widehat{P}_{\left[2^{j+k}\right]}-\widehat{P}_{\left[2^{j+k+1}\right]}\right)(\xi)
$$

for $t$ in the first half $\left[2^{j}, x_{j}\right]$ of the interval $I_{j}$, and

$$
t \mapsto m_{\left[2^{j+1}\right]}(\xi)\left(\widehat{P}_{\left[2^{k}\left(2^{j}+2\left(t-x_{j}\right)\right)\right]}-\widehat{P}_{\left[2^{k+1}\left(2^{j}+2\left(t-x_{j}\right)\right)\right]}\right)(\xi)
$$

in the second half. The first half "contains" the family $\widehat{n}_{k}(\xi, t)$ that we have to study, and adjoining the second half will allow us to exploit easily what has been done in Section 6.5 for the regular setting. We can describe more compactly the new setting if we define two motions going along $(0,+\infty)$ according to

$$
X(t)= \begin{cases}2^{j}+2\left(t-2^{j}\right), & 2^{j} \leqslant t \leqslant x_{j} \\ 2^{j+1}, & x_{j} \leqslant t \leqslant 2^{j+1}\end{cases}
$$

and

$$
Y(t)= \begin{cases}2^{j}, & 2^{j} \leqslant t \leqslant x_{j} \\ 2^{j}+2\left(t-x_{j}\right), & x_{j} \leqslant t \leqslant 2^{j+1}\end{cases}
$$

Then, the new function can be written as

$$
\begin{equation*}
\widetilde{m}_{k}(\xi, t):=m_{[X(t)]}(\xi)\left(\widehat{P}_{\left[2^{k} Y(t)\right]}-\widehat{P}_{\left[2^{k+1} Y(t)\right]}\right)(\xi) \tag{6.34}
\end{equation*}
$$

corresponding to the family of kernels $K_{t}=K_{(X(t))} *\left(P_{\left(2^{k} Y(t)\right)}-P_{\left(2^{k+1} Y(t)\right)}\right)$. The two functions $X, Y$ are non-decreasing, continuous, piecewise linear, and we have $X\left(2^{j}\right)=Y\left(2^{j}\right)=2^{j}$ for $j$ in $\mathbb{Z}$. Notice that $X(2 t)=2 X(t)$ and $Y(2 t)=2 Y(t)$ (make use of $2 x_{j}=x_{j+1}$ ). Also, $0 \leqslant X^{\prime}(t), Y^{\prime}(t) \leqslant 2$. Applying Remark 6.16, one sees easily that the functions $g_{\xi}(t)=\widetilde{m}_{k}(\xi, t) / t$ satisfy (6.26).

In the "dilation case" where $m_{0}(\xi, t)=m(t \xi)$, we have that $m_{0}(s \xi, t)=$ $m_{0}(\xi, s t)$ for every $s>0$, and it allowed us to restrict the study of the functions $t \mapsto m_{0}(\xi, t), \xi \in \mathbb{R}^{n}$, to the case $|\xi|=1$. This is not true anymore, but we still have that $m(2 \xi, t)=m(\xi, 2 t)$ for the two components $\Phi$ and $\Psi$ of $\widetilde{m}_{k}(\xi, t)$, defined by

$$
\Phi(\xi, t)=m_{[X(t)]}(\xi), \quad \Psi(\xi, t)=\left(\widehat{P}_{[Y(t)]}-\widehat{P}_{[Y(2 t)]}\right)(\xi)
$$

and this permits us to restrict to the case $1 \leqslant|\xi|<2$. Indeed,

$$
\Phi(2 \xi, t)=m_{[X(t)]}(2 \xi)=m(2 X(t) \xi)=m(X(2 t) \xi)=\Phi(\xi, 2 t)
$$

The same property holds true for $\Psi(\xi, t)$, with $Y$ replacing $X$.
Let us fix $\xi$ such that $1 \leqslant|\xi|<2$, and consider now

$$
\phi_{1}(u)=\Phi(\xi, u)=m(X(u) \xi), \quad \chi_{1}(u)=\Psi(\xi, u)=\mathrm{e}^{-2 \pi Y(u)|\xi|}-\mathrm{e}^{-4 \pi Y(u)|\xi|}
$$

Letting $\xi=|\xi| \theta$, we compare $\phi(u)=m(u \theta)$ with $\phi_{1}(u)=\phi(X(u)|\xi|)$. For every $u>0$, we have $u \leqslant X(u) \leqslant 2 u$ and $u / 2 \leqslant Y(u) \leqslant u$. We have therefore that $u \leqslant X(u)|\xi| \leqslant 4 u$ and $u / 2 \leqslant Y(u)|\xi| \leqslant 2 u$. Recall that $m$, difference of $m_{g}$ and $\widehat{P}$, satisfies (6.31a). It follows that

$$
\begin{aligned}
\delta^{-1}\left|\phi_{1}(u)\right| & =\left|\phi_{\delta}(X(u)|\xi|)\right| \leqslant 8\left[(X(u)|\xi|) \wedge\left(X(u)^{-1}|\xi|^{-1}\right)\right] \\
& \leqslant 32\left(|u| \wedge|u|^{-1}\right) \lesssim|u| \wedge|u|^{-1}
\end{aligned}
$$

We also have $\phi_{1}^{\prime}(u)=X^{\prime}(u) \phi^{\prime}(X(u)|\xi|)$, and since $X^{\prime}(u) \leqslant 2$,

$$
\delta^{-1}\left|\phi_{1}^{\prime}(u)\right| \leqslant 2\left|\phi_{\delta}^{\prime}(X(u)|\xi|)\right| \leqslant 16\left[1 \wedge\left(X(u)^{-1}|\xi|^{-1}\right)\right] \leqslant 16\left(1 \wedge|u|^{-1}\right)
$$

which can be written as $\delta^{-1}\left|\phi_{1}^{\prime}(u)\right| \lesssim 1 \wedge|u|^{-1}$. Using (6.31b), we have the same kind of inequalities for $\chi_{1}$. The proof in Section 6.5 depended only on these two bounds, so the result in (6.33) is also valid in the modified setting and gives the following lemma.

Lemma 6.19. - Suppose that $K_{g}$ is a probability density on $\mathbb{R}^{n}$ satisfying (6.1.H), that $m=m_{g}-\widehat{P}$ and that $\widetilde{m}_{k}$ is defined by (6.34). For $\alpha \in(1 / 2,1)$, one has

$$
\sup _{\xi \in \mathbb{R}^{n}}\left\|t \mapsto \widetilde{m}_{k}(\xi, t)\right\|_{L_{\alpha}^{2}} \leqslant \kappa\left(\delta_{0, g}+\delta_{1, g}\right) 2^{-(1-\alpha)|k|}, \quad k \in \mathbb{Z}
$$

### 6.6. Appendix: proof of Bourgain's $L^{2}$ theorem by Carbery's criterion

Proof. - This section is intended to illustrate the Fourier definition (6.7) of $D^{\alpha}$, and we shall have to perform some contortions in order to enter into the suitable setting. The kernel $K$ on $\mathbb{R}^{n}$ to which we want to apply the conclusion (1) of Carbery's Proposition 6.14 is again $K=K_{l c}-P$, as in Section 6.4, where $K_{l c}$ is a symmetric log-concave probability density on $\mathbb{R}^{n}$ normalized by variance. Let us fix a norm one vector $\theta \in \mathbb{R}^{n}$; here, the function $\varphi_{\theta}(s)=\int_{\theta \perp} K(y+s \theta) \mathrm{d}^{n-1} y$, for $s \in \mathbb{R}$, is the difference of two symmetric probability densities $\phi_{j}$, associated respectively to $K_{l c}$ and to the Poisson kernel $P$. The function $\phi_{1}$ of integrals of $K_{l c}$ on affine hyperplanes parallel to $\theta^{\perp}$ satisfies, according to Lemma 5.6, an estimate of exponential decay $\phi_{1}(s) \leqslant \kappa \mathrm{e}^{-|s| / \kappa}$, for $s \in \mathbb{R}$ and for a certain $\kappa>0$ universal. On the other hand, $\phi_{2}(s)$ is the Cauchy kernel (1.33.C) equal to $\pi^{-1}\left(1+s^{2}\right)^{-1}$, for which one has only $\phi_{2}(s) \leqslant 1 \wedge s^{-2}$, where $a \wedge b$ denotes the minimum of two real numbers $a$ and $b$. This estimate is valid also for $\phi_{1}$, up to some universal factor $\kappa$, and we shall remember for the absolute value of $\varphi_{\theta}$ that

$$
\begin{equation*}
\forall s \in \mathbb{R}, \quad\left|\varphi_{\theta}(s)\right| \leqslant \kappa\left(1 \wedge \frac{1}{s^{2}}\right) \tag{6.35}
\end{equation*}
$$

The Fourier transform $m$ of $K$ is given by

$$
m(t \theta)=\int_{\mathbb{R}} \varphi_{\theta}(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s
$$

Denote by $\Phi$ the antiderivative of $\varphi_{\theta}$ vanishing at 0 . The function $\Phi$ is odd, it vanishes also at infinity because $\varphi_{\theta}$ is even with integral zero. We deduce from (6.35), for some $\kappa>0$ and every $s \in \mathbb{R}$, that

$$
\begin{equation*}
|\Phi(s)| \leqslant \kappa\left(|s| \wedge|s|^{-1}\right) \tag{6.36}
\end{equation*}
$$

For $t \neq 0$, we could, performing an integration by parts, express $m(t \theta)$ by a simply converging integral

$$
m(t \theta)=2 \mathrm{i} \pi t \int_{-\infty}^{+\infty} \Phi(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s
$$

but we prefer to work with absolutely converging integrals, for example in this way: let us denote by $\widetilde{P}_{0}$ the $L^{1}$-normalized truncation $\widetilde{P}_{0}=$ $\left\|\mathbf{1}_{B} P\right\|_{L^{1}\left(\mathbb{R}^{n}\right)}^{-1} \mathbf{1}_{B} P$ of the Poisson kernel $P$ at a sufficiently large Euclidean ball $B$ in $\mathbb{R}^{n}$, so that $\left\|\mathbf{1}_{B} P\right\|_{1}>1 / 2$. We can see according to (1.35) that the radius of $B$ must be at least of order $\kappa \sqrt{n}$. Another possibility is to introduce a modified Poisson kernel

$$
\widetilde{P}(x)=2 P(x) \mathrm{e}^{-\varepsilon_{0}|x|^{2} / 2},
$$

where $\varepsilon_{0}>0$ is chosen so that the integral of $\widetilde{P}$ is equal to 1 . With both choices, one has $\widetilde{P}_{0}, \widetilde{P} \leqslant 2 P$, and the estimates of the maximal function for the kernel $P$ are thus clearly true for $\widetilde{P}$, with a bound simply doubled. For the same fixed $\theta$ of norm one, the modified function $\phi_{2}$ defined by

$$
\phi_{2}(s)=2 \int_{\theta^{\perp}} P(y+s \theta) \mathrm{e}^{-\varepsilon_{0}\left(|y|^{2}+s^{2}\right) / 2} \mathrm{~d}^{n-1} y \leqslant C(n) \mathrm{e}^{-\varepsilon_{0} s^{2} / 2}
$$

decays exponentially at infinity, and since $\phi_{2}(s) \leqslant 2 \pi^{-1}\left(1+s^{2}\right)^{-1}$, the modified function $\phi_{2}$ satisfies (6.35) and (6.36). The modified antiderivative $\Phi$ inherits now at infinity of the exponential decay of $\phi_{1}$ and of $\phi_{2}$, and this makes the integrals that follow absolutely convergent. However, the "universal" estimates remain given by (6.35) and (6.36).

The situation would be simpler using a Gaussian kernel, letting

$$
K(x)=K_{C}(x)-G(x), \quad x \in \mathbb{R}^{n}
$$

with $G$ being the $N\left(0, \mathrm{I}_{n}\right)$ density (1.17) on $\mathbb{R}^{n}$.
We apply here the Fourier definition (6.7) for $D^{\alpha}$. For every $t>0$ we write

$$
\frac{m(t \theta)}{t}=2 \mathrm{i} \pi \int_{\mathbb{R}} \Phi(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s
$$

where $|\Phi|$ decays exponentially at infinity. This ensures that $t \mapsto m(t \theta) / t$ is $C^{\infty}$ on the line, with bounded derivatives. By (6.7), we can express the fractional derivative appearing in Carbery's criterion as

$$
D_{t}^{\alpha}\left(\frac{m(t \theta)}{t}\right)=2 \mathrm{i} \pi(2 \pi)^{\alpha} \int_{\mathbb{R}}(\mathrm{i} s)^{\alpha} \Phi(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s
$$

For $0<\alpha<1$, we write

$$
\int_{0}^{\infty} s^{\alpha} \Phi(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s=\frac{1}{2 \mathrm{i} \pi t} \int_{0}^{\infty}\left(s^{\alpha} \Phi(s)\right)^{\prime} \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s
$$

and because $\left(s^{\alpha} \Phi(s)\right)^{\prime}$ vanishes at 0 , we see that

$$
\int_{0}^{\infty} s^{\alpha} \Phi(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s=-\frac{1}{4 \pi t^{2}} \int_{0}^{\infty}\left(s^{\alpha} \Phi(s)\right)^{\prime \prime} \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s
$$

The integrals on the side of negative $s$ ask for an analogous treatment, essentially already seen in Section 5.2, Lemma 5.8. We estimate the various parts (five parts) issued from the differentiations of $s^{\alpha} \Phi(s)$ to the first and second order, by applying the upper bounds (6.35) and (6.36) and the fact that $0<\alpha<1$. Notice that

$$
\int_{0}^{\infty}\left(s^{\alpha-1}+s^{\alpha-2}\right)\left(s \wedge s^{-1}\right) \mathrm{d} s=\frac{1}{1+\alpha}+\frac{1}{\alpha}+\frac{1}{1-\alpha}+\frac{1}{2-\alpha}=: \kappa_{\alpha} .
$$

Grouping two of the terms issued from $\left(s^{\alpha} \Phi\right)^{\prime},\left(s^{\alpha} \Phi\right)^{\prime \prime}$ and using (6.36), we have

$$
\left|\int_{0}^{\infty} s^{\alpha-1} \Phi(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s\right|+\left|\int_{0}^{\infty} s^{\alpha-2} \Phi(s) \mathrm{e}^{-2 \mathrm{i} \pi s t} \mathrm{~d} s\right| \leqslant \kappa \kappa_{\alpha}
$$

we also have $\int_{0}^{\infty}\left(s^{\alpha}+s^{\alpha-1}\right)\left|\varphi_{\theta}(s)\right| \mathrm{d} s \leqslant \kappa \kappa_{\alpha}$ for two other terms by (6.35), and finally for each $\phi=\phi_{j}, j=1,2$, decreasing on the positive side of the real line, we know by Lemma 5.9 that

$$
\int_{0}^{\infty} s^{\alpha}\left|\phi^{\prime}(s)\right| \mathrm{d} s=\alpha \int_{0}^{+\infty} s^{\alpha-1} \phi(s) \mathrm{d} s<+\infty
$$

which permits us to close this list of estimates for $\varphi_{\theta}=\phi_{1}-\phi_{2}$. It follows that for every $t>0$, we have

$$
\left|D_{t}^{\alpha}\left(\frac{m(t \theta)}{t}\right)\right| \leqslant \kappa_{\alpha}^{\prime}\left(t^{-1} \wedge t^{-2}\right)
$$

with $\kappa_{\alpha}^{\prime} \leqslant \kappa^{\prime}(2 \pi)^{\alpha} \kappa_{\alpha}$ independent of the direction $\theta$. Recalling the definition (6.24) and since $0<\alpha<1$, we get

$$
\begin{aligned}
C_{\alpha}(m)^{2} & =\sup _{\theta \in S^{n-1}}\|t \mapsto m(t \theta)\|_{L_{\alpha}^{2}}^{2}=\sup _{\theta \in S^{n-1}} \int_{0}^{+\infty}\left|t^{\alpha+1} D_{t}^{\alpha}\left(\frac{m(t \theta)}{t}\right)\right|^{2} \frac{\mathrm{~d} t}{t} \\
& \leqslant\left(\kappa_{\alpha}^{\prime}\right)^{2}\left(\int_{0}^{1}\left(t^{\alpha+1} t^{-1}\right)^{2} \frac{\mathrm{~d} t}{t}+\int_{1}^{+\infty}\left(t^{\alpha+1} t^{-2}\right)^{2} \frac{\mathrm{~d} t}{t}\right) \\
& =\left(\kappa_{\alpha}^{\prime}\right)^{2}\left(\int_{0}^{1} t^{2 \alpha-1} \mathrm{~d} t+\int_{1}^{+\infty} t^{2 \alpha-3} \mathrm{~d} t\right) \\
& =\left(\kappa_{\alpha}^{\prime}\right)^{2}\left(\frac{1}{2 \alpha}+\frac{1}{2-2 \alpha}\right)<+\infty
\end{aligned}
$$

One thus chooses $\alpha \in(1 / 2,1)$ arbitrary and applies Carbery's Proposition $6.14(1)$, which gives the boundedness on $L^{2}\left(\mathbb{R}^{n}\right)$ of the maximal operator associated to the difference kernel $K=K_{l c}-\widetilde{P}$. We get in this way that the maximal operator $\mathrm{M}_{K_{l c}}$ is bounded on $L^{2}\left(\mathbb{R}^{n}\right)$ by a constant independent of the dimension $n$.

## 7. The Detlef Müller article

Müller [59] introduces a geometrical parameter $Q(C)$ associated to every symmetric convex body $C$ in $\mathbb{R}^{n}$. When $C$ is isotropic of volume 1 , this parameter $Q(C)$ is equal to the maximum of $(n-1)$-dimensional volumes of hyperplane projections of $C$. Müller shows that in the class $\mathcal{C}(\lambda)$ consisting of $C$ s for which $Q(C)$ and the isotropy constant $L(C)$ are bounded by a given $\lambda$, the existence for the maximal operator $\mathrm{M}_{C}$ associated to $C$ of an $L^{p}\left(\mathbb{R}^{n}\right)$ bound, uniform in $n$, can be pushed to every value $p>1$ with a constant $\kappa(p, \lambda)$ depending on $p$ and $\lambda$ only. This removes - in a way the restriction $p>3 / 2$ imposed by Bourgain and Carbery.

We have seen in (5.1) and (5.3) that when $C_{0}$ is isotropic of volume 1 in $\mathbb{R}^{n}$, then the dilate $C_{1}=r_{0} C_{0}$ with $r_{0}=L\left(C_{0}\right)^{-1}$ is isotropic and normalized by variance. The proof of Müller will actually make use of a parameter $q\left(C_{1}\right)$ equal to the supremum in $\theta \in S^{n-1}$ of the masses of the signed measures $\theta$. $\nabla K_{C_{1}}$. We shall see that for $\theta$ of norm one, the mass of the measure $\theta \cdot \nabla K_{C_{1}}$, the directional derivative in the sense of distributions of the probability measure $\mu_{C_{1}}$, is given by

$$
\frac{2\left|P_{\theta} C_{1}\right|_{n-1}}{\left|C_{1}\right|_{n}}=2 r_{0}^{-n} r_{0}^{n-1}\left|P_{\theta} C_{0}\right|_{n-1} \leqslant \frac{2}{r_{0}} Q\left(C_{0}\right)=2 L\left(C_{0}\right) Q\left(C_{0}\right)
$$

where $P_{\theta}$ is the orthogonal projection onto the hyperplane $\theta^{\perp}$. For every symmetric convex set $C$, we let $C_{0}$ be an isotropic position of volume 1 for $C$ and we set

$$
\begin{equation*}
q(C)=2 L\left(C_{0}\right) Q\left(C_{0}\right) . \tag{7.1}
\end{equation*}
$$

Müller [59, Section 3] proves that $q(C)$ is uniformly bounded for the family of unit balls $B_{n}^{q}$ of $\ell_{n}^{q}, 1 \leqslant q<+\infty$ fixed and $n \in \mathbb{N}^{*}$. This is easy when $q=2$. By (5.4), we know that the Euclidean ball $B_{n, V}$ in $\mathbb{R}^{n}$ normalized by variance has a radius $r_{n, V}$ equal to $\sqrt{n+2}$, hence by the log-convexity of the Gamma function we get

$$
\begin{aligned}
q\left(B_{n}^{2}\right) & =\sup _{\theta \in S^{n-1}} \frac{2\left|P_{\theta} B_{n, V}\right|_{n-1}}{\left|B_{n, V}\right|_{n}}=\frac{2 \omega_{n-1}}{r_{n, V} \omega_{n}}=\frac{2 \Gamma(n / 2+1)}{\sqrt{\pi(n+2)} \Gamma(n / 2+1 / 2)} \\
& \leqslant \frac{2 \Gamma(n / 2+1 / 2)^{1 / 2} \Gamma(n / 2+3 / 2)^{1 / 2}}{\sqrt{\pi(n+2)} \Gamma(n / 2+1 / 2)}=2 \sqrt{\frac{n+1}{2 \pi(n+2)}}<\sqrt{\frac{2}{\pi}}
\end{aligned}
$$

Given a kernel $K$ integrable on $\mathbb{R}^{n}$ and having partial derivatives $\partial_{j} K$ in the sense of distributions that are (signed) measures $\mu_{j}$, for $j=1, \ldots, n$, we define the directional variation $V(K)$ of $K$ by

$$
\begin{equation*}
V(K):=\sup _{\theta \in S^{n-1}}\|\theta \cdot \nabla K\|_{1}=\sup _{\theta \in S^{n-1}}\left\|\sum_{j=1}^{n} \theta_{j} \mu_{j}\right\|_{1} \tag{7.2}
\end{equation*}
$$

We will show at Lemma 7.10 that $V\left(K_{C}\right)=q(C)$ when $C$ is an isotropic symmetric convex body normalized by variance. For the $N\left(0, \mathrm{I}_{n}\right)$ Gaussian density $\gamma_{n}$, we see that $V\left(\gamma_{n}\right)=\int_{\mathbb{R}^{n}}\left|x \cdot \mathbf{e}_{1}\right| \mathrm{d} \gamma_{n}(x)=\int_{\mathbb{R}}|u| \mathrm{d} \gamma_{1}(u)=\sqrt{2 / \pi}$. Notice that

$$
\begin{equation*}
V\left(K_{(t)}\right)=t^{-1} V(K), \quad t>0, \quad \text { and } \quad V(K * \mu) \leqslant V(K) \tag{7.3}
\end{equation*}
$$

for any probability measure $\mu$ on $\mathbb{R}^{n}$. Since $V\left(\gamma_{n}\right)$ is independent of $n$, it follows from the subordination formula (1.30) that the same is true for the Poisson kernel $P_{1}^{(n)}$ on $\mathbb{R}^{n}$ expressed in (1.32). Precisely, because $G_{s}$ in (1.30) is a $N\left(0, s \mathrm{I}_{n}\right)$ Gaussian measure, we have $V\left(G_{s}\right)=s^{-1 / 2} V\left(\gamma_{n}\right)$ by (7.3) and we first get

$$
\begin{equation*}
V\left(P_{1}^{(n)}\right) \leqslant \int_{0}^{+\infty} V\left(G_{s}\right) \frac{s^{-3 / 2}}{\sqrt{2 \pi}} \mathrm{e}^{-1 /(2 s)} \mathrm{d} s=\int_{0}^{+\infty} \frac{\mathrm{e}^{-1 /(2 s)}}{\pi} \frac{\mathrm{d} s}{s^{2}}=\frac{2}{\pi} \tag{7.4}
\end{equation*}
$$

but actually $V\left(P_{1}^{(n)}\right)=2 / \pi$ since for each $x \in \mathbb{R}^{n}$, all gradients $\nabla G_{s}(x)$, $s>0$, are nonnegative multiples of the same vector $-x$. This equality is of course also easy to derive by a direct calculation on the Poisson density.

Besides the appearance of the parameter $q(C)$, Müller's proof draws on estimates such as (6.1.H), but extended to more derivatives of the Fourier transform $m_{C}$ of $K_{C}$. That bounding more derivatives leads to improved results was already seen in Bourgain [11], who obtained a dimension free bound in $L^{p}\left(\mathbb{R}^{n}\right)$ for all $p>1$ in the case of the maximal operator $\mathrm{M}_{C}$ of $\ell_{n}^{q}$ balls when $q$ is an even integer. We shall consider a probability density $K_{g}$ on $\mathbb{R}^{n}$ or more generally an integrable kernel $K_{g}$, with a Fourier transform $m_{g}$ satisfying that for every integer $j \geqslant 0$, there exists a constant $\delta_{j, g}$ such that

$$
\left|\frac{\mathrm{d}^{j}}{\mathrm{~d} t^{j}} m_{g}(t \theta)\right| \leqslant \frac{\delta_{j, g}}{1+t}, \quad \theta \in S^{n-1}, \quad t>0
$$

Actually, for each specific value $p \in(1,3 / 2]$, bounding $\mathrm{M}_{C}$ in $L^{p}\left(\mathbb{R}^{n}\right)$, knowing that $q(C) \leqslant \lambda$, requires a certain finite number of estimates from the infinite list $\left(7.5 . \mathbf{H}_{\infty}\right)$, and this number increases to infinity when $p$ tends to 1. We let

$$
\begin{equation*}
\Delta_{k}=\sum_{j=0}^{k} \delta_{j, g} \tag{7.6}
\end{equation*}
$$

The "radial" estimate $\left(7.5 . \mathbf{H}_{\infty}\right)$ implies $\left|\mathrm{d}^{j} /\left(\mathrm{d} t^{j}\right) m_{g}(t \xi)\right| \leqslant \delta_{j, g}|\xi|^{j} /(1+|t \xi|)$ for $\xi \neq 0$. It is natural to disregard $\xi=0$ in a radial method, but when $j>0$, we can extend continuously $\xi \mapsto \mathrm{d}^{j} /\left(\mathrm{d} t^{j}\right) m_{g}(t \xi)$ by giving the value 0 at $\xi=0$.

Theorem 7.1 (Müller [59]). - For every $p \in(1,+\infty]$ and $\lambda>0$, there exists a constant $\kappa(p, \lambda)$ independent of $n$ such that

$$
\left\|\mathrm{M}_{K_{l c}} f\right\|_{L^{p}\left(\mathbb{R}^{n}\right)} \leqslant \kappa(p, \lambda)\|f\|_{L^{p}\left(\mathbb{R}^{n}\right)}
$$

if $K_{l c}$ is an isotropic symmetric log-concave probability density on $\mathbb{R}^{n}$, normalized by variance and with $V\left(K_{l c}\right) \leqslant \lambda$. In particular, for every symmetric convex body $C$ in $\mathbb{R}^{n}$ such that $q(C) \leqslant \lambda$, one has $\left\|\mathrm{M}_{C} f\right\|_{L^{p}\left(\mathbb{R}^{n}\right)} \leqslant$ $\kappa(p, \lambda)\|f\|_{L^{p}\left(\mathbb{R}^{n}\right)}$. When $p \in(1,2]$, we can write more precisely

$$
\left\|\mathrm{M}_{K_{l c}} f\right\|_{L^{p}\left(\mathbb{R}^{n}\right)} \leqslant \kappa(p)\left(1+\lambda^{2 / p-1}\right)
$$

If a probability density $K_{g}$ satisfies $\left(7.5 . \mathbf{H}_{\infty}\right)$ and if $p \in(1,2]$, then we have $\left\|\mathrm{M}_{K_{g}} f\right\|_{L^{p}\left(\mathbb{R}^{n}\right)} \leqslant \kappa_{p} \Delta_{k_{0}(p)}^{1-1 / p} \Delta_{1}^{1-1 / p}\left(1+V\left(K_{g}\right)^{2 / p-1}\right)$, with $k_{0}(p)<p /(p-1)$.

The subsequent proof furnishes for the constant $\kappa_{p}$ in the line above an order exponential in $q=p /(p-1)$ that is certainly not right, see Remarks 7.13 and 7.14. The case $p>3 / 2$ is already known, with $\kappa(p, \lambda)$ independent of $\lambda$, see Theorem 6.2 and Proposition 6.3. We know by Lemma 5.11 that isotropic symmetric log-concave probability densities satisfy (7.5.H ${ }_{\infty}$ ) with absolute constants $\left(\delta_{j, c}\right)_{j=0}^{\infty}$. We shall thus concentrate on the $K_{g}$ case and on values $p \in(1,3 / 2]$. Taking Carbery's results into account, the following proposition will be (essentially) enough for proving Müller's theorem.

Proposition 7.2 ([59, Proposition 1]). - Let $K_{g}$ be an integrable kernel on $\mathbb{R}^{n}$ satisfying $\left(7.5 . \mathbf{H}_{\infty}\right)$ and let $m_{g}$ be its Fourier transform. For every $\alpha \in(0,1)$ and every $p \in(1,+\infty)$, the multiplier $(\xi \cdot \nabla)^{\alpha} m_{g}(\xi)$ in $\left(6.18 . \nabla^{\alpha}\right)$ admits on $L^{p}\left(\mathbb{R}^{n}\right)$ a bound that depends upon $p, \alpha, \mathbf{d}=\left(\delta_{j, g}\right)_{j=0}^{\infty}$ and $V\left(K_{g}\right)$, but not on the dimension $n$. When $p \in(1,2]$ and if $\left\|K_{g}\right\|_{L^{1}\left(\mathbb{R}^{n}\right)} \leqslant 1$, we can write
$\left\|(\xi \cdot \nabla)^{\alpha} m_{g}(\xi)\right\|_{p \rightarrow p} \leqslant 1+\kappa(\alpha, p) \Delta_{k(p)}^{(4 / 3)(1-1 / p)}\left(1+\delta_{0, g}^{(2 / 3)(1-1 / p)} V\left(K_{g}\right)^{2 / p-1}\right)$, with $k(p)=\lceil 3 p /(4 p-4)\rceil$.

The case $p=2$ follows easily from Parseval (2.12.P) by (6.1.H) and (6.20). The result for $p \geqslant 2$ can be obtained by duality from the case $1<p \leqslant 2$.

Proof of Theorem 7.1. - Let $p \in(1,2)$ be given. We then choose $p_{0} \in$ $(1, p)$ and $\alpha \in\left(1 / p_{0}, 1\right)$ as being functions of $p$, for example $p_{0}=(2 p+$ $2) /(5-p)$ and $\alpha=(p+7) /(4 p+4)$. We apply in $L^{p_{0}}\left(\mathbb{R}^{n}\right)$ the part (2) of Proposition 6.14 to the kernel $K=K_{g}-P$. We know by Proposition 7.2
that $(\xi \cdot \nabla)^{\alpha} m_{g}(\xi)$ is bounded on $L^{p_{0}}\left(\mathbb{R}^{n}\right)$ by a function of $V\left(K_{g}\right)$ and we will check in Section 7.2 that $(\xi \cdot \nabla)^{\alpha} \widehat{P}(\xi)$ is also bounded on $L^{p_{0}}\left(\mathbb{R}^{n}\right)$ by some $\pi_{\alpha, p_{0}}=\pi(p)$. It follows for $m=m_{g}-\widehat{P}$ that
$\left\|(\xi \cdot \nabla)^{\alpha} m(\xi)\right\|_{p_{0} \rightarrow p_{0}} \leqslant \kappa_{0}(p, \mathbf{d})\left(1+V\left(K_{g}\right)^{2 / p_{0}-1}\right) \leqslant \kappa_{0}(p, \mathbf{d})\left(1+\lambda^{2 / p_{0}-1}\right)$,
with $\kappa_{0}(p, \mathbf{d}) \leqslant \kappa(p) \Delta_{k\left(p_{0}\right)}^{(4 / 3)\left(1-1 / p_{0}\right)} \delta_{0, g}^{(2 / 3)\left(1-1 / p_{0}\right)}$, where $\Delta_{j} \geqslant \delta_{0, g} \geqslant 1$ because $K_{g}$ here is a probability density. We obtain in this way that

$$
f \mapsto W_{1} f:=\sup _{1 \leqslant u \leqslant 2}\left|K_{(u)} * f\right|
$$

is bounded on $L^{p_{0}}\left(\mathbb{R}^{n}\right)$. This was the only missing information for deducing from Proposition 6.6 that $M_{K}$ is bounded on $L^{p}\left(\mathbb{R}^{n}\right)$ when $1<p \leqslant 3 / 2$. Indeed, with the notation of Section 6.1, let $T_{j, v}$ be the convolution with $K_{\left(2^{j} v\right)}, v \in[1,2]$ and let $T_{j}$ be as in (6.2). By Proposition $6.14(2)$, we have for every $j \in \mathbb{Z}$ that

$$
\begin{align*}
\left\|T_{j}\right\|_{p_{0} \rightarrow p_{0}}=\left\|T_{0}\right\|_{p_{0} \rightarrow p_{0}} & =\left\|W_{1}\right\|_{p_{0} \rightarrow p_{0}} \\
& \leqslant \kappa_{\alpha, p_{0}}\left(2+\left\|(\xi \cdot \nabla)^{\alpha} m(\xi)\right\|_{p_{0} \rightarrow p_{0}}\right) \tag{7.7}
\end{align*}
$$

with $\kappa_{\alpha, p_{0}}$ from (6.25). We bound it by $C_{p_{0}}^{\prime \prime}(\lambda):=\kappa_{\alpha, p_{0}}\left(2+\kappa_{0}(p, \mathbf{d})(1+\right.$ $\left.\lambda^{2 / p_{0}-1}\right)$ ). By (6.4), with $p_{0}$ already set and $r_{0}=2 p /\left(p+2-p_{0}\right)$ function of $p$ and $p_{0}$, we get

$$
\begin{equation*}
\left\|M_{K}\right\|_{p \rightarrow p} \leqslant\left(C_{r_{0}}\right)^{2 \gamma / p_{0}} C_{p_{0}}^{\prime \prime}(\lambda)^{\gamma}\left(\sum_{k \in \mathbb{Z}} a_{k}^{(1-\gamma) p / 2}\right)^{2 / p}+2 C_{p}^{\prime} \tag{7.8}
\end{equation*}
$$

where $\gamma=[1 / p-1 / 2] /\left[1 / p_{0}-1 / 2\right]=(p+1) /(2 p)$. The constants $C_{r_{0}}$ in $\left(A_{\mathbf{0}}\right)$, $C_{p}^{\prime}$ in $\left(A_{\mathbf{1}}\right)$ depend only on $p, p_{0}$ and $r_{0}$, hence on $p$ alone, and they exist regardless of $p>3 / 2$ or not. By Section 6.5, we know that under (6.1.H), the $\left(a_{k}\right)_{k \in \mathbb{Z}}$ in $\left(A_{\mathbf{3}}\right)$ satisfy $a_{k} \leqslant\left(\delta_{0, g}+\delta_{1, g}\right) a_{\alpha, k}$ with $\left(a_{\alpha, k}\right)_{k \in \mathbb{Z}}$ universal. We obtain

$$
\begin{aligned}
\left\|\mathrm{M}_{K_{g}}\right\|_{p \rightarrow p} & \leqslant\left\|M_{K}\right\|_{p \rightarrow p}+\kappa_{p} \\
& \leqslant \kappa(p, \mathbf{d})\left(1+\lambda^{2\left(1 / p_{0}-1 / 2\right) \gamma}\right)=\kappa(p, \mathbf{d})\left(1+\lambda^{2 / p-1}\right)
\end{aligned}
$$

with $1-1 / p_{0}=(3 p-3) /(2 p+2), k\left(p_{0}\right)=\lceil(p+1) /(2 p-2)\rceil<p /(p-1)$, and

$$
\begin{aligned}
\kappa(p, \mathbf{d}) & \leqslant \kappa(p)\left(\Delta_{k\left(p_{0}\right)\left(1-1 / p_{0}\right)}^{(4 / 3)} \delta_{0, g}^{(2 / 3)\left(1-1 / p_{0}\right)}\right)^{\gamma} \Delta_{1}^{1-\gamma} \\
& \leqslant \kappa(p) \Delta_{k_{0}(p)}^{1-1 / p} \Delta_{1}^{1-1 / p}
\end{aligned}
$$

### 7.1. The Müller strategy

Müller prefers to work with another version $i^{w}$ of the fractional integral $I^{w}$ from (6.9). This version is defined when $\operatorname{Re} w>0$, beginning this time with $f \in C^{\infty}(\mathbb{R})$, by the formula

$$
\left(i^{w} f\right)(t)=\frac{1}{\Gamma(w)} \int_{t}^{2}(u-t)^{w-1} f(u) \mathrm{d} u, \quad t \leqslant 2
$$

The chosen limit 2 is rather arbitrary, but will be quite convenient for the computations that follow, in particular because $(2-1)^{w}=1$ for every $w$. Integrating by parts as we did for $I^{w}$ in Section 6.2, we get

$$
\left(i^{w} f\right)(t)=\frac{(2-t)^{w} f(2)}{\Gamma(w+1)}-\frac{1}{\Gamma(w+1)} \int_{t}^{2}(u-t)^{w} f^{\prime}(u) \mathrm{d} u
$$

This new formula makes sense for $\operatorname{Re} w>-1$ and defines a fractional derivative $d^{z}$ if $z=-w$ and $\operatorname{Re} z<1$, by setting

$$
\begin{equation*}
\left(d^{z} f\right)(t)=\frac{(2-t)^{-z} f(2)}{\Gamma(1-z)}-\frac{1}{\Gamma(1-z)} \int_{t}^{2}(u-t)^{-z} f^{\prime}(u) \mathrm{d} u, \quad t \leqslant 2 \tag{7.9}
\end{equation*}
$$

Notice that $\left(d^{0} f\right)(t)=f(2)-\int_{t}^{2} f^{\prime}(u) \mathrm{d} u=f(t)$. Continuing integration by parts as in Section 6.2, we get successive formulas defining $d^{z} f$, for each integer $k$, which make sense for $\operatorname{Re} z<k$ and extend each other. Gluing them together, we can define entire functions of $z$ for every $t$ fixed and every given function $f \in C^{\infty}(\mathbb{R})$, for example $\left(d^{z} \mathbf{1}\right)(1)=1 / \Gamma(1-z)$ if $f=\mathbf{1}$. Suppose that $\operatorname{Re} z<0$. From

$$
\left(d^{z} f\right)(t)=\frac{1}{\Gamma(-z)} \int_{t}^{2}(u-t)^{-z-1} f(u) \mathrm{d} u
$$

we get for every integer $k \geqslant 1$ that

$$
\begin{equation*}
\left(d^{z} f\right)(t)=E_{k}(z, t)+(-1)^{k} \frac{1}{\Gamma(k-z)} \int_{t}^{2}(u-t)^{-z+k-1} f^{(k)}(u) \mathrm{d} u \tag{7.10}
\end{equation*}
$$

a formula to be compared with (6.16), and where $E_{k}(z, t)$ is equal to

$$
E_{k}(z, t)=\sum_{j=0}^{k-1}(-1)^{j} \frac{(2-t)^{-z+j} f^{(j)}(2)}{\Gamma(j+1-z)}
$$

If $z$ is in $\mathbb{C}, t \leqslant 2$ and $\operatorname{Re} z<k$, we can take (7.10) as definition for $\left(d^{z} f\right)(t)$.
When $-1<\operatorname{Re} z<0, f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$ and $t<2$, we see that

$$
\begin{align*}
\left(D^{z} f\right)(t)-\left(d^{z} f\right)(t) & =\left(\left[I^{-z}-i^{-z}\right] f\right)(t) \\
& =\frac{1}{\Gamma(-z)} \int_{2}^{+\infty}(u-t)^{-z-1} f(u) \mathrm{d} u \tag{7.11}
\end{align*}
$$

This equality can be extended by analytic continuation to every $z \in \mathbb{C}$ with $\operatorname{Re} z>-1$, or it can be proved by successive integrations by parts. In particular, one has $\left(d^{N} f\right)(t)=\left(D^{N} f\right)(t)=(-1)^{N} f^{(N)}(t)$ for every integer $N \geqslant 0$ because $\Gamma(-N)^{-1}=0$. As we did for $D^{\alpha}$, when the function of $t$ does not have an explicit name, we use the notation $d_{t}^{\alpha} f(2 t)$, and $\left.d_{t}^{\alpha} f(2 t)\right|_{t=1}$ for the value at $t=1$.

Lemma 7.3 (Müller [59]). - Let $m$ denote the Fourier transform of a kernel $K$ integrable on $\mathbb{R}^{n}$. For every $\alpha \in(0,1)$, the difference

$$
(\xi \cdot \nabla)^{\alpha} m(\xi)-\left.d_{t}^{\alpha} m(t \xi)\right|_{t=1}, \quad \xi \in \mathbb{R}^{n}
$$

is a multiplier on $L^{p}\left(\mathbb{R}^{n}\right), 1 \leqslant p \leqslant+\infty$, with a norm bounded by $\|K\|_{L^{1}\left(\mathbb{R}^{n}\right)}$.
Proof. - By $\left(6.18 . \nabla^{\alpha}\right)$ we have $(\xi \cdot \nabla)^{\alpha} m(\xi)=\left.D_{t}^{\alpha} m(t \xi)\right|_{t=1}$. From (7.11), we get

$$
(\xi \cdot \nabla)^{\alpha} m(\xi)-\left.d_{t}^{\alpha} m(t \xi)\right|_{t=1}=\frac{1}{\Gamma(-\alpha)} \int_{2}^{+\infty}(u-1)^{-\alpha-1} m(u \xi) \mathrm{d} u
$$

The result follows by Lemma 2.1, since

$$
\frac{1}{|\Gamma(-\alpha)|} \int_{2}^{+\infty}(u-1)^{-\alpha-1} \mathrm{~d} u=\frac{1}{|-\alpha \Gamma(-\alpha)|}=\frac{1}{\Gamma(1-\alpha)}<1
$$

Thanks to the reduction from $(\xi \cdot \nabla)^{\alpha} m(\xi)$ to $\left.d_{t}^{\alpha} m(t \xi)\right|_{t=1}$ given by Lemma 7.3, one can transform the condition (2) of Proposition 6.14. The objective now is to control the action on $L^{p}\left(\mathbb{R}^{n}\right)$ of the multiplier $\left.d_{t}^{\alpha} m_{g}(t \xi)\right|_{t=1}$, for some fixed $\alpha \in(1 / p, 1)$ denoted by $\alpha=1-\varepsilon$, where $\varepsilon>0$ gets arbitrarily small when $p$ tends to 1 . Müller embeds the "objective" into the holomorphic family of multipliers

$$
\begin{equation*}
m_{z}^{\varepsilon}(\xi)=\left.(1+|\xi|)^{1-\varepsilon-z} d_{t}^{z} m_{g}(t \xi)\right|_{t=1}, \quad \operatorname{Re} z>-1 \tag{7.12}
\end{equation*}
$$

and applies the complex interpolation scheme described in Section 3.2. For the value $z=\alpha=1-\varepsilon$, one has

$$
m_{\alpha}^{\varepsilon}(\xi)=m_{1-\varepsilon}^{\varepsilon}(\xi)=\left.d_{t}^{1-\varepsilon} m_{g}(t \xi)\right|_{t=1}=\left.d_{t}^{\alpha} m_{g}(t \xi)\right|_{t=1}
$$

which is the objective to be controlled. Müller studies this holomorphic family for $z \in \mathbb{C}$ varying in a strip of the form $-\varepsilon \leqslant \operatorname{Re} z \leqslant \nu$, with $\nu>0$ real. He shows by rather long and delicate calculations that the multipliers $m_{z}^{\varepsilon}(\xi)$ are bounded functions of $\xi \in \mathbb{R}^{n}$, for all $z$ in this strip, not uniformly in $z$, but with a $L^{\infty}\left(\mathbb{R}^{n}\right)$ norm of order $\Gamma(z)^{-1}$. This allows him to control the action on $L^{2}\left(\mathbb{R}^{n}\right)$, which is used for one end of the interpolation scale, the one corresponding to $\operatorname{Re} z=\nu$.

The other end of the scale is $\operatorname{Re} z=-\varepsilon$, where the operator associated to

$$
\begin{aligned}
m_{-\varepsilon+\mathrm{i} \tau}^{\varepsilon} & =\left.(1+|\xi|)^{1-\mathrm{i} \tau} d_{t}^{-\varepsilon+\mathrm{i} \tau} m_{g}(t \xi)\right|_{t=1} \\
& =\left.(1+|\xi|)^{-\mathrm{i} \tau}(1+|\xi|) d_{t}^{-\varepsilon+\mathrm{i} \tau} m_{g}(t \xi)\right|_{t=1}
\end{aligned}
$$

involves a "small" fractional integration $d^{-\varepsilon+\mathrm{i} \tau}$ of order $\varepsilon$, and a multiplication on the Fourier side by $1+|\xi|$. We will show that these multipliers $m_{-\varepsilon+\mathrm{i} \tau}^{\varepsilon}$ are bounded on all the spaces $L^{r}\left(\mathbb{R}^{n}\right), 1<r<+\infty$. In order to do it, we shall have to work mainly on the multiplier $|\xi| m_{g}(\xi)$. The parameter $V\left(K_{g}\right)$ appears when bounding the action of this multiplier on $L^{r}\left(\mathbb{R}^{n}\right)$, and the proof will use the dimensionless estimates for the Riesz transforms given in (2.22). Next, given $p$ in $(1,2]$, we choose $p_{0} \in(1, p), \alpha \in\left(1 / p_{0}, 1\right)$, and $\nu>\alpha$ which is a function of $p, p_{0}, \alpha$. By interpolation between $L^{2}\left(\mathbb{R}^{n}\right)$ (when $\operatorname{Re} z=\nu$ ) and $L^{p_{0}}\left(\mathbb{R}^{n}\right)$ (when $\operatorname{Re} z=-\varepsilon$ ), we shall obtain for the value $\alpha=1-\varepsilon$ the boundedness on $L^{p}\left(\mathbb{R}^{n}\right)$ of the multiplier $m_{\alpha}^{\varepsilon}(\xi)$ that is our "objective", thus proving Proposition 7.2.

Let us comment on the formulas for the Müller multipliers. We know by (5.19) in Corollary 5.13 that differentiating $N$ times the function $t \mapsto$ $m_{g}(t \xi)$ introduces a factor of order $(1+|\xi|)^{N-1}$, which must be compensated for being in a position to apply Parseval for the $L^{2}$ bound, using (2.12.P) as usual. This is done by multiplying by $(1+|\xi|)^{1-\varepsilon-\nu}$ when $z=\nu$. On the other hand, we do not want a compensating factor when $z=\alpha$, where we want to precisely recover our objective. The compensation will thus be of the form $(1+|\xi|)^{a z+b}$, with $a \nu+b=1-\varepsilon-\nu$ and $a \alpha+b=0$. We then get a "compensating factor" with a positive power of $|\xi|$ for $\operatorname{Re} z<\alpha$, which becomes an additional problem and requires more work.

The interpolation strip technique has been often employed by Stein. For example, in [73, Chap. III, §3], for studying the maximal function $\sup _{t>0}\left|P_{t} f\right|$ of general semi-groups, Stein works on a strip $S$ of the form $-1 \leqslant \operatorname{Re} z \leqslant$ $N$. If $z=-1$, he considers that the maximal inequality of Hopf concerns the derivative of order -1 of the semi-group, that is to say, its antiderivative (multiplied by $t^{z}=t^{-1}$ )

$$
t^{-1} D_{t}^{-1}\left(P_{t} f\right)=\frac{1}{t} \int_{0}^{t}\left(P_{s} f\right) \mathrm{d} s
$$

By Hopf, this operator is known to be $L^{p}$ bounded, $1<p<+\infty$. Stein must check in addition that the extension to complex values in the vertical line $z=-1+\mathrm{i} \tau$ also gives bounded operators on $L^{p}\left(\mathbb{R}^{n}\right)$.

Stein's objective is to study the maximal function of the semi-group itself, which corresponds to the derivative of order $z=0$. In order to do this, he interpolates between Hopf in $L^{p_{0}}, p_{0}<p<2$, for $\operatorname{Re} z=$ -1 , and an $L^{2}$ estimate of derivatives of the semi-group, for $\operatorname{Re} z=N$.

For each integer $k$, the quantity $t^{k} D_{t}^{k}\left(P_{t} f\right)$ appears in the LittlewoodPaley function $g_{k}(f)$, so one can control in $L^{2}$ its maximal function, see Section 2.1.1. The holomorphic family is then defined by $z \mapsto t^{z} D_{t}^{z}\left(P_{t} f\right)$, $z \in S$, for a suitable version $D^{z}$ of fractional differentiation.

The general strategy above was already applied in [71] to the discrete case.

### 7.2. Model of proof: the Poisson case

For proving Theorem 7.1, we have to apply Carbery's Proposition 6.14 (2) to the difference $K=K_{g}-P$. Müller shows that $(\xi \cdot \nabla)^{\alpha} m_{g}(\xi)$ acts on $L^{p}\left(\mathbb{R}^{n}\right)$ when $0<\alpha<1$ and $1<p<+\infty$, and we need to verify that the corresponding multiplier $(\xi \cdot \nabla)^{\alpha} \widehat{P}(\xi)$ for the Poisson kernel $P$ also acts on $L^{p}\left(\mathbb{R}^{n}\right), 1<p<+\infty$, with bounds independent of the dimension $n$. This could be covered by Proposition 7.2 , by observing that the Poisson kernel $P_{1}^{(n)}$ in (1.32), with Fourier transform $\mathrm{e}^{-2 \pi|\xi|}$, clearly satisfies (7.5. $\mathbf{H}_{\infty}$ ) and has $V\left(P_{1}^{(n)}\right)$ bounded independently of $n$ according to (7.4). We actually prefer to take an opportunity to examine the structure of Müller's proof in a simple case. When $\alpha \in(0,1)$, we could find a shorter specific proof, but the longer one that is given below provides a better introduction to what follows in this Section 7.

One sees that $(\xi \cdot \nabla)^{\alpha} \widehat{P}(\xi)=(2 \pi|\xi|)^{\alpha} \mathrm{e}^{-2 \pi|\xi|}$, either by applying (6.13) that gives $D_{t}^{\alpha} \mathrm{e}^{-\lambda|t|}=\lambda^{\alpha} \mathrm{e}^{-\lambda|t|}$ for $\lambda, t>0$, or by making use of the residue theorem.

Indeed, according to $\left(6.18 . \nabla^{\alpha}\right)$ with $\xi=|\xi| \theta$, one has

$$
(\xi \cdot \nabla)^{\alpha} \widehat{P}(\xi)=\int_{\mathbb{R}}(2 \mathrm{i} \pi s|\xi|)^{\alpha} \varphi_{\theta}(s) \mathrm{e}^{-2 \mathrm{i} \pi s|\xi|} \mathrm{d} s=\int_{\mathbb{R}}(2 \mathrm{i} \pi s|\xi|)^{\alpha} \frac{\mathrm{e}^{-2 \mathrm{i} \pi s|\xi|}}{\pi\left(1+s^{2}\right)} \mathrm{d} s
$$

that can be computed using a contour formed of $[-R, R]$ with $R>1$, and of a half-circle of radius $R$ centered at 0 , located in the lower complex half-plane.

We are going to bound the action on $L^{p}\left(\mathbb{R}^{n}\right)$ of the multiplier $|\xi|^{\alpha} \mathrm{e}^{-|\xi|}$ by the interpolation scheme of Section 3.2. Consider the holomorphic family of multipliers

$$
\mathcal{P}_{z}(\xi)=|\xi|^{z} \mathrm{e}^{-|\xi|}, \quad \operatorname{Re} z \geqslant 0, \quad \xi \in \mathbb{R}^{n}
$$

We will interpolate between $L^{2}\left(\mathbb{R}^{n}\right)$ and $L^{p_{0}}\left(\mathbb{R}^{n}\right), p_{0}>1$ close to 1 . For proving the boundedness on $L^{2}\left(\mathbb{R}^{n}\right)$, it is enough by (2.12.P) to see that the function $\xi \mapsto|\xi|^{z} \mathrm{e}^{-|\xi|}$ is bounded when $\xi$ varies in $\mathbb{R}^{n}$, and since this
function is radial, its supremum is independent of $n$. If we write $z=a+\mathrm{i} b$, $a \geqslant 0$, we have

$$
\begin{align*}
\sup _{\xi \in \mathbb{R}^{n}} \sup _{\operatorname{Re} z=a}\left|\mathcal{P}_{z}(\xi)\right| & =\sup _{\xi \in \mathbb{R}^{n}, b \in \mathbb{R}^{2}}\left\{\left.| | \xi\right|^{a+\mathrm{i} b} \mid \mathrm{e}^{-|\xi|}\right\}  \tag{7.13}\\
& =\sup _{r \geqslant 0}\left\{r^{a} \mathrm{e}^{-r}\right\}=a^{a} \mathrm{e}^{-a}
\end{align*}
$$

We work on a line $\operatorname{Re} z=\nu$, with $\nu$ "large", for dealing with the $L^{2}$ boundedness, and the other line is $\operatorname{Re} z=0$. For the values $z=0+\mathrm{i} b$, $b$ real, we know by (2.18) when $1<r<+\infty$ that the norm on $L^{r}\left(\mathbb{R}^{n}\right)$ of the multiplier $|\xi|^{i b}$ is bounded by $\lambda_{r} \mathrm{e}^{\pi|b| / 2}$, with $\lambda_{r}$ independent of the dimension $n$. The multiplier $\mathrm{e}^{-|\xi|}$ corresponds to the convolution with a Poisson probability measure, so it is bounded by 1 on $L^{r}\left(\mathbb{R}^{n}\right)$ when $1 \leqslant r \leqslant$ $+\infty$ by (2.13).

Let $\alpha \in(0,1)$ be given. Consider $p \in(1,2)$, introduce $p_{0}=2 p /(p+1) \in$ $(1, p)$, making $1 / p_{0}$ the midpoint between 1 and $1 / p$. Then with $\theta=p-1 \in$ $(0,1)$ we can check that $1 / p=(1-\theta) / p_{0}+\theta / 2$, and we define $\nu$ by the condition $\alpha=(1-\theta) .0+\theta \nu$, namely, we set $\nu=\alpha /(p-1)$. Let $T_{z}$ be the operator associated to the multiplier $\mathcal{P}_{z}$. We have to estimate the norm of $T_{\alpha}$ on $L^{p}$ by bounding $\left\langle T_{\alpha} f, g\right\rangle$ uniformly for $f$ in the unit ball of $L^{p}\left(\mathbb{R}^{n}\right)$ and $g$ in the unit ball of the dual $L^{q}\left(\mathbb{R}^{n}\right)$, where $1 / q+1 / p=1$. Consider the holomorphic function

$$
H: z \mapsto\left\langle T_{z} f_{z}, g_{z}\right\rangle
$$

where $f_{z}, g_{z}$ are as in (3.23). The bounds obtained for the family $T_{z}$ do not allow us to apply directly the three lines Lemma 3.1, but Corollary 3.4 will do the job. We got at the boundary of the strip, for the norms $\left\|T_{z}\right\|_{p_{0} \rightarrow p_{0}}$ when $\operatorname{Re} z=0$, a bound of the form $O\left(\mathrm{e}^{\kappa|\operatorname{Im} z|}\right)$. For every real number $\tau$, the function $H$ satisfies

$$
|H(0+\mathrm{i} \tau)| \leqslant \lambda_{p_{0}} \mathrm{e}^{\pi|\tau| / 2} \quad \text { and also }|H(\nu+\mathrm{i} \tau)| \leqslant \nu^{\nu} \mathrm{e}^{-\nu}
$$

By Corollary 3.4, the value $H(\alpha)$ is bounded uniformly by a quantity $\eta$ depending on $p_{0}, \theta$ and on the width $w=\nu$ of the strip, hence on $\alpha, p$ only. As explained in (3.26), this gives then for the action of $T_{\alpha}$ on $L^{p}\left(\mathbb{R}^{n}\right)$ a bound $\left\|T_{\alpha}\right\|_{p \rightarrow p} \leqslant \eta$.

For applying Corollary 3.4, it remains to check that $H$ has an admissible growth in $S=\{z: 0<\operatorname{Re} z<\nu\}$. We may actually reduce the discussion to a function $H$ bounded in the strip (but without universal estimate). Indeed, one can observe that all operators $T_{z}, z \in S$, are uniformly bounded on $L^{2}\left(\mathbb{R}^{n}\right)$, since $\left|\mathcal{P}_{z}(\xi)\right|$ is bounded by $\nu^{\nu}$ for all $\xi \in \mathbb{R}^{n}$ and $z$ in $S$ by (7.13). We may limit ourselves to $f, g$ continuous with compact support, so that $f_{z}, g_{z}, z \in S$, stay in a bounded subset of $L^{2}$, according to (3.24), implying that $H=H_{f, g}$ is bounded in the strip.

### 7.3. The interpolation part of Carbery's proof for Theorem 6.2

Proof. - In order to complete the proof of Theorem 6.2 and Proposition 6.3, it remains to show that the multiplier $(\xi \cdot \nabla)^{\alpha} m(\xi)$, where $m$ is the Fourier transform of $K=K_{g}-P$, is bounded on $L^{p}\left(\mathbb{R}^{n}\right)$ for at least one value $\alpha>1 / p$ when $p>3 / 2$. We have seen in the preceding Section 7.2 that $(\xi \cdot \nabla)^{\alpha} \widehat{P}(\xi)$ is bounded on $L^{p}\left(\mathbb{R}^{n}\right)$, we need only consider now $(\xi \cdot \nabla)^{\alpha} m_{g}(\xi)$. We will obtain the result by interpolating between the boundedness on $L^{1}\left(\mathbb{R}^{n}\right)$, for $\alpha_{0}=-\varepsilon$, and the boundedness on $L^{2}\left(\mathbb{R}^{n}\right)$, for $\alpha_{1}=1-\varepsilon$, of a certain holomorphic family $N_{z}(\xi)$ such that $N_{\alpha}(\xi)$ controls $(\xi \cdot \nabla)^{\alpha} m_{g}(\xi)$. If $p>3 / 2$ is fixed, its conjugate $q$ is $<3$. We write

$$
\frac{2}{3}>\frac{1}{p}=\frac{1-\theta}{1}+\frac{\theta}{2}=1-\frac{\theta}{2}, \quad \text { thus } \quad \frac{\theta}{2}=1-\frac{1}{p}=\frac{1}{q}
$$

and $\theta=2 / q>2 / 3>1-\theta / 2$. One can then find $\varepsilon \in(0,1)$ small enough, and independent of the dimension $n$, so that

$$
\alpha:=(1-\theta)(-\varepsilon)+\theta(1-\varepsilon)=\theta-\varepsilon>1-\frac{\theta}{2}=\frac{1}{p} .
$$

We need $0<\varepsilon<3 \theta / 2-1$, we can set for example $\varepsilon=3 \theta / 4-1 / 2=$ $(p-3 / 2) / p$. By Lemma 7.3, it is enough to show that $\left.d_{t}^{\alpha} m_{g}(t \xi)\right|_{t=1}$ is bounded on $L^{p}$. Consider the holomorphic family of multipliers $\left(N_{z}\right)$, simpler than that of Müller, namely, $N_{z}(\xi):=\left.d_{t}^{z} m_{g}(t \xi)\right|_{t=1}$ in the strip $-\varepsilon \leqslant \operatorname{Re} z \leqslant 1-\varepsilon$. When $\operatorname{Re} z<0$, we have

$$
\begin{equation*}
N_{z}(\xi)=\frac{1}{\Gamma(-z)} \int_{1}^{2}(u-1)^{-z-1} m_{g}(u \xi) \mathrm{d} u \tag{7.14}
\end{equation*}
$$

and in particular

$$
N_{-\varepsilon+\mathrm{i} \tau}(\xi)=\frac{1}{\Gamma(\varepsilon-\mathrm{i} \tau)} \int_{1}^{2}(u-1)^{\varepsilon-\mathrm{i} \tau-1} m_{g}(u \xi) \mathrm{d} u
$$

We see that

$$
\int_{1}^{2}\left|(u-1)^{\varepsilon-\mathrm{i} \tau-1}\right| \mathrm{d} u=\int_{1}^{2}(u-1)^{\varepsilon-1} \mathrm{~d} u=\varepsilon^{-1}<+\infty
$$

thus $N_{-\varepsilon+\mathrm{i} \tau}$ acts on $L^{1}$, with norm $\leqslant 2 \varepsilon^{-1}\left(1+\tau^{2}\right)^{1 / 4-\varepsilon / 2} \mathrm{e}^{\pi|\tau| / 2}$, according to Lemma 2.1, to the inequality (3.12. $\Gamma$ ) for the Gamma function and since the $L^{1}$ norm of the kernel $K_{g}$ is equal to 1 . When $\operatorname{Re} z=1-\varepsilon$, we have by (7.9) that

$$
N_{1-\varepsilon+\mathrm{i} \tau}(\xi)=\frac{m_{g}(2 \xi)}{\Gamma(\varepsilon-\mathrm{i} \tau)}-\frac{1}{\Gamma(\varepsilon-\mathrm{i} \tau)} \int_{1}^{2}(u-1)^{\varepsilon-\mathrm{i} \tau-1} \xi \cdot \nabla m_{g}(u \xi) \mathrm{d} u
$$

The kernel $N_{1-\varepsilon+\mathrm{i} \tau}$ is a bounded function of $\xi$, because we have $\left|m_{g}(2 \xi)\right| \leqslant$ $\delta_{0, g}$ and $\left|u \xi \cdot \nabla m_{g}(u \xi)\right| \leqslant \delta_{1, g}$ by (6.1.H). Using (3.12.Г) we obtain

$$
\begin{aligned}
\left|N_{1-\varepsilon+\mathrm{i} \tau}(\xi)\right| & \leqslant \frac{1}{|\Gamma(\varepsilon-\mathrm{i} \tau)|}\left(\delta_{0, g}+\int_{1}^{2}\left|(u-1)^{\varepsilon-\mathrm{i} \tau-1}\right| \frac{\delta_{1, g}}{u} \mathrm{~d} u\right) \\
& \leqslant 2\left(\delta_{0, g}+\delta_{1, g}\right) \varepsilon^{-1}\left(1+\tau^{2}\right)^{1 / 4-\varepsilon / 2} \mathrm{e}^{\pi|\tau| / 2}
\end{aligned}
$$

This shows that the operator associated to $N_{1-\varepsilon+\mathrm{i} \tau}(\xi)$ is bounded on $L^{2}\left(\mathbb{R}^{n}\right)$ with a bound $O\left(\mathrm{e}^{\kappa|\tau|}\right)$. We deal with this estimate as in the preceding Section 7.2 , and we obtain by interpolation that $N_{\alpha}(\xi)$ is a $L^{p}\left(\mathbb{R}^{n}\right)$-multiplier. Remark 3.6 takes care of the polynomial factor $\left(1+\tau^{2}\right)^{1 / 4-\varepsilon / 2} \leqslant\left(1+\tau^{2}\right)^{1 / 4}$. By Lemma 7.3 and (3.22) with $w=1, c_{j}=1 / 4, u_{j}=\pi / 2$, and since $\delta_{0, g} \geqslant 1$ we get

$$
\begin{aligned}
\left\|(\xi \cdot \nabla)^{\alpha} m_{g}(\xi)\right\|_{p \rightarrow p} & \leqslant 1+\frac{2 p}{p-3 / 2}\left(\frac{3}{2}\right)^{1 / 2} \mathrm{e}^{\pi / 4}\left(\delta_{0, g}+\delta_{1, g}\right)^{2-2 / p} \\
& \leqslant \kappa_{p} \Delta_{1}^{2-2 / p}
\end{aligned}
$$

We now check that the function $H(z)=\left\langle N_{z} f_{z}, g_{z}\right\rangle$ of (3.25) has an admissible growth in $S=\{-\varepsilon \leqslant \operatorname{Re} z \leqslant 1-\varepsilon\}$. We may again observe that all kernels $N_{z}(\xi)$ are bounded functions of $\xi$. Indeed, $N_{z}(\xi)$ can be expressed in the whole strip by

$$
N_{z}(\xi)=\frac{m_{g}(2 \xi)}{\Gamma(-z+1)}-\frac{1}{\Gamma(-z+1)} \int_{1}^{2}(u-1)^{-z} \xi \cdot \nabla m_{g}(u \xi) \mathrm{d} u, \quad \xi \in \mathbb{R}^{n}
$$

so that $\left|N_{z}(\xi)\right| \leqslant \kappa_{\varepsilon, \delta}\left(1+\tau^{2}\right)^{1 / 4} \mathrm{e}^{\pi|\tau| / 2}$. Next, we can assume that the two functions $f, g$ appearing in the definition of $H$ are bounded with bounded support, and argue with (3.24) as at the end of Section 7.2, obtaining that $|H(z)| \leqslant \kappa\left\|N_{z}\right\|_{2 \rightarrow 2} \leqslant \kappa_{\varepsilon}^{\prime}\left(1+\tau^{2}\right)^{1 / 4} \mathrm{e}^{\pi|\tau| / 2}$, a growth admissible for applying Corollary 3.4.

We see pretty well why Müller finds a better result than the one given by the preceding argument, which suffices for Carbery's theorem. It is because Müller is able to make use of multipliers more difficult to handle, which contain an extra factor $|\xi|$ on the line $\operatorname{Re} z=-\varepsilon$, for example $m_{-\varepsilon}^{\varepsilon}(\xi)=$ $(1+|\xi|) N_{-\varepsilon}(\xi)$ when $z=-\varepsilon$. This factor $|\xi|$ is precisely the one that will be treated by the geometrical parameter $q(C)$. On the other hand, Müller's approach is not better when $p>3 / 2$, since the result is known in this case without assumption on $q(C)$.

Remark. - The factor $1 / \Gamma(-z)$ in (7.14) is not purely decorative. Without it, $N_{z}(\xi)$ would have a "pole" at $z=0$, which is compensated by the zero of $1 / \Gamma(z)$ at 0 . One could perhaps get away here with a less sophisticated factor such as $z /(a-z)$, with $a$ real and $>1$. See also Remark 7.13.

### 7.4. Upper bounds for the functions $\xi \mapsto m_{z}^{\varepsilon}(\xi)$

We present a version of Müller's upper bounds for the functions $m_{z}^{\varepsilon}$ defined in (7.12). Müller's bounds in [59] are not fully explicit since they use asymptotic estimates, but they do not contain the annoying factor $\varepsilon^{-1}$ that our somewhat shorter proof introduces below.

Lemma 7.4 ([59, Lemma 2]). - Assume that the kernel $K_{g}$ integrable on $\mathbb{R}^{n}$ satisfies $\left(7.5 . \mathbf{H}_{\infty}\right)$, let $\varepsilon \in(0,1)$, let $\nu \geqslant 1-\varepsilon$ and set $\ell=\lceil\nu+\varepsilon\rceil$. For every $z \in \mathbb{C}$ such that $-\varepsilon \leqslant \operatorname{Re} z \leqslant \nu$, one has that

$$
\forall \xi \in \mathbb{R}^{n}, \quad\left|m_{z}^{\varepsilon}(\xi)\right| \leqslant \kappa_{\nu} \varepsilon^{-1} \Delta_{\ell}\left(1+(\operatorname{Im} z)^{2}\right)^{\nu / 2-1 / 4} \mathrm{e}^{\pi|\operatorname{Im} z| / 2}
$$

where $\kappa_{\nu}=4 \Gamma(\max (\nu, 2))(3 / 2)^{\nu-1 / 2} \mathrm{e}^{\pi / 4}$ and where $\Delta_{\ell}$ is defined at (7.6).
One of the difficulties in Müller's article is the following: with the operator $D^{\alpha}$, we have been able to compute certain integrals by the residue theorem, on entire half-lines. The corresponding values for $d^{\alpha}$ are less pleasant, because they involve bounded segments, and quarters of circle at finite distance whose contribution is not zero. Let us mention another difficulty, somewhat related to the latter. If we know that $\left|D_{t}^{z} m(t \theta)\right| \leqslant \kappa(1+|t|)^{-1}$ for every $t$ real and $\theta \in S^{n-1}$, then by the homogeneity relation (6.8) we get $\left|D_{t}^{z} m(t \xi)\right| \leqslant$ $\kappa|\xi|^{\operatorname{Re} z}(1+|t \xi|)^{-1}$ for $\xi \in \mathbb{R}^{n}$, but this kind of behavior is not clear for $d^{z}$. The more delicate analysis of [59] will not be given here, but some special cases are rather easy. Indeed, the computation is not difficult when $\operatorname{Re} z=k-\varepsilon$, for every integer $k \geqslant 0$. We will however be able to deduce Lemma 7.4 from the easy cases that are treated in the next lemma.

Lemma 7.5. - Assume that $K_{g}$ is integrable on $\mathbb{R}^{n}$ and satisfies (7.5. $\left.\mathbf{H}_{\infty}\right)$. For every $\varepsilon \in(0,1)$, every integer $k \geqslant 0$ and $z \in \mathbb{C}$ such that $\operatorname{Re} z=k-\varepsilon$, one has

$$
\forall \xi \in \mathbb{R}^{n}, \quad\left|m_{z}^{\varepsilon}(\xi)\right| \leqslant 2 \kappa_{k}^{\prime} \varepsilon^{-1} \Delta_{k}\left(1+(\operatorname{Im} z)^{2}\right)^{k^{*} / 2-\varepsilon / 2-1 / 4} \mathrm{e}^{\pi|\operatorname{Im} z| / 2}
$$

where $k^{*}=\max (k, 1), \kappa_{k}^{\prime} \leqslant \Gamma(k-\varepsilon)$ for $k \geqslant 3$ and $\kappa_{0}^{\prime}, \kappa_{1}^{\prime}, \kappa_{2}^{\prime} \leqslant 1$.
Proof. - We first give the proof for $k=0$, when $z=-\varepsilon+\mathrm{i} \tau$. We have

$$
m_{-\varepsilon+\mathrm{i} \tau}^{\varepsilon}(\xi)=(1+|\xi|)^{1-\varepsilon-(-\varepsilon+\mathrm{i} \tau)} \frac{1}{\Gamma(\varepsilon-\mathrm{i} \tau)} \int_{1}^{2}(u-1)^{\varepsilon-\mathrm{i} \tau-1} m_{g}(u \xi) \mathrm{d} u
$$

and it follows that

$$
\left|m_{-\varepsilon+\mathrm{i} \tau}^{\varepsilon}(\xi)\right| \leqslant\left|\frac{1}{\Gamma(\varepsilon-\mathrm{i} \tau)}\right| \int_{1}^{2}(u-1)^{\varepsilon-1}(1+|\xi|)\left|m_{g}(u \xi)\right| \mathrm{d} u
$$

By $\left(7.5 . \mathbf{H}_{\infty}\right)$, we know that $\left|m_{g}(u \xi)\right| \leqslant \delta_{0, g}(1+|\xi|)^{-1}$ when $u \geqslant 1$, thus

$$
\left|\Gamma(\varepsilon-\mathrm{i} \tau) m_{-\varepsilon+\mathrm{i} \tau}^{\varepsilon}(\xi)\right| \leqslant \delta_{0, g} \int_{1}^{2}(u-1)^{\varepsilon-1} \mathrm{~d} u=\frac{\delta_{0, g}}{\varepsilon}
$$

Using also (3.12.Г), this simplest case reads as

$$
\left\|m_{-\varepsilon+\mathrm{i} \tau}^{\varepsilon}\right\|_{\infty} \leqslant 2 \delta_{0, g} \varepsilon^{-1}\left(1+\tau^{2}\right)^{1 / 4-\varepsilon / 2} \mathrm{e}^{\pi|\tau| / 2}, \quad \tau \in \mathbb{R}
$$

and $1 / 4=k^{*} / 2-1 / 4$ here. For $k>0$, we have by (7.10) with $z=k-\varepsilon+\mathrm{i} \tau$ that

$$
\left.d_{t}^{k-\varepsilon+\mathrm{i} \tau} m_{g}(t \xi)\right|_{t=1}=E_{k}+(-1)^{k} \frac{1}{\Gamma(\varepsilon-\mathrm{i} \tau)} \int_{1}^{2}(u-1)^{\varepsilon-\mathrm{i} \tau-1} \frac{\mathrm{~d}^{k}}{\mathrm{~d} u^{k}} m_{g}(u \xi) \mathrm{d} u
$$

where

$$
E_{k}=\sum_{j=0}^{k-1}(-1)^{j} \frac{\left.\frac{\mathrm{~d}^{j}}{\mathrm{~d} u^{j}} m_{g}(u \xi)\right|_{u=2}}{\Gamma(j+1-k+\varepsilon-\mathrm{i} \tau)}
$$

By our assumption $\left(7.5 . \mathbf{H}_{\infty}\right)$, the function $u \mapsto m_{g}(u \xi)$ satisfies

$$
\begin{equation*}
\forall u \geqslant 1, \quad\left|\frac{\mathrm{~d}^{j}}{\mathrm{~d} u^{j}} m_{g}(u \xi)\right|=\left|\frac{\mathrm{d}^{j}}{\mathrm{~d} u^{j}} m_{g}(u|\xi| \theta)\right| \leqslant \delta_{j, g} \frac{|\xi|^{j}}{1+|\xi|} \tag{7.15}
\end{equation*}
$$

for each integer $j \geqslant 0$, if $\xi \neq 0$ and $\theta=|\xi|^{-1} \xi$. This yields

$$
\begin{aligned}
\left|\int_{1}^{2}(u-1)^{\varepsilon-\mathrm{i} \tau-1} \frac{\mathrm{~d}^{k}}{\mathrm{~d} u^{k}} m_{g}(u \xi) \mathrm{d} u\right| & \leqslant \delta_{k, g} \int_{1}^{2}(u-1)^{\varepsilon-1} \frac{|\xi|^{k}}{1+|\xi|} \mathrm{d} u \\
& =\frac{\delta_{k, g}}{\varepsilon} \frac{|\xi|^{k}}{1+|\xi|}
\end{aligned}
$$

For the terms in the expression $E_{k}$, we have by (7.15) that

$$
\begin{aligned}
\left|\frac{\mathrm{d}^{j}}{\mathrm{~d} u^{j}} m_{g}(u \xi)\right|_{u=2} \left\lvert\, \leqslant \delta_{j, g} \frac{|\xi|^{j}}{1+|\xi|} \leqslant \delta_{j, g} \frac{(1+|\xi|)^{j}}{1+|\xi|} \leqslant \delta_{j, g}(1+|\xi|)^{k-1}\right. \\
j=0, \ldots, k-1
\end{aligned}
$$

Recalling $\Delta_{k}=\sum_{j=0}^{k} \delta_{j, g}$ and (3.12. $\Gamma$ ) with $a=-k+1+\varepsilon$, we get

$$
\begin{align*}
\left|m_{k-\varepsilon+\mathrm{i} \tau}^{\varepsilon}(\xi)\right|= & (1+|\xi|)^{1-\varepsilon-(k-\varepsilon)}\left|d_{t}^{k-\varepsilon+\mathrm{i} \tau} m_{g}(t \xi)\right|_{t=1} \mid \\
\leqslant & \Delta_{k} \varepsilon^{-1}(1+|\xi|)^{1-k}(1+|\xi|)^{k-1} \\
& \quad \times \max \left\{\left|\Gamma\left(\varepsilon-\mathrm{i} \tau-j_{1}\right)\right|^{-1}: 0 \leqslant j_{1} \leqslant k-1\right\}  \tag{7.16}\\
\leqslant & \beta_{a} \Delta_{k} \varepsilon^{-1}\left(1+\tau^{2}\right)^{1 / 4+(k-1-\varepsilon) / 2} \mathrm{e}^{\pi|\tau| / 2}
\end{align*}
$$

We may take $\beta_{a}=2$ when $k \leqslant 2$ and $\beta_{a}=2 \Gamma(k-\varepsilon)$ otherwise.
Remark. - One could not make the same simple computation for $k-\varepsilon^{\prime}$ when $\varepsilon^{\prime}>\varepsilon$. Indeed, we have then

$$
m_{k-\varepsilon^{\prime}}^{\varepsilon}(\xi)=\left.(1+|\xi|)^{1-\left(k-\varepsilon^{\prime}\right)-\varepsilon} d_{t}^{k-\varepsilon^{\prime}+\mathrm{i} \tau} m_{g}(t \xi)\right|_{t=1}
$$

so $\left|m_{k-\varepsilon^{\prime}}^{\varepsilon}(\xi)\right|$ contains the factor $(1+|\xi|)^{1-k+\varepsilon^{\prime}-\varepsilon}$, that is not controllable by the preceding proof when $\varepsilon^{\prime}>\varepsilon$. With one more integration by parts in the log-concave case we obtain

$$
\frac{\mathrm{d}^{j}}{\mathrm{~d} u^{j}} m_{l c}(u \xi)=\frac{(-2 \mathrm{i} \pi|\xi|)^{j-2}}{u^{2}} \int_{\mathbb{R}}\left(s^{j} \varphi_{\theta}(s)\right)^{\prime \prime} \mathrm{e}^{-2 \mathrm{i} \pi s u|\xi|} \mathrm{d} s
$$

that seems to give an additional improvement, able to swallow the bad factor $|\xi|^{\varepsilon^{\prime}-\varepsilon}$ above. However, we would need now for $\int_{\mathbb{R}}|s|^{j}\left|\varphi_{\theta}^{\prime \prime}(s)\right| \mathrm{d} s$ a universal bound that does not exist, and actually, this integral does not make sense in general.

When $k \geqslant 1$, the kernel $\widetilde{m}_{k-\varepsilon+\mathrm{i} \tau}^{\varepsilon}:=\left.(1+|\xi|)^{1-k-\mathrm{i} \tau} D_{u}^{k-\varepsilon+\mathrm{i} \tau} m_{g}(u \xi)\right|_{u=1}$ is even easier to bound since we can write directly

$$
\left|\int_{1}^{+\infty}(u-1)^{\varepsilon-\mathrm{i} \tau-1} \frac{\mathrm{~d}^{k}}{\mathrm{~d} u^{k}} m_{g}(u \xi) \mathrm{d} u\right| \leqslant \delta_{k, g}\left(\int_{1}^{+\infty}(u-1)^{\varepsilon-1} \frac{\mathrm{~d} u}{u}\right)|\xi|^{k-1}
$$

but $\left.D_{u}^{-\varepsilon+\mathrm{i} \tau} m_{g}(u \xi)\right|_{u=1}$ is not a bounded function of $\xi$ in the neighborhood of $\xi=0$. For example, we have $\left.D_{u}^{-\varepsilon} \mathrm{e}^{-u|\xi|}\right|_{u=1}=|\xi|^{-\varepsilon} \mathrm{e}^{-|\xi|}$. Thus $\widetilde{m}_{-\varepsilon}^{\varepsilon}$ is not an $L^{2}$ multiplier, nor an $L^{p}$ multiplier for any $p \neq 2$, and this justifies working with $d_{t}^{z}$ instead.

Proof of Lemma 7.4. - Let $\nu>1-\varepsilon$ and $\ell=\lceil\nu+\varepsilon\rceil \geqslant 2$, so that $\nu \leqslant \ell-\varepsilon$. If $\operatorname{Re} z<\ell$, we have by (7.10) that

$$
\left.d_{t}^{z} m_{g}(t \xi)\right|_{t=1}=E_{\ell}(z)+(-1)^{\ell} \frac{1}{\Gamma(\ell-z)} \int_{1}^{2}(u-1)^{-z+\ell-1} \frac{\mathrm{~d}^{\ell}}{\mathrm{d} u^{\ell}} m_{g}(u \xi) \mathrm{d} u
$$

with $E_{\ell}(z)=\left.\sum_{i=0}^{\ell-1}(-1)^{i} \Gamma(i+1-z)^{-1} \frac{\mathrm{~d}^{i}}{\mathrm{~d} u^{i}} m_{g}(u \xi)\right|_{u=2}$. We fix $\xi \in \mathbb{R}^{n}$ and consider the holomorphic function

$$
H_{\xi}: z \mapsto m_{z}^{\varepsilon}(\xi)=\left.(1+|\xi|)^{1-\varepsilon-z} d_{t}^{z} m_{g}(t \xi)\right|_{t=1}
$$

in the strip $-\varepsilon<\operatorname{Re} z<\ell-\varepsilon$. We have $\left|\left(\mathrm{d}^{i} / \mathrm{d} u^{i}\right) m_{g}(u \xi)\right| \leqslant \delta_{i, g}|\xi|^{i}$ by $\left(7.5 . \mathbf{H}_{\infty}\right)$, and it follows from (3.12. $)$ that $\left|H_{\xi}(z)\right| \leqslant \kappa \mathrm{e}^{\kappa|\tau|}$ in the strip, with $\kappa$ depending on $|\xi|$. Consider an arbitrary $z_{0}$ such that $1-\varepsilon<\nu_{0}:=$ $\operatorname{Re} z_{0} \leqslant \nu$. Let $k$ integer be such that $k-\varepsilon<\nu_{0} \leqslant k+1-\varepsilon$, thus $1 \leqslant k<\ell$. By Lemma 7.5 , when $\operatorname{Re} z=k-\varepsilon$ or $\operatorname{Re} z=k+1-\varepsilon$, we have for $H_{\xi}(z)$ the good bound

$$
\begin{equation*}
\left|H_{\xi}(z)\right| \leqslant 2 \kappa_{\operatorname{Re} z+\varepsilon}^{\prime} \Delta_{\operatorname{Re} z+\varepsilon} \varepsilon^{-1}\left(1+(\operatorname{Im} z)^{2}\right)^{\operatorname{Re} z / 2-1 / 4} \mathrm{e}^{\pi|\operatorname{Im} z| / 2} \tag{7.17}
\end{equation*}
$$

We write $\Delta_{k}<\Delta_{k+1} \leqslant \Delta_{\ell}$ and $\nu_{0}+\varepsilon=(1-\theta) k+\theta(k+1)$. When $k \geqslant 3$, we have

$$
\begin{aligned}
&{\kappa_{k}^{\prime 1-\theta}{\kappa^{\prime}}_{k+1}^{\theta}}^{\leqslant} \leqslant \Gamma(k-\varepsilon)^{1-\theta} \Gamma(k+1-\varepsilon)^{\theta}=\frac{\Gamma(k+1-\varepsilon)}{(k-\varepsilon)^{1-\theta}} \\
& \leqslant \frac{\Gamma\left(\nu_{0}\right)^{\theta} \Gamma\left(\nu_{0}+1\right)^{1-\theta}}{(k-\varepsilon)^{1-\theta}}=\frac{\nu_{0}^{1-\theta}}{(k-\varepsilon)^{1}-\theta} \Gamma\left(\nu_{0}\right)<2 \Gamma(\nu)
\end{aligned}
$$

and $\kappa_{1}^{\prime 1-\theta} \kappa^{\prime}{ }_{2}^{\theta} \leqslant 1, \kappa^{\prime 1-\theta} \kappa^{\prime}{ }_{3}^{\theta} \leqslant 2$. By Corollary 3.4 and Remark 3.6, (3.22) with $w=1$ and $c_{j}=(k+j-\varepsilon) / 2-1 / 4, j=0,1$, we get for $\left|H_{\xi}\left(z_{0}\right)\right|$ a bound
$4 \Gamma(\max (\nu, 2))(3 / 2)^{\operatorname{Re} z_{0}-1 / 2} \mathrm{e}^{\pi / 4} \varepsilon^{-1} \Delta_{\ell}\left(1+\left(\operatorname{Im} z_{0}\right)^{2}\right)^{\operatorname{Re} z_{0} / 2-1 / 4} \mathrm{e}^{\pi\left|\operatorname{Im} z_{0}\right| / 2}$.
This proves Lemma 7.4 when $1-\varepsilon<\operatorname{Re} z_{0} \leqslant \nu$. The case $-\varepsilon \leqslant \operatorname{Re} z_{0} \leqslant \nu=$ $1-\varepsilon$ is left to the reader, one has $k=0$ and the polynomial component of the bound is then $\left(1+\tau^{2}\right)^{\nu / 2-1 / 4}$ on both sides of the strip $-\varepsilon<\operatorname{Re} z<1-\varepsilon$.

An alternative proof could go like this: divide the integral $\int_{1}^{2}$ in the definition of $\left.d_{t}^{z} m_{g}(t \xi)\right|_{t=1}$ into $\int_{1}^{1+a}$ and $\int_{1+a}^{2}$, for some suitable $a \in[0,1]$. For the first integral $\int_{1}^{1+a}$, we modify (7.10) and get when $-1<\operatorname{Re} z<0$ that

$$
\begin{aligned}
d_{z, 1}(a):= & \frac{1}{\Gamma(-z)} \int_{1}^{1+a}(u-1)^{-z-1} m_{g}(u \xi) \mathrm{d} u \\
= & E_{k+2, a}(z)+(-1)^{k+2} \frac{1}{\Gamma(k+2-z)} \\
& \times \int_{1}^{1+a}(u-1)^{-z+k+1} \frac{\mathrm{~d}^{k+2}}{\mathrm{~d} u^{k+2}} m_{g}(u \xi) \mathrm{d} u
\end{aligned}
$$

for every integer $k \geqslant-1$, where $E_{k+2, a}(z)$ is equal to

$$
E_{k+2, a}(z)=\sum_{j=0}^{k+1}(-1)^{j} \frac{\left.a^{-z+j} \frac{\mathrm{~d}^{j}}{\mathrm{~d} u^{j}} m_{g}(u \xi)\right|_{u=1+a}}{\Gamma(j+1-z)}
$$

Let now $-1<\operatorname{Re} z \leqslant \nu$ and write $z=k+\sigma+\mathrm{i} \tau$ with $k$ integer and $0<\sigma \leqslant 1$. Applying the preceding formulas it follows by (7.5.H $\mathbf{H}_{\infty}$ ) that

$$
\left|d_{z, 1}(a)\right| \leqslant \sum_{j=0}^{k+1} \frac{a^{-k-\sigma+j} \delta_{j, g}|\xi|^{j}}{|\Gamma(j+1-z)|(1+|\xi|)}+\frac{(2-\sigma)^{-1} a^{2-\sigma} \delta_{k+2, g}|\xi|^{k+2}}{|\Gamma(k+2-z)|(1+|\xi|)} .
$$

When $|\xi| \leqslant 1$, we choose $a=1$ and obtain $\left|d_{z, 1}(a)\right| \leqslant C_{k}(z)(1+|\xi|)^{-1}$ where

$$
C_{k}(z)=\Delta_{k+2} \max \left\{|\Gamma(i-z)|^{-1}: 0 \leqslant i \leqslant k+2\right\} .
$$

When $|\xi|>1$, we let $a=|\xi|^{-1}$ and get $\left|d_{z, 1}(a)\right| \leqslant C_{k}(z)|\xi|^{k+\sigma}(1+|\xi|)^{-1}$. The other term $d_{z, 2}(a)$, corresponding to $\int_{1+a}^{2}$, is zero when $|\xi| \leqslant 1$ since
$a=1$ in this case. Otherwise, we have $a=|\xi|^{-1}$ and assuming $k+\sigma \neq 0$, we get

$$
\begin{aligned}
\left|d_{z, 2}(a)\right| & =\left|\frac{1}{\Gamma(-z)} \int_{1+|\xi|^{-1}}^{2}(u-1)^{-z-1} m_{g}(u \xi) \mathrm{d} u\right| \\
& \leqslant \frac{1}{|\Gamma(-z)|} \frac{|\xi|^{k+\sigma}-1}{k+\sigma} \frac{\delta_{0, g}}{1+|\xi|} .
\end{aligned}
$$

There is no problem as long as $\operatorname{Re} z=k+\sigma$ is not close to 0 . Otherwise, we can apply $\left||\xi|^{k+\sigma}-1\right| \leqslant|k+\sigma||\xi|^{(k+\sigma)^{+}} \ln |\xi|$, where $t^{+}=\max (t, 0)$. Summing up and letting $L(\xi)=1+(\ln |\xi|)^{+}$, we have when $-1<\operatorname{Re} z=$ : $s<\nu$ that

$$
\left|d_{t}^{z} m_{g}(t \xi)\right|_{t=1} \mid \leqslant C_{\nu}(z)\left(\left[1+|s|^{-1}\right] \wedge L(\xi)\right)(1+|\xi|)^{s^{+}-1}
$$

giving bounds multiple of $(1+|\xi|)^{-1},(1+|\xi|)^{-1} L(\xi)$ for $s$ in $[-1,-\varepsilon / 2]$ and $[-\varepsilon / 2,0]$ respectively, $(1+|\xi|)^{s-1} L(\xi)$ and $(1+|\xi|)^{s-1}$ in $[0, \varepsilon / 2]$ and $[\varepsilon / 2, \nu]$ respectively. For $m_{z}^{\varepsilon}(\xi)$, we get bounds multiple of $1,(1+|\xi|)^{-\varepsilon / 2} L(\xi)$ and $(1+|\xi|)^{-\varepsilon}$ for $s$ in $[-\varepsilon,-\varepsilon / 2],[-\varepsilon / 2, \varepsilon / 2]$ and $[\varepsilon / 2, \nu]$ respectively. This shows that $m_{z}^{\varepsilon}(\xi)$ is a bounded function of $\xi \in \mathbb{R}^{n}$.

### 7.5. Lemma 4 of Müller's article

We must control the action on $L^{p}\left(\mathbb{R}^{n}\right), p>1$ close to 1 , of multipliers $m_{z}^{\varepsilon}$ when $\operatorname{Re} z=-\varepsilon$. If $z=-\varepsilon+\mathrm{i} \tau$, we have

$$
\Gamma(\varepsilon-\mathrm{i} \tau) m_{-\varepsilon+\mathrm{i} \tau}^{\varepsilon}(\xi)=(1+|\xi|)^{1-\mathrm{i} \tau} \int_{1}^{2}(s-1)^{\varepsilon-\mathrm{i} \tau-1} m_{g}(s \xi) \mathrm{d} s
$$

Since $\int_{1}^{2}\left|(s-1)^{\varepsilon-\mathrm{i} \tau-1}\right| \mathrm{d} s=\varepsilon^{-1}$, it is enough to bound uniformly in $s \in[1,2]$ the norm of $n_{s}(\xi):=(1+|\xi|)^{1-\mathrm{i} \tau} m_{g}(s \xi)$. This multiplier can be decomposed into several parts: first $(1+|\xi|)^{-\mathrm{i} \tau}$, which is taken care of by Proposition 2.2 on multipliers of Laplace type. Indeed, replacing $\lambda$ by $1+\lambda$ in (2.17) and integrating by parts, one finds that $(1+\lambda)^{-\mathrm{i} \tau}=\lambda \int_{0}^{+\infty} \mathrm{e}^{-\lambda t} a_{\tau}(t) \mathrm{d} t$, with

$$
\begin{equation*}
a_{\tau}(t)=\frac{1}{\Gamma(1+\mathrm{i} \tau)}\left(t^{\mathrm{i} \tau} \mathrm{e}^{-t}+\int_{0}^{t} s^{\mathrm{i} \tau} \mathrm{e}^{-s} \mathrm{~d} s\right) \tag{7.18}
\end{equation*}
$$

and $\left|a_{\tau}(t)\right| \leqslant|\Gamma(1+\mathrm{i} \tau)|^{-1} \leqslant\left(1+\tau^{2}\right)^{-1 / 4} \mathrm{e}^{\pi|\tau| / 2}$ according to (3.4). Next, in $n_{s}(\xi)$, we have $(1+|\xi|) m_{g}(s \xi)$, which is formed of $m_{g}(s \xi)$, multiplier bounded by $\left\|K_{g}\right\|_{L^{1}\left(\mathbb{R}^{n}\right)}$ on all $L^{p}\left(\mathbb{R}^{n}\right)$ spaces, and of $s^{-1}|s \xi| m_{g}(s \xi), s>1$, with a multiplier norm less than that of $|\xi| m_{g}(\xi)$, according to (2.10).

Given an integrable kernel $K$ on $\mathbb{R}^{n}$ and its Fourier transform $m$, the question boils down to handling the crucial multiplier

$$
\begin{equation*}
m^{\#}(\xi):=|\xi| m(\xi) \tag{7.19}
\end{equation*}
$$

We summarize the latter discussion in the lemma that follows, where we include the bound $2\left(1+\tau^{2}\right)^{1 / 4} \mathrm{e}^{\pi|\tau| / 2}$ from $(3.12 . \Gamma)$ for the factor $|\Gamma(\varepsilon-\mathrm{i} \tau)|^{-1}$ that was left apart above. So far, the kernel $K$ can be arbitrary in $L^{1}\left(\mathbb{R}^{n}\right)$.

Lemma 7.6. - Let $p$ belong to $(1,2]$. One has that

$$
\left\|m_{-\varepsilon+\mathrm{i} \tau}^{\varepsilon}\right\|_{p \rightarrow p} \leqslant 2 \varepsilon^{-1} \lambda_{p} \mathrm{e}^{\pi|\tau|}\left(\left\|K_{g}\right\|_{L^{1}\left(\mathbb{R}^{n}\right)}+\left\|m^{\#}\right\|_{p \rightarrow p}\right), \quad \tau \in \mathbb{R}
$$

where $\lambda_{p}$ is the constant appearing in Proposition 2.2.

The serious work will be done in the proof of the following essential lemma.

Lemma 7.7. - Let $K_{g}$ be a kernel integrable on $\mathbb{R}^{n}$ satisfying (6.1.H), and $m_{g}$ its Fourier transform. Let $m_{g}^{\#}$ be defined by (7.19) and $p \in(1,2]$. One has that

$$
\left\|m_{g}^{\#}\right\|_{p \rightarrow p} \leqslant(2 \pi)^{1-2 / p} \rho_{p} \delta_{0, g}^{2-2 / p} V\left(K_{g}\right)^{-1+2 / p}
$$

where $\rho_{p}$ is the constant from (2.22) and where $V\left(K_{g}\right)$ is defined at (7.2).

The proof of Lemma 7.7 will be broken into several easy statements. Some of them are used again in Section 8. To begin with, we merely assume that $K$ is an integrable kernel on $\mathbb{R}^{n}$ having partial derivatives $\partial_{j} K$ in the sense of distributions that are (signed) measures $\mu_{j}$, and we let $m=\widehat{K}$. We can express $m^{\#}(\xi)$ with the help of the Riesz transforms $\left(R_{j}\right)_{j=1}^{n}$ introduced in Section 2.3, writing

$$
2 \pi m^{\#}(\xi)=\sum_{j=1}^{n} \frac{-\mathrm{i} \xi_{j}}{|\xi|}\left(2 \mathrm{i} \pi \xi_{j}\right) m(\xi)
$$

The functions $\left(2 \mathrm{i} \pi \xi_{j}\right) m(\xi), j=1, \ldots, n$, are the Fourier transforms of the measures $\mu_{j}=\partial_{j} K$. When $K$ is the uniform probability density $K_{C}$ on a symmetric convex set $C$, the $\mu_{j}$ s are supported on the boundary of $C$, and we shall see below that $V\left(K_{C}\right)=q(C)$ if $C$ is isotropic and normalized by variance.

The convolution operator $T_{m \#}$ can thus be written under the form

$$
T_{m \#}: f \mapsto T_{m \#} f=(2 \pi)^{-1} \sum_{j=1}^{n} R_{j} \mu_{j} * f
$$

Riesz transforms commute with convolutions. If $g$ is in the dual $L^{q}$ of $L^{p}$, we have

$$
\begin{aligned}
2 \pi\left|\left\langle T_{m \#}^{\#} f, g\right\rangle\right| & =\left|\sum_{j=1}^{n}\left\langle R_{j} \mu_{j} * f, g\right\rangle\right|=\left|\sum_{j=1}^{n}\left\langle\left(R_{j} f\right) * \mu_{j}, g\right\rangle\right|=\left|\sum_{j=1}^{n}\left\langle R_{j} f, \widetilde{\mu}_{j} * g\right\rangle\right| \\
& \leqslant \int_{\mathbb{R}^{n}}\left(\sum_{j=1}^{n}\left|R_{j} f\right|^{2}\right)^{1 / 2}\left(\sum_{j=1}^{n}\left|\widetilde{\mu}_{j} * g\right|^{2}\right)^{1 / 2}
\end{aligned}
$$

where $\widetilde{\mu}_{j}$ denotes the image of $\mu_{j}$ under the symmetry $x \mapsto-x$ of $\mathbb{R}^{n}$. By (2.22), the Riesz transforms are "collectively bounded" in $L^{p}\left(\mathbb{R}^{n}\right)$ by a constant $\rho_{p}$ independent of the dimension $n$, and we obtain therefore that

$$
2 \pi\left|\left\langle T_{m \#} f, g\right\rangle\right| \leqslant \rho_{p}\|f\|_{p}\left\|\left(\sum_{j=1}^{n}\left|\widetilde{\mu}_{j} * g\right|^{2}\right)^{1 / 2}\right\|_{q}
$$

Noticing that $\widetilde{\mu}_{j} * g=\left(\mu_{j} * \widetilde{g}\right)^{\sim}$ and $\left(\sum_{j=1}^{n}\left|\mu_{j} * g\right|^{2}\right)^{1 / 2}=|\nabla K * g|$, we are led to study the operator

$$
\begin{equation*}
U_{K}: g \in L^{q}\left(\mathbb{R}^{n}\right) \mapsto \nabla K * g \in L^{q}\left(\mathbb{R}^{n}, \mathbb{R}^{n}\right) \tag{7.20}
\end{equation*}
$$

given by the vector-valued convolution with $\nabla K$. Let us state what we have got.

Lemma 7.8. - Let $K$ be an integrable kernel on $\mathbb{R}^{n}$, $m$ its Fourier transform and let $m^{\#}$ be defined by (7.19). For every $p \in(1,2]$ and $q=p /(p-1)$, one has

$$
\left\|T_{m \#}\right\|_{p \rightarrow p} \leqslant(2 \pi)^{-1} \rho_{p} \sup _{\|g\|_{q} \leqslant 1}\|\nabla K * g\|_{L^{q}\left(\mathbb{R}^{n}\right)}=(2 \pi)^{-1} \rho_{p}\left\|U_{K}\right\|_{q \rightarrow q}
$$

When $K=K_{g}$ satisfies (6.1.H), we shall estimate $\left\|U_{K_{g}}\right\|_{q \rightarrow q}$ by interpolation between $L^{2}$ and $L^{\infty}$. Contrary to the $L^{2}$ estimate which will make use of (6.1.H), the $L^{\infty}$ estimate is a straightforward observation following from the definition of $V(K)$. In the special case $K_{g}=K_{C}$ of a convex body $C$, this $L^{\infty}$ case will bring in the geometrical parameter $q(C)=2 Q\left(C_{0}\right) L\left(C_{0}\right)$, equal to $V\left(K_{C}\right)$.

Lemma 7.9. - Let $K$ be an integrable kernel on $\mathbb{R}^{n}$ having a finite directional variation $V(K)$, and let $U_{K}$ be defined by (7.20). One has that

$$
\begin{equation*}
\left\|U_{K}\right\|_{\infty \rightarrow \infty} \leqslant V(K) \tag{7.21}
\end{equation*}
$$

Proof. - For each $x \in \mathbb{R}^{n}$, the Euclidean norm of the vector $(\nabla K *$ $g)(x) \in \mathbb{R}^{n}$ is given by the supremum over $\theta \in S^{n-1}$ of

$$
\begin{aligned}
\left|\theta \cdot\left(\int_{\mathbb{R}^{n}} g(x-y) \mathrm{d}(\nabla K)(y)\right)\right| & =\left|\int_{\mathbb{R}^{n}} g(x-y) \mathrm{d}(\theta \cdot \nabla K)(y)\right| \\
& \leqslant\|g\|_{\infty}\|\theta \cdot \nabla K\|_{1} \leqslant V(K)\|g\|_{\infty}
\end{aligned}
$$

Lemma 7.10. - For every symmetric convex body $C$, isotropic and normalized by variance, one has that $V\left(K_{C}\right)=q(C)$.

Proof. - Let $\theta$ belong to $S^{n-1}$ and let $y$ be in $\theta^{\perp}$. For each line $y+\mathbb{R} \theta$ that meets the set $C$, the jumps of the density $K_{C}=|C|^{-1} \mathbf{1}_{C}$ of $\mu_{C}$, when traveling on the line in the direction of increasing real numbers, are equal to $|C|^{-1}$ when we enter $C$, and to $-|C|^{-1}$ when leaving $C$, implying that the mass of the directional derivative is equal to $2 /|C|$ times the measure of the projection of $C$ onto $\theta^{\perp}$. More precisely, suppose without loss of generality that $\theta$ is the first basis vector $\mathbf{e}_{1}$ of $\mathbb{R}^{n}$ and let $\pi_{1}$ be the orthogonal projection onto $\mathbf{e}_{1}^{\perp}$. Let $\psi \in \mathcal{S}\left(\mathbb{R}^{n}\right)$ be given, and write each $x \in \mathbb{R}^{n}$ as $x=(s, y)$ with $s \in \mathbb{R}$ and $y \in \mathbb{R}^{n-1}$. Using Fubini, we get

$$
\begin{aligned}
-\left\langle\mathbf{e}_{1} \cdot \nabla \mu_{C}, \psi\right\rangle & =-\left\langle\frac{\partial \mu_{C}}{\partial x_{1}}, \psi\right\rangle=\left\langle\mu_{C}, \frac{\partial \psi}{\partial x_{1}}\right\rangle \\
& =\int_{\mathbb{R}^{n-1}}\left(\int_{\mathbb{R}}|C|^{-1} \mathbf{1}_{C}(s, y) \frac{\partial \psi}{\partial x_{1}}(s, y) \mathrm{d} s\right) \mathrm{d} y
\end{aligned}
$$

The inside integral is 0 if $L_{y}=y+\mathbb{R} \mathbf{e}_{1}$ does not meet the convex set $C$. Otherwise, the line $L_{y}$ cuts $C$ along a segment $\left[y+s_{1}(y) \mathbf{e}_{1}, y+s_{2}(y) \mathbf{e}_{1}\right]$, $s_{1}(y) \leqslant s_{2}(y)$, and
$-\left\langle\mathbf{e}_{1} \cdot \nabla \mu_{C}, \psi\right\rangle=\frac{1}{|C|_{n}} \int_{\pi_{1} C}\left(\psi\left(s_{2}(y), y\right)-\psi\left(s_{1}(y), y\right)\right) \mathrm{d} y \leqslant \frac{2\left|\pi_{1} C\right|_{n-1}}{|C|_{n}}\|\psi\|_{\infty}$.
Going back to a general $\theta \in S^{n-1}$ and according to (7.1), we conclude that

$$
\left\|\theta \cdot \nabla \mu_{C}\right\|_{1} \leqslant \frac{2}{|C|_{n}}\left|P_{\theta} C\right|_{n-1} \leqslant 2 Q\left(C_{0}\right) L\left(C_{0}\right)=q(C)
$$

We get $V\left(K_{C}\right) \leqslant q(C)$, which suffices for our purpose. Müller [59, Lemma 3] shows that this inequality is actually an equality.

When $K_{g}=K_{C}$, we have $\left\|U_{K_{C}}\right\|_{\infty \rightarrow \infty} \leqslant q(C)$, specifying the estimate (7.21) obtained in the general case. We complete now the interpolation for $U_{K_{g}}$. We formulate the next Lemma so that we can apply it again in Section 8.

Lemma 7.11. - Let $K$ be an isotropic log-concave probability density on $\mathbb{R}^{n}$ with variance $\sigma^{2}$. For $2 \leqslant q \leqslant+\infty$, one has that

$$
\left\|U_{K} f\right\|_{q}=\|\nabla K * f\|_{q} \leqslant 2^{1 / q} \sigma^{-2 / q} V(K)^{1-2 / q}\|f\|_{q}, \quad f \in L^{q}\left(\mathbb{R}^{n}\right)
$$

If $K_{g}$ is an integrable kernel on $\mathbb{R}^{n}$ satisfying (6.1.H), then

$$
\left\|U_{K_{g}}\right\|_{q \rightarrow q} \leqslant\left(2 \pi \delta_{0, g}\right)^{2 / q} V\left(K_{g}\right)^{1-2 / q}
$$

Proof. - Let $m=\widehat{K}$ and consider first $q=2$. By Parseval (2.12.P) we have

$$
\|\nabla K * f\|_{2}^{2}=\left\|\left(\sum_{j=1}^{n}\left|\partial_{j} K * f\right|^{2}\right)^{1 / 2}\right\|_{2}^{2}=\int_{\mathbb{R}^{n}} \sum_{j=1}^{n} 4 \pi^{2} \xi_{j}^{2}|m(\xi)|^{2}|\widehat{f}(\xi)|^{2} \mathrm{~d} \xi,
$$

and $\sum_{j=1}^{n} 4 \pi^{2} \xi_{j}^{2}|m(\xi)|^{2}=4 \pi^{2}|\xi|^{2}|m(\xi)|^{2} \leqslant 2 \sigma^{-2}$ by (5.17.B) (or by (6.1.H), it is $\leqslant 4 \pi^{2} \delta_{0, g}^{2}$ ), hence $\left\|U_{K}\right\|_{2 \rightarrow 2} \leqslant \sqrt{2} \sigma^{-1}$ (or $\leqslant 2 \pi \delta_{0, g}$ ). If $q \in(2,+\infty)$, we write $1 / q=(1-\theta) / 2+\theta / \infty$ with $\theta=1-2 / q$. We get that $\left\|U_{K}\right\|_{q \rightarrow q} \leqslant$ $\left(\sqrt{2} \sigma^{-1}\right)^{2 / q} V(K)^{1-2 / q}$ (or we get $\leqslant\left(2 \pi \delta_{0, g}\right)^{2 / q} V(K)^{1-2 / q}$ ) by Lemma 7.9 and interpolation $\left(L^{2}, L^{\infty}\right)$.

End of the proof of Lemma 7.7. - We use Lemma 7.8, then apply Lemma 7.11 to $K_{g}$ with $1 / q=1-1 / p$ and obtain that

$$
\left\|T_{m_{g}^{\#}}\right\|_{p \rightarrow p} \leqslant(2 \pi)^{-1} \rho_{p}\left\|U_{K_{g}}\right\|_{q \rightarrow q} \leqslant(2 \pi)^{1-2 / p} \rho_{p} \delta_{0, g}^{2-2 / p} V\left(K_{g}\right)^{-1+2 / p} .
$$

### 7.5.1. Conclusion

We finish the proof of Proposition 7.2. We first run over half of the way in the following lemma, which we shall refer to again in Section 8.

Lemma 7.12. - Let $K_{g}$ be an integrable kernel on $\mathbb{R}^{n}$ satisfying $\left(7.5 . \mathbf{H}_{\infty}\right)$, and let $m_{g}^{\#}$ be defined by (7.19). Let $\alpha \in(0,1)$ and suppose that $1<p_{0}<p<2$. There exists a constant $\kappa\left(p, p_{0}\right)$, independent of $n$, such that

$$
\begin{aligned}
& \left\|(\xi \cdot \nabla)^{\alpha} m_{g}(\xi)\right\|_{p \rightarrow p} \\
& \quad \leqslant\left\|K_{g}\right\|_{L^{1}\left(\mathbb{R}^{n}\right)}+\frac{\kappa\left(p, p_{0}\right)}{1-\alpha}\left(\left\|K_{g}\right\|_{L^{1}\left(\mathbb{R}^{n}\right)}+\left\|m_{g}^{\#}\right\|_{p_{0} \rightarrow p_{0}}\right)^{1-\theta} \Delta_{k(\theta)}^{\theta},
\end{aligned}
$$

where $\theta \in(0,1)$ is defined by $1 / p=(1-\theta) / p_{0}+\theta / 2$ and $k(\theta)=\lceil 1 / \theta\rceil$.
Proof. - Lemma 7.3 gives

$$
\left\|(\xi \cdot \nabla)^{\alpha} m_{g}(\xi)\right\|_{p \rightarrow p} \leqslant\left\|K_{g}\right\|_{L^{1}\left(\mathbb{R}^{n}\right)}+\left\|\left.d_{t}^{\alpha} m_{g}(t \xi)\right|_{t=1}\right\|_{p \rightarrow p}
$$

Let $\varepsilon=1-\alpha>0$. We apply complex interpolation to the Müller family $\left(m_{z}^{\varepsilon}\right)$ in the strip $S=\{z \in \mathbb{C}:-\varepsilon \leqslant z \leqslant \nu\}$ of width $w:=\nu+\varepsilon$. We bound $m_{\alpha}^{\varepsilon}(\xi)=\left.d_{t}^{\alpha} m_{g}(t \xi)\right|_{t=1}$ on $L^{p}\left(\mathbb{R}^{n}\right)$, using $L^{p_{0}}$ estimates of $m_{z}^{\varepsilon}$ for $\operatorname{Re} z=-\varepsilon$ and $L^{2}$ estimates when $\operatorname{Re} z=\nu$. The value $\nu$ must satisfy $\alpha=(1-\theta)(-\varepsilon)+\theta \nu$, hence $\nu=1 / \theta-\varepsilon>1-\varepsilon$ and $w=1 / \theta$. It follows from Lemma 7.6 that

$$
\left\|m_{-\varepsilon+\mathrm{i} \tau}^{\varepsilon}\right\|_{p_{0} \rightarrow p_{0}} \leqslant 2 \varepsilon^{-1} \lambda_{p_{0}} \mathrm{e}^{\pi|\tau|}\left(\left\|K_{g}\right\|_{L^{1}\left(\mathbb{R}^{n}\right)}+\left\|m_{g}^{\#}\right\|_{p_{0} \rightarrow p_{0}}\right) .
$$

By Lemma 7.4, each operator $m_{\nu+\mathrm{i} \tau}^{\varepsilon}$ is bounded by

$$
\kappa_{\nu} \varepsilon^{-1} \Delta_{\ell}\left(1+\tau^{2}\right)^{\nu / 2-1 / 4} \mathrm{e}^{\pi|\tau| / 2}
$$

on $L^{2}\left(\mathbb{R}^{n}\right)$, with $\kappa_{\nu} \leqslant \kappa_{w}=4 \Gamma(\max (w, 2))(3 / 2)^{w-1 / 2} \mathrm{e}^{\pi / 4}=: c_{\theta}$, a function of $\theta$ alone, and $\ell=\lceil\nu+\varepsilon\rceil=\lceil w\rceil$. If we check the admissible growth condition in $S$, we can rely on Corollary 3.4, Remark 3.6 and (3.22) in order to get a bound

$$
\left\|m_{\alpha}^{\varepsilon}(\xi)\right\|_{p \rightarrow p} \leqslant \kappa\left(p, p_{0}\right) \varepsilon^{-1}\left(\left\|K_{g}\right\|_{L^{1}\left(\mathbb{R}^{n}\right)}+\left\|m_{g}^{\#}\right\|_{p_{0} \rightarrow p_{0}}\right)^{1-\theta} \Delta_{k(\theta)}^{\theta}
$$

with $k(\theta)=\ell=\lceil 1 / \theta\rceil$ and $\kappa\left(p, p_{0}\right) \leqslant(1+w / 2)^{\theta(\nu-1 / 2)} \mathrm{e}^{\pi w / 2}\left(2 \lambda_{p_{0}}\right)^{1-\theta} c_{\theta}^{\theta}$. Observing that $\theta \nu<1, w>1$ and $\lambda_{p_{0}} \geqslant 1$, we may simplify this bound as

$$
\begin{equation*}
\kappa\left(p, p_{0}\right) \leqslant \kappa w \mathrm{e}^{\pi w / 2} \lambda_{p_{0}} \Gamma(\max (w, 2))^{1 / w} \leqslant \kappa w^{2} \mathrm{e}^{\pi w / 2} \lambda_{p_{0}} . \tag{7.22}
\end{equation*}
$$

We now verify that the holomorphic function $H(z)=\left\langle m_{z}^{\varepsilon} f_{z}, g_{z}\right\rangle$ has an admissible growth in $S$. Since the kernels are bounded functions of $\xi$ by Lemma 7.4, all multipliers $m_{z}^{\varepsilon}, z \in S$, are $L^{2}$-bounded with a bound of the form $\kappa \mathrm{e}^{\kappa|\operatorname{Im} z|}$. If we restrict to functions $f$ and $g$ bounded with bounded support, we have by (3.24) that $f_{z}, g_{z}$ are uniformly bounded in $L^{2}\left(\mathbb{R}^{n}\right)$, and we can conclude as in Section 7.3.

End of the proof of Proposition 7.2. - Given $p \in(1,2)$ and $\alpha=1-\varepsilon \in$ $(0,1)$, we select $p_{0} \in(1, p)$ and let $\theta \in(0,1)$ satisfy $1 / p=(1-\theta) / p_{0}+\theta / 2$. Since $1<p_{0}<p<2$, we have that $0<\theta<2(1-1 / p)<1$. It follows from Lemma 7.7 that

$$
\left\|m_{g}^{\#}\right\|_{p_{0} \rightarrow p_{0}} \leqslant(2 \pi)^{1-2 / p_{0}} \rho_{p_{0}} \delta_{0, g}^{2-2 / p_{0}} V\left(K_{g}\right)^{2 / p_{0}-1} .
$$

By Lemma 7.12 , and because $\left\|K_{g}\right\|_{L^{1}\left(\mathbb{R}^{n}\right)} \leqslant 1, \rho_{p_{0}} \geqslant 1$ (see Remark 2.3), we get
$\left\|(\xi \cdot \nabla)^{\alpha} m_{g}(\xi)\right\|_{p \rightarrow p} \leqslant 1+\kappa\left(p, p_{0}\right) \varepsilon^{-1} \rho_{p_{0}} \Delta_{k(\theta)}^{\theta}\left(1+\delta_{0, g}^{1-\theta-(2 / p-1)} V\left(K_{g}\right)^{2 / p-1}\right)$.
We still have a choice of $\theta \in(0,2-2 / p)$. If $\theta$ gets small, then the power of $\Delta_{k(\theta)}$ gets small, but the number $k(\theta)$ of constants $\delta_{j, g}$ involved increases to infinity. In the log-concave case, the estimate (5.18) indicates a growth of order $\Delta_{k, c} \sim k$ ! yielding $\Delta_{k(\theta)}^{\theta} \sim 1 / \theta$. Furthermore, the width $w=1 / \theta$ of the strip and the associated interpolation constants also tend to $+\infty$ in this case, and we should thus keep $\theta$ away from 0 , as much as possible. If $\theta$ approaches its upper limit $2(1-1 / p)$, then $p_{0}$ tends to 1 and the constants such as $\lambda_{p_{0}}, \rho_{p_{0}}$ tend to infinity. Choosing $\theta=(4 / 3)(1-1 / p)$ has the merit to provide the relatively simple bound

$$
\begin{align*}
& \left\|(\xi \cdot \nabla)^{\alpha} m_{g}(\xi)\right\|_{p \rightarrow p} \\
& \quad \leqslant 1+\kappa(\alpha, p) \Delta_{k(p)}^{(4 / 3)(1-1 / p)}\left(1+\delta_{0, g}^{(2 / 3)(1-1 / p)} V\left(K_{g}\right)^{2 / p-1}\right) \tag{7.23}
\end{align*}
$$

with $\kappa(\alpha, p)=\kappa\left(p, p_{0}\right) \varepsilon^{-1} \rho_{p_{0}} \leqslant \kappa \varepsilon^{-1} w^{2} \mathrm{e}^{\pi w / 2} \lambda_{p_{0}} \rho_{p_{0}}$ by (7.22), with $p_{0}-$ $1=(p-1) /(5-2 p)$ and $k(p)=\lceil 1 / \theta\rceil=\lceil 3 p /(4 p-4)\rceil$.

Remark 7.13. - It is usual to have a factor $1 / \Gamma(z)$ in fractional derivatives, which led us to seeing $\mathrm{e}^{\pi|\tau| / 2}$ in many places, ending with $\mathrm{e}^{\pi w / 2}$ in our estimate (7.22) of $\kappa\left(p, p_{0}\right)$, with $w=1 / \theta \sim q:=p /(p-1)$ after the final choice of $\theta=4 /(3 q)$ above. We could avoid this exponential though. Consider the modified Müller family

$$
\widetilde{m}_{z}^{\varepsilon}(\xi)=\Gamma(1+\varepsilon+z) \frac{\Gamma(2 \varepsilon+z)}{\Gamma(1+\varepsilon)} m_{z}^{\varepsilon}(\xi), \quad \operatorname{Re} z \geqslant-\varepsilon, \xi \in \mathbb{R}^{n}
$$

which coincides with the former at $z=\alpha$ since $\varepsilon+\alpha=1$ and $\Gamma(2)=1$. For the $L^{p_{0}}$ bound when $z=-\varepsilon+\mathrm{i} \tau$, we decompose $\widetilde{m}_{-\varepsilon+\mathrm{i} \tau}^{\varepsilon}(\xi)$ as

$$
\frac{1}{\Gamma(1+\varepsilon)}\left[\Gamma(1+\mathrm{i} \tau)(1+|\xi|)^{-\mathrm{i} \tau}\right]\left[\frac{\Gamma(\varepsilon+\mathrm{i} \tau)}{\Gamma(\varepsilon-\mathrm{i} \tau)}(1+|\xi|) \int_{1}^{2}(s-1)^{\varepsilon-\mathrm{i} \tau-1} m_{g}(s \xi) \mathrm{d} s\right]
$$

Introducing $\Gamma(1+\mathrm{i} \tau)$ in the Laplace type multiplier (7.18), we obtain a new function $\widetilde{a}_{\tau}(t)$ bounded by 1 , and $\Gamma(2 \varepsilon+z)$ is used for the bound of $\left.d_{t}^{-\varepsilon+\mathrm{i} \tau} m_{g}(t \xi)\right|_{t=1}$ because $|\Gamma(\varepsilon+\mathrm{i} \tau)|=|\Gamma(\varepsilon-\mathrm{i} \tau)|$. We get in this way for $\widetilde{m}_{z}^{\varepsilon}(\xi)$ a bound

$$
\left\|\widetilde{m}_{z}^{\varepsilon}(\xi)\right\|_{p_{0} \rightarrow p_{0}} \leqslant 2 \varepsilon^{-1} \lambda_{p_{0}}\left(\left\|K_{g}\right\|_{L^{1}\left(\mathbb{R}^{n}\right)}+\left\|m^{\#}\right\|_{p_{0} \rightarrow p_{0}}\right)
$$

that replaces Lemma 7.6 (we use again $1 / \Gamma(1+\varepsilon) \leqslant 2$ ). The $L^{\infty}$ bounds obtained in (7.16) when $z=k-\varepsilon+\mathrm{i} \tau, k>0$, have now a largest factor of $\Delta_{k} \varepsilon^{-1}$ equal to $\Gamma(k+1+\mathrm{i} \tau) \Gamma(k+\varepsilon+\mathrm{i} \tau) /[\Gamma(1+\varepsilon) \Gamma(-k+1+\varepsilon-\mathrm{i} \tau)]$ (when the index $j_{1}$ in (7.16) is equal to $k-1$ ). The modulus of this factor is the same as that of

$$
\frac{\Gamma(k+1+\mathrm{i} \tau) \Gamma(k+\varepsilon+\mathrm{i} \tau)}{\Gamma(1+\varepsilon) \Gamma(-k+1+\varepsilon+\mathrm{i} \tau)}=\frac{\Gamma(1+\mathrm{i} \tau)}{\Gamma(1+\varepsilon)}\left(\prod_{j=1}^{k}(j+\mathrm{i} \tau)\right)\left(\prod_{j=-k+1}^{k-1}(j+\varepsilon+\mathrm{i} \tau)\right)
$$

This is a bounded function of $\tau$ according to (3.1), with a rough bound given by $2 \sqrt{2 \pi}(k+|\tau|)^{3 k} \mathrm{e}^{-\pi|\tau| / 2} \leqslant 6.2^{k} k^{3 k}$ (use $x / \sinh x \leqslant(1+2 x) \mathrm{e}^{-x}$ for $x \geqslant 0)$. One need not be too careful here since this term will be raised to the power $\theta=1 / w \lesssim 1 / k$. We use it as in (7.17) for two values $k, k+1$ such that $k \leqslant \nu_{0}+\varepsilon \leqslant k+1 \leqslant \ell=\lceil w\rceil<w+1$. One has then for the $L^{2}$ bound of $\widetilde{m}_{\nu+\mathrm{i} \tau}^{\varepsilon}(\xi)$ an estimate by $\kappa 2^{w} w^{3(w+1)} \Delta_{\ell} \varepsilon^{-1}$. By interpolation we have

$$
\left\|m_{\alpha}^{\varepsilon}(\xi)\right\|_{p \rightarrow p} \leqslant \kappa \varepsilon^{-1} \lambda_{p_{0}}^{1-\theta}\left(2^{w} w^{3(w+1)}\right)^{\theta}\left(\left\|K_{g}\right\|_{L^{1}\left(\mathbb{R}^{n}\right)}+\left\|m_{g}^{\#}\right\|_{p_{0} \rightarrow p_{0}}\right)^{1-\theta} \Delta_{\ell}^{\theta}
$$

We thus get for $\kappa\left(p, p_{0}\right)$ in (7.22) a new estimate $\kappa^{\prime}\left(p, p_{0}\right) \leqslant \kappa w^{3} \lambda_{p_{0}}$, leading in (7.23) to $\kappa^{\prime}(\alpha, p) \leqslant \kappa \varepsilon^{-1} q^{3} \lambda_{p_{0}} \rho_{p_{0}}$. The final choice in the proof of Proposition 7.2 gives $p_{0}-1=(p-1) /(5-2 p)$ of order $p-1$ as $p \rightarrow 1$, and since $\lambda_{p_{0}}, \rho_{p_{0}}$ are $O\left(\left(p_{0}-1\right)^{-1}\right)$ as $p_{0} \rightarrow 1$ (see (2.20) and (2.24)), we end up with $\kappa^{\prime}(\alpha, p) \leqslant \kappa \varepsilon^{-1} q^{5}$, a bound which is polynomial but has no reason
to be accurate. After these modifications, we have for Proposition 7.2 when $1<p \leqslant 2$ a new form

$$
\begin{align*}
& \left\|(\xi \cdot \nabla)^{\alpha} m_{g}(\xi)\right\|_{p \rightarrow p} \\
& \quad \leqslant 1+\kappa \varepsilon^{-1} q^{5} \Delta_{k(p)}^{(4 / 3)(1-1 / p)}\left(1+\delta_{0, g}^{(2 / 3)(1-1 / p)} V\left(K_{g}\right)^{2 / p-1}\right) \tag{7.24}
\end{align*}
$$

with $k(p)=\lceil 3 p /(4 p-4)\rceil$ and $q=p /(p-1)$.
Remark 7.14. - With the new information above, we go back to the proof of Theorem 7.1. One has chosen there $\varepsilon=1-\alpha=3(p-1) /(4 p+4)$, and $p_{0} \in(1, p)$ such that $p_{0}-1=3(p-1) /(5-p)$. Both $\varepsilon$ and $p_{0}-1$ behave as multiples of $p-1$ when $p \rightarrow 1$. If we consider the Poisson kernel $P$ as another $K_{g}$ satisfying (7.5.H $\boldsymbol{H}_{\infty}$ ), and with $V(P) \leqslant 2 / \pi$ by (7.4), we can apply to it (7.24) for the value $p_{0}$ and obtain that $\left\|(\xi \cdot \nabla)^{\alpha} \widehat{P}(\xi)\right\|_{p_{0} \rightarrow p_{0}} \leqslant \kappa \varepsilon^{-1} q^{5}$. Applying also (7.24) to $m_{g}$ and $p_{0}$, we get for $m=m_{g}-\widehat{P}$ that

$$
\left\|(\xi \cdot \nabla)^{\alpha} m(\xi)\right\|_{p_{0} \rightarrow p_{0}} \leqslant \kappa q^{6} \Delta_{k\left(p_{0}\right)}^{(4 / 3)\left(1-1 / p_{0}\right)} \delta_{0, g}^{(2 / 3)\left(1-1 / p_{0}\right)}\left(1+\lambda^{2 / p_{0}-1}\right)=: B
$$

We have $\alpha-1 / p_{0}=3(p-1) /(4 p+4)$ which again is of order $p-1$. The constant $\kappa_{\alpha, p_{0}}$ from (6.25), seen in (7.7), behaves thus as $(p-1)^{-1 / p_{0}} \simeq q$, so $C_{p_{0}}^{\prime \prime}(\lambda)$ is bounded by $\kappa q(2+B)$. Also $r_{0}-1 \simeq\left(p+p_{0}-2\right) / 2$ in (7.8) is of order $p-1 \simeq 1 / q$. In (7.8), the constants $C_{r_{0}}$ and $C_{p}^{\prime}$ are of order $q$. Indeed, we can take $C_{r_{0}}=\mathrm{q}_{r_{0}}$ from (2.4), that was estimated by $r_{0} /\left(r_{0}-1\right)$ in (2.5), and $C_{p}^{\prime}$ can be the bound for the maximal function of the Poisson kernel, see $\left(1.31 . P^{*}\right)$. Also, $1-\gamma=(p-1) /(2 p)$, and with Lemma 6.19 we know that

$$
\sum_{k \in \mathbb{Z}} a_{\alpha, k}^{(1-\gamma) p / 2} \leqslant \kappa \sum_{k \in \mathbb{Z}} 2^{-(1-\alpha)(1-\gamma) p|k| / 2} \leqslant \frac{\kappa^{\prime}}{(1-\alpha)(1-\gamma)} \leqslant \kappa^{\prime \prime} q^{2}
$$

Finally, we obtain for Theorem 7.1 another bizarre polynomial estimate

$$
\begin{aligned}
\left\|\mathrm{M}_{K_{g}}\right\|_{p \rightarrow p} & \leqslant\left\|M_{K}\right\|_{p \rightarrow p}+O(q) \leqslant C_{r_{0}}^{2} C_{p_{0}}^{\prime \prime}(\lambda)\left(\kappa^{\prime \prime} q^{2}\right)^{2}+O(q) \\
& \leqslant \kappa q^{13} \Delta_{\lceil q\rceil}^{1-1 / p} \Delta_{1}^{1-1 / p}\left(1+\lambda^{2 / p-1}\right), \quad 1<p \leqslant 2 .
\end{aligned}
$$

## 8. Bourgain's article on cubes

In this section, $Q$ is a cube in dimension $n$, more precisely, the symmetric cube

$$
Q=Q_{n}=\left[-\frac{1}{2}, \frac{1}{2}\right]^{n}
$$

of volume 1 in $\mathbb{R}^{n}$. It is isotropic, but if we look for a multiple $b Q$ normalized by variance, we would need that the half-side $a=b / 2$ of $b Q$ satisfy $\sigma_{b Q}^{2}=1$, where

$$
\sigma_{b Q}^{2}=\frac{1}{|b Q|} \int_{b Q} x_{1}^{2} \mathrm{~d} x=\frac{1}{2 a} \int_{-a}^{a} s^{2} \mathrm{~d} s=\frac{1}{a} \int_{0}^{a} s^{2} \mathrm{~d} s=\frac{a^{2}}{3}
$$

and where $x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n}$. This gives $a=\sqrt{3}$ in every dimension $n$, but the cube $[-\sqrt{3}, \sqrt{3}]^{n}$ is not very pleasant to manipulate, and we shall rather follow Bourgain [13] and keep the volume 1 cube $Q$. With $a=1 / 2$, the covariance for $Q$ is given by $(12)^{-1} \mathrm{I}_{n}$. Since the variance $\sigma_{Q}^{2}=1 / 12$ is independent of the dimension, we shall have no problem with the estimates (5.17.B) or (5.19). The Fourier transform of the probability measure $\mu_{Q}$ is given by

$$
m_{Q}(\xi)=\widehat{\mu_{Q}}(\xi)=\widehat{K_{Q}}(\xi)=\prod_{j=1}^{n} \frac{\sin \left(\pi \xi_{j}\right)}{\pi \xi_{j}}, \quad \xi=\left(\xi_{1}, \ldots, \xi_{n}\right) \in \mathbb{R}^{n}
$$

Bourgain observes that a decay better than the usual (5.17.B) for a Fourier transform $m_{C}$ would allow to relax the limitation $p>3 / 2$ of Theorem 6.2, and that this better decay is achieved by $m_{Q}$ is most directions. He says that his proof proceeds therefore to diverse localizations in Fourier space.

Theorem 8.1 (Bourgain [13]). - For every $p$ in $(1,+\infty]$, there exists a constant $\kappa_{p}$ such that $\left\|\mathrm{M}_{Q_{n}}\right\|_{p \rightarrow p} \leqslant \kappa_{p}$ for every integer $n \geqslant 1$.

We shall approach the maximal function problem for the cube by summing expressions such as $K^{R}-K^{2 R}$, with

$$
K^{R}=K_{Q} * G_{(1 / R)}
$$

where $G$ is a Gaussian probability kernel, $G_{(1 / R)}$ its dilate (2.7), and where $R$ takes the values $1,2, \ldots, 2^{j}, \ldots$ with $j$ being any integer $\geqslant 0$. This is a Littlewood-Paley-type decomposition, similar to what we have seen before. By Prékopa-Leindler, $K^{R}$ is a log-concave probability density. We shall set $m^{R}=\widehat{K^{R}}$ in what follows.

We will call the Carbery-Müller artillery and obtain when $1<r<2$, for every $\delta>0$ and $R=2^{j}$ with $j \geqslant 0$, bounds of the form

$$
\left\|\sup _{1 \leqslant t \leqslant 2}\left|K_{(t)}^{R} * f\right|\right\|_{r} \leqslant \kappa_{\delta, r} R^{\delta}\|f\|_{r}, \quad \text { where } K_{(t)}^{R}:=\left(K^{R}\right)_{(t)}
$$

Why this may be a decisive step will be explained below. According to Carbery's Proposition $6.14(2)$, this bound will be consequence of the $L^{r}\left(\mathbb{R}^{n}\right)$ boundedness of the multiplier $(\xi \cdot \nabla)^{\alpha} m^{R}(\xi)$ for a value of $\alpha \in(1 / r, 1)$. Next, following Müller, it will be enough to estimate in $L^{s}\left(\mathbb{R}^{n}\right)$, with $1<s<r$, the "crucial" multiplier $|\xi| m^{R}(\xi)$. This is what Bourgain does along several
pages, in a series of reductions bringing in many tools that are specific to the product structure of the cube.

### 8.1. Holding on Müller and Carbery

Let us specify the preceding rough outline. The final objective is to bound in $L^{p}\left(\mathbb{R}^{n}\right)$ the maximal operator $\mathrm{M}_{Q}$ for $p$ below the limit $3 / 2$ that is known so far, proving that

$$
\left\|\sup _{t>0}\left|\left(K_{Q}\right)_{(t)} * f\right|\right\|_{p} \leqslant \kappa_{p}\|f\|_{p}, \quad 1<p<2, \quad f \in L^{p}\left(\mathbb{R}^{n}\right)
$$

We fix a value $p \in(1,2)$ in all that follows. In order to obtain the property $\left(A_{\mathbf{2}}\right)$, needed for applying Carbery's Proposition 6.6, we must show that

$$
\begin{equation*}
\left\|\sup _{1 \leqslant t \leqslant 2}\left|K_{(t)} * f\right|\right\|_{p_{2}} \leqslant \kappa\|f\|_{p_{2}}, \quad 1<p_{2}<p<2 \tag{A}
\end{equation*}
$$

where $K=K_{Q}-P$ and $P$ is the Poisson kernel (1.32). This is the only missing fact for lowering the limitation $p>3 / 2$ down to $p>1$, as explained in the proof of Müller's Theorem 7.1. For the Poisson side it is fine, it remains to work on $K_{Q}$. We introduce the Gaussian kernel $G=\left(\gamma_{n}\right)_{(\sqrt{2} / \sqrt{\pi})}$ on $\mathbb{R}^{n}$. The variance of $G$ is equal to $2 / \pi$, thus independent of $n$, and $\widehat{G}(\xi)=\mathrm{e}^{-4 \pi|\xi|^{2}}$ for every $\xi \in \mathbb{R}^{n}$. With this normalization for $G$, we have by (7.3) that

$$
\begin{equation*}
V(G)=\sqrt{\pi / 2} V\left(\gamma_{n}\right)=1 \tag{8.1}
\end{equation*}
$$

We decompose the Dirac probability measure $\delta_{0}$ at the origin, in the sense of distributions, by means of the simple telescopic series

$$
\delta_{0}=G_{(1)}+\left(G_{(1 / 2)}-G_{(1)}\right)+\cdots+\left(G_{\left(2^{-k-1}\right)}-G_{\left(2^{-k}\right)}\right)+\cdots
$$

and we decompose $K_{Q}$ accordingly, using the approximations $K^{R}=K_{Q} *$ $G_{(1 / R)}$, for $R=2^{j} \geqslant 1$ and $j$ nonnegative integer, under the form

$$
K_{Q}=K^{1}+\left(K^{2}-K^{1}\right)+\cdots+\left(K^{2^{j+1}}-K^{2^{j}}\right)+\cdots
$$

By Prékopa-Leindler, each $K^{R}$ is a log-concave symmetric probability density on $\mathbb{R}^{n}$. It is isotropic, with a variance $\sigma_{R}^{2}$ satisfying

$$
\begin{equation*}
12^{-1}<\sigma_{R}^{2}=12^{-1}+2 \pi^{-1} R^{-2}<1, \quad R \geqslant 1 \tag{8.2}
\end{equation*}
$$

We set $\mathrm{d} \mu^{R}(x)=K^{R}(x) \mathrm{d} x, m^{R}=\widehat{K^{R}}=\widehat{\mu^{R}}$. It follows from (5.19) that $m^{R}$ satisfies $\left(7.5 . \mathbf{H}_{\infty}\right)$ with constants independent of $n$. We get

$$
\begin{equation*}
\left|\frac{\mathrm{d}^{j}}{\mathrm{~d} t^{j}} m^{R}(t \theta)\right| \leqslant \frac{\delta_{j, c}}{1+\pi|t| / \sqrt{3}} \leqslant \frac{\delta_{j, c}}{1+|t|}, \quad \theta \in S^{n-1}, \quad t \in \mathbb{R}, \quad j \geqslant 0 \tag{8.3}
\end{equation*}
$$

We shall obtain the desired estimate (A) for $p_{2}$ by interpolating between $p_{1}$ and 2 , where $1<p_{1}<p_{2}<p<2$. As we have said previously, we will show that for every $\delta>0$, we have for all $R=2^{j} \geqslant 1$ that

$$
\begin{equation*}
\left\|\sup _{1 \leqslant t \leqslant 2}\left|K_{(t)}^{R} * f\right|\right\|_{p_{1}} \leqslant \kappa_{\delta} R^{\delta}\|f\|_{p_{1}}, \quad f \in L^{p_{1}}\left(\mathbb{R}^{n}\right) \tag{B}
\end{equation*}
$$

and on the other hand, we prove that for every $f \in L^{2}\left(\mathbb{R}^{n}\right)$ we have

$$
\begin{align*}
\left\|\sup _{t>0}\left|\left(K^{R}-K^{2 R}\right)_{(t)} * f\right|\right\|_{2} & \leqslant \kappa R^{-1 / 2}\|f\|_{2}, \quad\left\|\sup _{t>0}\left|K_{(t)}^{1} * f\right|\right\|_{2}  \tag{C}\\
& \leqslant \kappa\|f\|_{2}
\end{align*}
$$

The second inequality in $(\mathbf{C})$ is the log-concave version of Bourgain's $L^{2}$ theorem, Theorem 5.2. One can obtain the first part of $(\mathbf{C})$ by the $\Gamma_{B}(K)$ criterion, Lemma 5.14. We have indeed, uniformly in $\theta \in S^{n-1}$ and in the dimension $n$ (observe that $\widehat{G}$ has a radial expression independent of $n$ ), that

$$
\begin{aligned}
|\widehat{G}(u \theta)-\widehat{G}(2 u \theta)| & \lesssim u^{2} \wedge \mathrm{e}^{-4 \pi u^{2}} \leqslant|u| \wedge|u|^{-1} \quad \text { and } \\
|\theta \cdot \nabla \widehat{G}(u \theta)-\theta \cdot \nabla \widehat{G}(2 u \theta)| & \lesssim|u|\left(1 \wedge \mathrm{e}^{-4 \pi u^{2}}\right) \leqslant 1 \wedge|u|^{-1}, \quad u \in \mathbb{R} .
\end{aligned}
$$

We apply Lemma 6.18 with $K_{1}=K_{Q}, K_{2}=G$, replacing $2^{|k|}$ with $R$ and obtaining

$$
\sum_{j \in \mathbb{Z}}\left(\alpha_{j}\left(K^{R}\right)+\sqrt{\alpha_{j}\left(K^{R}\right) \beta_{j}\left(K^{R}\right)}\right) \lesssim R^{-1 / 2}
$$

If $\delta>0$ is sufficiently small we deduce by interpolation between (B) and $(\mathbf{C})$ that there exists $\delta_{1}>0$ such that

$$
\left\|\sup _{1 \leqslant t \leqslant 2}\left|\left(K_{(t)}^{R}-K_{(t)}^{2 R}\right) * f\right|\right\|_{p_{2}} \leqslant \kappa_{\delta} R^{-\delta_{1}}\|f\|_{p_{2}}
$$

and we get Property (A) by summing on the values $R=2^{j}$ for all integers $j \geqslant 0$. We fix thus a value $\delta_{*}=\delta_{*}\left(p, p_{2}, p_{1}\right)>0$ of $\delta$, sufficiently small for implying that $\delta_{1}>0$ whenever $0<\delta \leqslant \delta_{*}$. Precisely, if $\lambda \in(0,1)$ is such that

$$
\frac{1}{p_{2}}=\frac{1-\lambda}{p_{1}}+\frac{\lambda}{2}
$$

we need to choose $\delta_{*}>0$ so that $-\delta_{1}=(1-\lambda) \delta_{*}-\lambda / 2<0$, i.e., we select a value $\delta_{*}=\delta_{*}\left(p, p_{2}, p_{1}\right)$ such that $0<\delta_{*}\left(p, p_{2}, p_{1}\right)<\left(p_{2}-p_{1}\right) /\left(2 p_{1}-p_{2} p_{1}\right)$.

For obtaining ( $\mathbf{B}$ ) we shall use the conclusion (2) of Carbery's Proposition 6.14 and also apply Müller's analysis. We need to show that for some $\alpha \in\left(1 / p_{1}, 1\right)$ and $0<\delta \leqslant \delta_{*}$, we have

$$
\begin{equation*}
2\left\|m^{R}\right\|_{p_{1} \rightarrow p_{1}}+\left\|(\xi \cdot \nabla)^{\alpha} m^{R}(\xi)\right\|_{p_{1} \rightarrow p_{1}} \leqslant \kappa R^{\delta} \tag{8.4}
\end{equation*}
$$

for all $R=2^{j}, j \in \mathbb{N}$. There is no problem for $m^{R}$, which corresponds to convolution with a probability density, and for the other term we shall apply

Lemma 7.12 with $1<p_{0}<p_{1}<2$. For technical reasons, the value $p_{0}$, close to 1 , is chosen in a way that its conjugate $q_{0}$ is an integer of the form $2^{\nu}$, with $\nu$ integer $>0$. If we can prove that for a fixed $\delta>0$ and for every $R=2^{j}$, we have

$$
\begin{equation*}
\left\||\xi| m^{R}(\xi)\right\|_{p_{0} \rightarrow p_{0}} \leqslant \kappa_{\delta} R^{\delta} \tag{8.5}
\end{equation*}
$$

it follows from Lemma 7.12 that $\left\|(\xi \cdot \nabla)^{\alpha} m^{R}(\xi)\right\|_{p_{1} \rightarrow p_{1}} \leqslant \kappa_{\delta}^{\prime}\left(1+R^{\delta \beta}\right) \leqslant \kappa_{\delta}^{\prime \prime} R^{\delta}$ for some $\beta \in(0,1)$, uniformly in the dimension $n$ according to (8.3). The conclusion (8.4) is then obtained.

By exploiting the inequality (2.22) on Riesz transforms, Müller's plan went on with a reduction to estimating the expression $\left\|\nabla \mu^{R} * g\right\|_{q_{0}}$ when $g \in L^{q_{0}}\left(\mathbb{R}^{n}\right)$ and $1 / q_{0}+1 / p_{0}=1$. We must show that for every $R=2^{j}$ we have

$$
\left\|\nabla \mu^{R} * g\right\|_{q_{0}} \leqslant \kappa_{p_{0}, \delta} R^{\delta}\|g\|_{q_{0}}
$$

yielding (8.5) by Lemma 7.8. We use (8.1) and (7.3), which give

$$
\begin{equation*}
V\left(K^{R}\right)=V\left(\mu_{Q} * G_{(1 / R)}\right) \leqslant V\left(G_{(1 / R)}\right)=R \tag{8.6}
\end{equation*}
$$

By Lemma 7.9 , this bound for the mass of $\theta \cdot \nabla \mu^{R}$ when $\theta \in S^{n-1}$ implies that $\left\|\nabla \mu^{R} * g\right\|_{L^{\infty}\left(\mathbb{R}^{n}\right)} \leqslant R\|g\|_{L^{\infty}\left(\mathbb{R}^{n}\right)}$. Then, by interpolation with the $L^{2}$ case given by $(\mathbf{C})$, we can find when $2<q<+\infty$ a bound in $L^{q}\left(\mathbb{R}^{n}\right)$ of the form

$$
\begin{aligned}
\left\|\nabla\left(\mu^{R}-\mu^{2 R}\right) * g\right\|_{L^{q}\left(\mathbb{R}^{n}\right)} & \leqslant \kappa\left(R^{-1 / 2}\right)^{2 / q} R^{1-2 / q}\|g\|_{L^{q}\left(\mathbb{R}^{n}\right)} \\
& =\kappa R^{1-3 / q}\|g\|_{L^{q}\left(\mathbb{R}^{n}\right)} .
\end{aligned}
$$

This interpolation $\left(L^{\infty}, L^{2}\right)$ does not give the desired bound $R^{\delta}$ in $L^{q_{0}}\left(\mathbb{R}^{n}\right)$, with $\delta$ small, when $q_{0} \geqslant 3$. However, it does give the right ingredient for the Bourgain-Carbery Theorem 6.2 when $3 / 2<p \leqslant 2$, since $1-3 / q<0$ in this case.

For going farther than Müller, one has to prove inequalities that allow one to work in $L^{r}\left(\mathbb{R}^{n}\right), 2<r<+\infty$, instead of $L^{\infty}\left(\mathbb{R}^{n}\right)$. This is done with the help of certain analytic semi-groups (Section 8.2), as well as a ad hoc method a la Bourgain, which he says inspired from martingale techniques (Section 8.3). Theorem 8.1 will be obtained once we have the following proposition, which we can apply with a value $\delta \leqslant \delta_{*}\left(p, p_{2}, p_{1}\right)$. We then conclude by the preceding discussion.

Proposition 8.2. - For every $\delta>0$ and $q_{0}=2^{\nu}$, with $\nu$ an integer $\geqslant 1$, there exists a constant $\kappa\left(q_{0}, \delta\right)$ such that for every $n \geqslant 1$ and $R=2^{k}$, $k=0,1, \ldots$, one has

$$
\left\|\nabla \mu^{R} * g\right\|_{L^{q_{0}}\left(\mathbb{R}^{n}\right)} \leqslant \kappa\left(q_{0}, \delta\right) R^{\delta}\|g\|_{L^{q_{0}}\left(\mathbb{R}^{n}\right)}, \quad g \in L^{q_{0}}\left(\mathbb{R}^{n}\right)
$$

We shall keep $\delta>0, p_{0}=q_{0} /\left(q_{0}-1\right)$ and $R=2^{k_{0}}$ fixed in the rest of Section 8.

### 8.1.1. A priori estimate

The proof will play with an a priori estimate

$$
\begin{equation*}
\left\|\nabla \mu^{R} * g\right\|_{L^{q_{0}}\left(\mathbb{R}^{n}\right)} \leqslant B\left(q_{0}, R, n\right)\|g\|_{L^{q_{0}}\left(\mathbb{R}^{n}\right)}, \quad g \in L^{q_{0}}\left(\mathbb{R}^{n}\right) \tag{8.7}
\end{equation*}
$$

and will aim to find a relation of the form $B\left(q_{0}, R, n\right) \leqslant c\left(q_{0}, \delta\right) R^{\delta}+$ $\varepsilon B\left(q_{0}, R, n\right)$ for some $\varepsilon<1$ and for $R$ larger than some $R_{1}$, for example with $\varepsilon=1 / 2$, thus reaching the conclusion that $B\left(q_{0}, R, n\right) \leqslant 2 c\left(q_{0}, \delta\right) R^{\delta}$ when $R>R_{1}$. We know that $B\left(q_{0}, R, n\right)$ is finite for every dimension $n$, for instance as a consequence of the trivial bound $\left\|\nabla \mu^{R}\right\|_{1} \leqslant\left\|\nabla G^{R}\right\|_{1} \leqslant \kappa \sqrt{n} R$.

We must notice that the a priori estimate in $\mathbb{R}^{n}$ yields the same estimate for the dimensions $\ell \leqslant n$, with a smaller or equal constant, precisely, we must know that $B\left(q_{0}, R, \ell\right) \leqslant B\left(q_{0}, R, n\right)$ when $1 \leqslant \ell \leqslant n$. Indeed, the forthcoming proof in dimension $n$ will bring the question down to dimensions $\ell \leqslant n$, where we shall use the a priori bound by $B\left(q_{0}, R, \ell\right)$. For justifying the validity of the same bound when $\ell \leqslant n$, apply the case $n$ to a function $g$ of the form $g_{1} \otimes \varphi$, namely

$$
g\left(x_{1}, x_{2}\right)=g_{1}\left(x_{1}\right) \varphi\left(x_{2}\right),
$$

where $x_{1}$ is in $\mathbb{R}^{\ell}, g_{1} \in L^{q_{0}}\left(\mathbb{R}^{\ell}\right), x_{2} \in \mathbb{R}^{n-\ell}$ and where $\varphi$ is a fixed $C^{\infty}$ function with compact support in $\mathbb{R}^{n-\ell}$, not identically zero. The indicator of the cube and the Gaussian density have a product structure, which allows us to write

$$
K^{R}\left(x_{1}, x_{2}\right)=K_{1}^{R}\left(x_{1}\right) \psi\left(x_{2}\right), \quad \mathrm{d} \mu^{R}\left(x_{1}, x_{2}\right)=\mathrm{d} \mu_{1}^{R}\left(x_{1}\right) \otimes\left(\psi\left(x_{2}\right) \mathrm{d} x_{2}\right)
$$

where $K_{1}, K_{1}^{R}$ and $\mathrm{d} \mu_{1}^{R}\left(x_{1}\right)=K_{1}^{R}\left(x_{1}\right) \mathrm{d} x_{1}$ correspond to the cube in $\mathbb{R}^{\ell}$, and $\psi$ is a probability density on $\mathbb{R}^{n-\ell}$ corresponding to the cube in $\mathbb{R}^{n-\ell}$. We also have

$$
\mu^{R} * g=\left(\mu_{1}^{R} * g_{1}\right) \otimes(\psi * \varphi) .
$$

The gradient of $\mu^{R} * g$ contains $\left(\nabla \mu_{1}^{R} * g_{1}\right) \otimes(\psi * \varphi)$ in its first $\ell$ coordinates, thus

$$
\begin{aligned}
&\left\|\nabla \mu_{1}^{R} * g_{1}\right\|_{L^{q_{0}}\left(\mathbb{R}^{\ell}\right)}\|\psi * \varphi\|_{L^{q_{0}\left(\mathbb{R}^{n-\ell}\right)}}=\left\|\left(\nabla \mu_{1}^{R} * g_{1}\right) \otimes(\psi * \varphi)\right\|_{L^{q_{0}}\left(\mathbb{R}^{n}\right)} \\
& \leqslant\left\|\nabla \mu^{R} * g\right\|_{L^{q_{0}}\left(\mathbb{R}^{n}\right)} \leqslant B\left(q_{0}, R, n\right)\|g\|_{L^{q_{0}}\left(\mathbb{R}^{n}\right)} \\
&=B\left(q_{0}, R, n\right)\left\|g_{1}\right\|_{L^{q_{0}}\left(\mathbb{R}^{\ell}\right)}\|\varphi\|_{L^{q_{0}}\left(\mathbb{R}^{n-\ell}\right)} .
\end{aligned}
$$

This yields $B\left(q_{0}, R, \ell\right) \leqslant B\left(q_{0}, R, n\right)\|\varphi\|_{L^{q_{0}}\left(\mathbb{R}^{n-\ell}\right)} /\|\psi * \varphi\|_{L^{q_{0}\left(\mathbb{R}^{n-\ell}\right)}}$ and by spreading $\varphi$, replacing it with $\varphi_{k}: x \mapsto \varphi(x / k), k \rightarrow+\infty$, one makes the quotient of norms tend to 1 , thus proving that $B\left(q_{0}, R, \ell\right) \leqslant B\left(q_{0}, R, n\right)$.

### 8.2. First reduction

One applies a result of Pisier [62] about holomorphic semi-groups. If $\mathbf{T}=$ $\left(T_{j}\right)_{j=1}^{n}$ is a family of bounded linear operators on $L^{q}(X, \Sigma, \mu), 1 \leqslant q \leqslant+\infty$, we introduce for every subset $J \subset N=\{1, \ldots, n\}$ the operators

$$
\mathbf{T}^{J}=\prod_{j \in J} T_{j}, \quad \mathbf{T}^{\sim J}=\mathbf{T}^{N \backslash J}=\prod_{j \notin J} T_{j}
$$

and $\mathbf{T}^{\sim j}$ will be a short form for $\mathbf{T}^{\sim\{j\}}, 1 \leqslant j \leqslant n$. We found the notation $\mathbf{T}^{\sim J}$ convenient, but it might be ambiguous, since it depends on the ambient set $N$.

Given commuting projectors $\left(E_{j}\right)_{j=1}^{n}$, one can consider the semi-group

$$
T_{t}=\prod_{j=1}^{n}\left(E_{j}+\mathrm{e}^{-t}\left(I-E_{j}\right)\right), \quad t \geqslant 0
$$

where $I$ denotes the identity operator. If we set $z=\mathrm{e}^{-t}$ and expand the product, we can arrange it according to powers of $z$, displaying in this way homogeneous parts $z^{k} H_{k}$ of degree $k$. We see that

$$
T_{t}=\sum_{k=0}^{n} z^{k}\left(\sum_{|J|=k} \mathbf{E}^{\sim J}(\mathbf{I}-\mathbf{E})^{J}\right)=\sum_{k=0}^{n} z^{k} H_{k}=\sum_{k=0}^{n} \mathrm{e}^{-k t} H_{k}
$$

Letting $\Sigma_{k}$ denote the family of subsets $J \subset N$ of cardinality $k$, we have

$$
\begin{equation*}
H_{k}=\sum_{J \in \Sigma_{k}} \mathbf{E}^{\sim J}(\mathbf{I}-\mathbf{E})^{J}, \quad k=0, \ldots, n, \quad \text { and } \quad \sum_{k=0}^{n} H_{k}=T_{0}=I \tag{8.8}
\end{equation*}
$$

Proposition 8.3 (after Pisier [62]). - Let $\left(E_{j}\right)_{j=1}^{n}$ be a family of commuting conditional expectation projectors on $L^{q}(X, \Sigma, \mu), 1<q<+\infty$, and consider the semi-group

$$
P_{t}=\prod_{j=1}^{n}\left(\mathrm{e}^{-t} I+\left(1-\mathrm{e}^{-t}\right) E_{j}\right)=\prod_{j=1}^{n}\left(E_{j}+\mathrm{e}^{-t}\left(I-E_{j}\right)\right), \quad t \geqslant 0
$$

This semi-group is analytic on $L^{q}(X, \Sigma, \mu), 1<q<+\infty$, with an extension $\left(P_{z}\right)_{z \in \Omega_{\varphi_{q}}}$ to a sector $\Omega_{\varphi_{q}}=\left\{z=r \mathrm{e}^{\mathrm{i} \theta}: r>0,|\theta|<\varphi_{q}\right\}$ in $\mathbb{C}$, where $\varphi_{q}>0$ depends on $q$ only. The extension is bounded uniformly in $q$ on every compact subset of $\Omega_{\varphi_{q}}$. There exists $h_{q} \geqslant 1$ independent of $n$ such that whenever $0 \leqslant k \leqslant n$, the homogeneous part $H_{k}$ in (8.8) is bounded on $L^{q}(X, \Sigma, \mu)$ by $\left(h_{q}\right)^{k}$.

That $h_{q} \geqslant 1$ can be seen on any example $P_{t} f=E_{1} f+\mathrm{e}^{-t}\left(f-E_{1} f\right)$ with $n=1$ and $E_{1} \neq I$. Then $H_{1}$ is the projector $I-E_{1} \neq 0$, hence
$h_{q} \geqslant\left\|H_{1}\right\|_{q \rightarrow q} \geqslant 1$. If $\left(E_{j, s}\right)_{j=1}^{n}, s \in[0,1]$, is a family of such conditional expectations, where $E_{j, s}$ and $E_{k, t}$ commute for all $j \neq k$ and all $s, t \in[0,1]$, and if we set for example

$$
U_{j}=\int_{0}^{1} E_{j, s} \mathrm{~d} s, \quad j=1, \ldots, n
$$

then we see that

$$
Q_{t}=\prod_{j=1}^{n}\left(\mathrm{e}^{-t} I+\left(1-\mathrm{e}^{-t}\right) U_{j}\right)=\int_{[0,1]^{n}} P_{t, s_{1}, \ldots, s_{n}} \mathrm{~d} s_{1} \ldots \mathrm{~d} s_{n}
$$

where each $P_{t, s_{1}, \ldots, s_{n}}=\prod_{j=1}^{n}\left(\mathrm{e}^{-t} I+\left(1-\mathrm{e}^{-t}\right) E_{j, s_{j}}\right)$ is of "Pisier type". Also, the corresponding homogeneous parts are of the form

$$
\begin{aligned}
\widetilde{H}_{k} & =\sum_{J \in \Sigma_{k}} \mathbf{U}^{\sim J}(\mathbf{I}-\mathbf{U})^{J} \\
& =\int_{[0,1]^{n}} \sum_{J \in \Sigma_{k}}\left(\prod_{i \notin J} E_{i, s_{i}}\right)\left(\prod_{j \in J}\left(I-E_{j, s_{j}}\right)\right) \mathrm{d} s_{1} \mathrm{~d} s_{2} \ldots \mathrm{~d} s_{n}
\end{aligned}
$$

that are averages of terms $H_{k}\left(s_{1}, \ldots, s_{n}\right)$ bounded by $h_{q}^{k}$ according to Proposition 8.3. The result of Proposition 8.3 generalizes thus to families such as $\left(U_{j}\right)_{j=1}^{n}$.

We shall apply Proposition 8.3 to operators $\left(E_{j}\right)_{j=1}^{n}$ of conditional expectation on $L^{q}\left(\mathbb{R}^{n}\right)$, where each $E_{j}$ is acting in the $x_{j}$ variable and $1 \leqslant j \leqslant n$. For one variable and $s_{0} \in \mathbb{R}$ fixed, we associate to a locally integrable function $f$ on $\mathbb{R}$ its averages on length one intervals $I_{r}=\left[s_{0}+r, s_{0}+r+1\right)$, $r \in \mathbb{Z}$, defining $E_{s_{0}}$ by

$$
\left(E_{s_{0}} f\right)(v)=\sum_{r \in \mathbb{Z}}\left(\int_{I_{r}} f(s) \mathrm{d} s\right) \mathbf{1}_{I_{r}}(v), \quad v \in \mathbb{R} .
$$

This operator is a conditional expectation, as considered in Remark 1.2. We define operators $E_{j, s_{0}}, j=1, \ldots, n$, on $L_{\mathrm{loc}}^{1}\left(\mathbb{R}^{n}\right)$ by the analogous formula, acting on the $x_{j}$ variable. When $j=1$ for example, we let

$$
\left(E_{1, s_{0}} f\right)\left(x_{1}, x_{2}, \ldots, x_{n}\right)=\sum_{r \in \mathbb{Z}}\left(\int_{I_{r}} f\left(s, x_{2}, \ldots, x_{n}\right) \mathrm{d} s\right) \mathbf{1}_{I_{r}}\left(x_{1}\right)
$$

Averaging on values of $s_{0}$, one can replace the $E_{j}$ s by convolution operators with probability densities $\chi$ on $\mathbb{R}$ of the form

$$
\begin{equation*}
\chi(x)=\int_{\mathbb{R}} \mathbf{1}_{[s, s+1]}(x) \mathrm{d} \nu(s), \quad x \in \mathbb{R} \tag{8.9}
\end{equation*}
$$

where $\nu$ is a probability measure on the line. We see that $\chi(x)=F(x)-$ $F(x-1)$, with $F(x)=\nu[(-\infty, x)]$ non-decreasing, $F(-\infty)=0$ and
$F(+\infty)=1$. One can also proceed to changes of scale. Summarizing, we have the lemma that follows.

Lemma 8.4 (Bourgain [13], Lemma 5). - Let $\chi$ be a compactly supported probability density on $\mathbb{R}$ of the form (8.9). Denote by $T_{j}$ the convolution operator with $\chi_{\left(t_{j}\right)}$ in the $x_{j}$ variable, $t_{j}>0, j=1, \ldots, n$. For $0 \leqslant k \leqslant n$, the norm of the operator

$$
H_{k}:=\sum_{S \in \Sigma_{k}} \mathbf{T}^{\sim S}(\mathbf{I}-\mathbf{T})^{S}
$$

on $L^{q}\left(\mathbb{R}^{n}\right)$ is bounded by $h_{q}^{k}$, with $1<q<+\infty$ and $h_{q}$ from Proposition 8.3.
In what follows, we denote by $T_{j}, j=1, \ldots, n$, the convolution in the $x_{j}$ variable on $L^{q}\left(\mathbb{R}^{n}\right)$ by $\eta_{\left(w_{0}\right)}\left(x_{j}\right)$, where $w_{0}=R^{-\delta / 2}$ will stay fixed and where
$\eta(x)=(1-|x|)_{+}=\int_{-1 / 2}^{1 / 2} \mathbf{1}_{[-1 / 2+s, 1 / 2+s]}(x) \mathrm{d} s=\left(\mathbf{1}_{[-1 / 2,1 / 2]} * \mathbf{1}_{[-1 / 2,1 / 2]}\right)(x)$.
Since $\eta$ is a convolution square, $\widehat{\eta}$ is real and nonnegative. We have

$$
\widehat{\eta}(t)=\left(\frac{\sin (\pi t)}{\pi t}\right)^{2}, \quad \text { and } \quad \widehat{\eta}^{\prime \prime}(t)=-4 \pi^{2} \int_{\mathbb{R}} s^{2}(1-|s|)_{+} \cos (2 \pi s t) \mathrm{d} s
$$

for every $t \in \mathbb{R}$, thus

$$
\left|\widehat{\eta}^{\prime \prime}(t)\right| \leqslant 8 \pi^{2} \int_{0}^{1} s^{2}(1-s) \mathrm{d} s=\frac{8 \pi^{2}}{12}<8
$$

By the Taylor formula we get

$$
\begin{equation*}
0 \leqslant 1-\widehat{\eta}(t) \leqslant\left(4 t^{2}\right) \wedge 1 \tag{8.10}
\end{equation*}
$$

For every subset $S \subset N:=\{1, \ldots, n\}$ let us set

$$
\begin{equation*}
\Gamma^{S}=\mathbf{T}^{\sim S}(\mathbf{I}-\mathbf{T})^{S} \tag{8.11}
\end{equation*}
$$

The homogeneous parts $\left(H_{k}\right)$ in $Q_{t}=\prod_{j=1}^{n}\left(\mathrm{e}^{-t} I+\left(1-\mathrm{e}^{-t}\right) T_{j}\right)$ have the form

$$
H_{k}=\sum_{S \in \Sigma_{k}} \Gamma^{S}, \quad 0 \leqslant k \leqslant n, \quad \text { and } \quad \sum_{k=0}^{n} H_{k}=I
$$

In particular, $H_{0}=\Gamma^{\emptyset}=\mathbf{T}^{N}=\prod_{j=1}^{n} T_{j}$ has norm $\leqslant 1$ on every space $L^{q}\left(\mathbb{R}^{n}\right)$, for $1 \leqslant q \leqslant+\infty$, since $H_{0}$ is the convolution with the product probability density $\prod_{j=1}^{n} \eta_{\left(w_{0}\right)}\left(x_{j}\right)$. When $1<q<+\infty$ and $1 \leqslant k \leqslant n$, we have $\left\|H_{k}\right\|_{q \rightarrow q} \leqslant h_{q}^{k}$ by Proposition 8.3. It is convenient to set $H_{k}=0$ below when $k>n$.

To every given function $g$ in $L^{q}\left(\mathbb{R}^{n}\right)$, we shall apply a decomposition of the form $g=H_{0} g+\cdots+H_{M-1} g+h$, and consider the corresponding expression

$$
\begin{equation*}
\nabla \mu^{R} * g=\nabla \mu^{R} * H_{0} g+\cdots+\nabla \mu^{R} * H_{M-1} g+\nabla \mu^{R} * h \tag{8.12}
\end{equation*}
$$

where $M \geqslant 1$ will be chosen as a function of the already fixed $p_{0}$ and $\delta>0$. We have to estimate in $L^{q_{0}}\left(\mathbb{R}^{n}\right)$ the successive terms in (8.12). The function $h$ is considered as a small rest, the mapping $g \mapsto \nabla \mu^{R} * h$ will be handled in $L^{2}\left(\mathbb{R}^{n}\right)$ by a Fourier estimate, and in some $L^{q_{1}}\left(\mathbb{R}^{n}\right), q_{1}>q_{0}$, as a consequence of Proposition 8.3. We choose $M$ large enough for deducing from $\left\|\nabla \mu^{R} * h\right\|_{2} \leqslant \kappa R^{1-\delta M / 2}\|g\|_{2}$ and $\left\|\nabla \mu^{R} * h\right\|_{q_{1}} \leqslant \kappa R\|g\|_{q_{1}}$ that one has by interpolation

$$
\begin{equation*}
\left\|\nabla \mu^{R} * h\right\|_{q_{0}} \leqslant \kappa\left(q_{0}, \delta\right)\|g\|_{q_{0}} \tag{8.13}
\end{equation*}
$$

which is just perfect in the direction of (8.7). Recall that $\mu_{j}^{R}$ denotes the $j$ th partial derivative $\partial_{j} \mu^{R}=\left(\partial_{j} \mu_{Q}\right) * G^{R}$ of $\mu^{R}$, so that $\left|\nabla \mu^{R} * h\right|^{2}=$ $\sum_{j=1}^{n}\left|\mu_{j}^{R} * h\right|^{2}$.

We factor the mapping $g \mapsto \nabla \mu^{R} * h$ into $U_{K^{R}}: h \mapsto \nabla \mu^{R} * h$ and $A: g \mapsto h$, i.e., $A=I-H_{0}-\cdots-H_{M-1}=\sum_{k \geqslant M}^{K^{R}} H_{k}$. We look for estimates in $L^{2}$ and $L^{q}, q_{0}<q<+\infty$. For $U_{K^{R}}$ we use Lemma 7.11 and get by (8.2) and (8.6) that

$$
\left\|U_{K^{R}}\right\|_{q \rightarrow q} \leqslant 2^{1 / q} \sigma_{R}^{-2 / q} V\left(K^{R}\right)^{1-2 / q} \leqslant(24)^{1 / q} R^{1-2 / q}<5 R
$$

since $q \geqslant 2$. On the other hand, by Lemma 8.4, the mapping $A: g \mapsto h$ is bounded in $L^{q}\left(\mathbb{R}^{n}\right)$ by $1+\sum_{k=0}^{M-1} h_{q}^{k} \leqslant(M+1) h_{q}^{M-1}$. It follows that

$$
\begin{equation*}
\left\|\nabla \mu^{R} * h\right\|_{q}=\left\|U_{K^{R}} h\right\|_{q} \leqslant 5 R\|h\|_{q} \leqslant 5 R(M+1) h_{q}^{M-1}\|g\|_{q} \tag{8.14}
\end{equation*}
$$

This is also valid when $q=2$, but the point is that we will then get a much better bound by factoring now $g \mapsto \nabla \mu^{R} * h$ as $U_{G^{R}} \circ B$, with $U_{G^{R}}: f \mapsto \nabla G^{R} * f$ and $B: g \mapsto \mu_{Q} * A g=\mu_{Q} * h$. We begin by estimating

$$
\left\|\mu_{Q} * h\right\|_{2}=\left\|\mu_{Q} *\left(\sum_{k \geqslant M} H_{k}\right) g\right\|_{2}=\left\|\mu_{Q} *\left(\sum_{|S| \geqslant M} \Gamma^{S}\right) g\right\|_{2}
$$

One needs to control the $L^{\infty}\left(\mathbb{R}^{n}\right)$ norm of the function $\xi \mapsto L(\xi)$, where $L$ is the multiplier associated to the mapping $B$. It is the aim of the next lemma. One sees that

$$
L(\xi):=\left(\prod_{j=1}^{n} \frac{\sin \left(\pi \xi_{j}\right)}{\pi \xi_{j}}\right)\left(\sum_{|S| \geqslant M} \prod_{j \notin S} \widehat{\eta}\left(w_{0} \xi_{j}\right) \prod_{j \in S}\left(1-\widehat{\eta}\left(w_{0} \xi_{j}\right)\right)\right)
$$

Lemma 8.5 (see [13, Equations (2.9), (2.11)]). - For $0 \leqslant u \leqslant 1 / 4$ and every $\xi \in \mathbb{R}^{n}$, one has that

$$
\left|\left(\prod_{j=1}^{n} \frac{\sin \left(\pi \xi_{j}\right)}{\pi \xi_{j}}\right)\left(\sum_{|S| \geqslant M} \prod_{j \notin S} \widehat{\eta}\left(u \xi_{j}\right) \prod_{j \in S}\left(1-\widehat{\eta}\left(u \xi_{j}\right)\right)\right)\right| \leqslant u^{M}
$$

Proof. - We know from (8.10) that $0 \leqslant \widehat{\eta}(t) \leqslant 1$ and $1-\widehat{\eta}(t) \leqslant\left(4 t^{2}\right) \wedge 1$. We introduce $v=1 / u \geqslant 4$ and begin by checking that for every $t \geqslant 0$, we have

$$
X(t):=\left|\frac{\sin (\pi t)}{\pi t}\right|\left(1+v\left[\left(4 u^{2} t^{2}\right) \wedge 1\right]\right) \leqslant 1
$$

Consider first the case $0 \leqslant t \leqslant 1 /(2 u)$. One has then $4 u^{2} t^{2} \leqslant 1$ and it follows that $1+v\left[\left(4 u^{2} t^{2}\right) \wedge 1\right]=1+4 u t^{2}$. If in addition $0 \leqslant t \leqslant 1$, then, for example by the Euler product formula (3.2.E), we have $|\sin (\pi t) / \pi t| \leqslant 1-t^{2}$, and since $4 u \leqslant 1$ by assumption, we get

$$
X(t) \leqslant\left(1-t^{2}\right)\left(1+4 u t^{2}\right) \leqslant\left(1-t^{2}\right)\left(1+t^{2}\right) \leqslant 1
$$

When $1<t \leqslant 1 /(2 u)$, we have

$$
\left|\frac{\sin (\pi t)}{\pi t}\right|\left(1+4 u t^{2}\right) \leqslant \frac{1+4 u t^{2}}{\pi t}=\frac{1}{\pi}(1 / t+4 u t) \leqslant \frac{3}{\pi}<1
$$

In the second case, when $2 u t>1$, we can write

$$
X(t) \leqslant \frac{1+v}{\pi t} \leqslant \frac{2 u(1+v)}{\pi} \leqslant \frac{1 / 2+2}{\pi}<1
$$

Expanding the product $\prod_{j=1}^{n} X\left(\xi_{j}\right)$ and since $X$ is even, one sees that

$$
\begin{aligned}
1 & \geqslant \prod_{j=1}^{n} X\left(\xi_{j}\right)=\prod_{j=1}^{n}\left|\frac{\sin \left(\pi \xi_{j}\right)}{\pi \xi_{j}}\right|\left(1+v\left[\left(4 u^{2} \xi_{j}^{2}\right) \wedge 1\right]\right) \\
& \geqslant \prod_{j=1}^{n}\left|\frac{\sin \left(\pi \xi_{j}\right)}{\pi \xi_{j}}\right|\left(\widehat{\eta}\left(u \xi_{j}\right)+v\left(1-\widehat{\eta}\left(u \xi_{j}\right)\right)\right) \\
& \geqslant v^{M}\left|\left(\prod_{j=1}^{n} \frac{\sin \left(\pi \xi_{j}\right)}{\pi \xi_{j}}\right)\left(\sum_{|S| \geqslant M} \prod_{j \notin S} \widehat{\eta}\left(u \xi_{j}\right) \prod_{j \in S}\left(1-\widehat{\eta}\left(u \xi_{j}\right)\right)\right)\right| .
\end{aligned}
$$

By Lemma 7.11, we have that $\left\|U_{G^{R}}\right\|_{2 \rightarrow 2} \leqslant \sqrt{\pi} R<2 R$, because the variance of $G^{R}$ is $2 \pi^{-1} R^{-2}$. Let us define $R_{0}$ by $R_{0}^{\delta / 2}=4$. If $R \geqslant R_{0}$, then $w_{0}=R^{-\delta / 2} \leqslant 1 / 4$, we obtain from Lemma 8.5 with $u=w_{0}$ the final control

$$
\left\|\nabla \mu^{R} * h\right\|_{2}=\left\|U_{G^{R}}\left(\mu_{Q} * h\right)\right\|_{2} \leqslant 2 R\left\|\mu_{Q} * h\right\|_{2} \leqslant 2 R R^{-\delta M / 2}\|g\|_{2}
$$

We use now (8.14) with for example $q=q_{1}=2 q_{0} / p_{0}>q_{0}$. Letting $\theta=1 / p_{0}$, we have $(1-\theta) / 2+\theta / q_{1}=1 / q_{0}$ and we see by interpolation for $g \mapsto \nabla \mu^{R} * h$
that

$$
\left\|\nabla \mu^{R} * h\right\|_{q_{0}} \leqslant\left(2 R^{-\delta M / 2}\right)^{1 / q_{0}} R\left(5(M+1) h_{q_{1}}^{M-1}\right)^{1 / p_{0}}\|g\|_{q_{0}}
$$

We select $M=M(\delta)=\left\lceil 2 q_{0} / \delta\right\rceil$, so that $\delta M /\left(2 q_{0}\right) \geqslant 1$. When $R \geqslant R_{0}$ we get

$$
\begin{equation*}
\left\|\nabla \mu^{R} * h\right\|_{q_{0}} \leqslant \kappa_{q_{0}, \delta}\|g\|_{q_{0}} \text { with } \kappa_{q_{0}, \delta} \leqslant 5\left(2+2 q_{0} / \delta\right)^{1 / p_{0}} h_{2 q_{0} / p_{0}}^{2 q_{0} /\left(\delta p_{0}\right)} \tag{8.15}
\end{equation*}
$$

In what follows we assume that $R \geqslant R_{0}$, hence $R^{\delta} \geqslant 16$. In the conclusion section, we shall need the following bound for a Fourier transform.

Lemma 8.6. - For every $r \in \mathbb{R}, \ell \geqslant 1$ and all $\xi=\left(\xi_{1}, \ldots, \xi_{\ell}\right) \in \mathbb{R}^{\ell}$, one has that

$$
\left(1-\mathrm{e}^{-r^{2}|\xi|^{2}}\right) \prod_{j=1}^{\ell} \widehat{\eta}\left(\xi_{j}\right) \leqslant r^{2}
$$

Proof. - We observe first that

$$
\widehat{\eta}(t)=\left(\frac{\sin (\pi t)}{\pi t}\right)^{2} \leqslant \frac{1}{1+t^{2}}
$$

This is clear when $|t| \geqslant 1$ because $\widehat{\eta}(t) \leqslant(\pi t)^{-2}$ and $1+t^{2}<\pi^{2} t^{2}$ in this case. When $|t| \leqslant 1$ we have $\widehat{\eta}(t) \leqslant|\sin (\pi t)| /|\pi t| \leqslant\left(1-t^{2}\right) \leqslant\left(1+t^{2}\right)^{-1}$ by (3.2.E). It suffices thus to bound for $x \in \mathbb{R}^{\ell}$ the expression

$$
F(x)=\left(1-\mathrm{e}^{-r^{2}|x|^{2}}\right) \prod_{j=1}^{\ell} \frac{1}{1+x_{j}^{2}} \geqslant 0
$$

The function $F$ tends to 0 at infinity, we have at any maximum $\bar{x} \neq 0$ that

$$
\frac{2 r^{2} \bar{x}_{j} \mathrm{e}^{-r^{2}|\bar{x}|^{2}}}{1-\mathrm{e}^{-r^{2}|\bar{x}|^{2}}}=\frac{2 \bar{x}_{j}}{1+\bar{x}_{j}^{2}}, \quad j=1, \ldots, \ell
$$

The nonzero coordinates of $\bar{x}$ have the same square $\bar{x}_{j}^{2}=: y>0$, and if $k$ denotes their cardinality, we have $0<k \leqslant \ell$ and $|\bar{x}|^{2}=k y$. It follows that

$$
k r^{2} y \leqslant \mathrm{e}^{k r^{2} y}-1=r^{2}(1+y) \leqslant r^{2}(1+y)^{k}
$$

Finally, we have $F(\bar{x})=\left(1-\mathrm{e}^{-k r^{2} y}\right)(1+y)^{-k} \leqslant k r^{2} y(1+y)^{-k} \leqslant r^{2}$.

### 8.2.1. Decoupling

We have to analyze each of the expressions $\nabla \mu^{R} * H_{k} g$ in (8.12), for $0 \leqslant k<M$. When $1 \leqslant k<M$, we handle this by a decoupling argument that will allow us to essentially reduce to the cases where $k=0,1$, but
in a dimension $\ell \leqslant n$. Before proceeding by a Bourgainian technique of "selectors", we split

$$
\left|\nabla \mu^{R} * H_{k} g\right|=\left(\sum_{j=1}^{n}\left|\mu_{j}^{R} * H_{k} g\right|^{2}\right)^{1 / 2}=\left(\sum_{j=1}^{n}\left|\mu_{j}^{R} *\left(\sum_{S \in \Sigma_{k}} \Gamma^{S} g\right)\right|^{2}\right)^{1 / 2}
$$

into two. For each $j$ in $\{1, \ldots, n\}$, let $\Sigma_{k}^{j}$ and $\Sigma_{k}^{\sim j}$ denote respectively the family of subsets $S$ of $\{1, \ldots, n\}$ with cardinality $|S|=k$ containing $j$, resp. such that $j \notin S$. Then $\left|\nabla \mu^{R} * H_{k} g\right|$ is bounded by the sum of the two expressions

$$
\begin{equation*}
\mathbf{E}_{k}(R, n, g):=\left(\sum_{j=1}^{n}\left|\mu_{j}^{R} *\left(\sum_{S \in \Sigma_{k}^{\sim j}} \Gamma^{S} g\right)\right|^{2}\right)^{1 / 2} \tag{8.16a}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathbf{F}_{k}(R, n, g):=\left(\sum_{j=1}^{n}\left|\mu_{j}^{R} *\left(\sum_{S \in \Sigma_{k}^{j}} \Gamma^{S} g\right)\right|^{2}\right)^{1 / 2} \tag{8.16b}
\end{equation*}
$$

Assume that $1 \leqslant k<M=M(\delta)$. Let $\left(\gamma_{i}\right)_{1 \leqslant i \leqslant n}$ be independent $\{0,1\}$ valued random variables with mean $1 /(k+1)$ on some probability space $(\Omega, \mathcal{F}, P)$. For each $j$ in $\{1, \ldots, n\}$ and $S \in \Sigma_{k}^{\sim j}$, let $\sigma_{S, j}=\gamma_{j} \prod_{i \in S}\left(1-\gamma_{i}\right)$. We have that

$$
\mathrm{E} \sigma_{S, j}=\frac{1}{k+1}\left(1-\frac{1}{k+1}\right)^{k}=\frac{k^{k}}{(k+1)^{k+1}}=: e_{k}, \quad j=1, \ldots, n
$$

and $e_{k}^{-1} \leqslant \mathrm{e}(k+1) \leqslant \mathrm{e} M$ because $\mathrm{e}^{1 / k}>1+1 / k$. By convexity, we see that

$$
\begin{aligned}
e_{k} \mathbf{E}_{k}(R, n, g) & =\left(\sum_{j=1}^{n}\left|\mu_{j}^{R} *\left(\sum_{S \in \Sigma_{k}^{\sim j}}\left[\mathrm{E}_{\omega} \sigma_{S, j}(\omega)\right] \Gamma^{S} g\right)\right|^{2}\right)^{1 / 2} \\
& \leqslant \mathrm{E}_{\omega}\left[\left(\sum_{j=1}^{n}\left|\mu_{j}^{R} *\left(\sum_{S \in \Sigma_{k}^{\sim j}} \sigma_{S, j}(\omega) \Gamma^{S} g\right)\right|^{2}\right)^{1 / 2}\right]
\end{aligned}
$$

Let $q \geqslant 1$ be given. It follows that for some $\omega_{0} \in \Omega$, we have

$$
\begin{align*}
& \left\|\mathbf{E}_{k}(R, n, g)\right\|_{L^{q}\left(\mathbb{R}^{n}\right)} \\
& \quad \leqslant \mathrm{e} M\left\|\left(\sum_{j=1}^{n}\left|\mu_{j}^{R} *\left(\sum_{S \in \Sigma_{k}^{\sim j}} \sigma_{S, j}\left(\omega_{0}\right) \Gamma^{S} g\right)\right|^{2}\right)^{1 / 2}\right\|_{L^{q}\left(\mathbb{R}^{n}\right)} \tag{8.17}
\end{align*}
$$

Let $J_{0}=\left\{j: \gamma_{j}\left(\omega_{0}\right)=1\right\}$. Then $\sigma_{S, j}\left(\omega_{0}\right)=0$ whenever $S$ meets $J_{0}$ or $j \notin J_{0}$. The $L^{q}\left(\mathbb{R}^{n}\right)$ norm at the right-hand side of (8.17) is therefore the
norm of

$$
E\left(J_{0}, g\right):=\left(\sum_{j \in J_{0}}\left|\mu_{j}^{R} *\left(\sum_{S \in \Sigma_{k}^{\sim J_{0}}} \Gamma^{S} g\right)\right|^{2}\right)^{1 / 2}
$$

where $\Sigma_{k}^{\sim J_{0}}$ denotes the family of subsets $S$ of $\{1, \ldots, n\}$ such that $|S|=k$ and that are disjoint from $J_{0}$. Let us introduce the operator

$$
\mathbf{U}=\sum_{S \in \Sigma_{k}^{\sim J_{0}}} \mathbf{T}^{\sim\left(J_{0} \cup S\right)}(\mathbf{I}-\mathbf{T})^{S} \quad \text { and the function } \quad \Psi=\mathbf{U} g
$$

on $\mathbb{R}^{n}$. We see that $\mathbf{T}^{J_{0}} \mathbf{U}=\sum_{S \in \Sigma_{k}^{\sim J_{0}}} \Gamma^{S}$, and the operator $\mathbf{U}$ acts on the variables not in $J_{0}$ as does the $k$ th homogeneous part $H_{k}$ relative to $\mathbb{R}^{\{1, \ldots, n\} \backslash J_{0}}$. Consequently, applying Proposition 8.3 in the variables $\mathbf{x}^{\sim J_{0}}=$ $\left(x_{i}\right)_{i \notin J_{0}}$, we get

$$
\begin{equation*}
\left\|\Psi_{\mathbf{x}^{J_{0}}}\right\|_{L^{q}\left(\mathbb{R}^{\sim J_{0}}\right)} \leqslant h_{q}^{k}\left\|g_{\mathbf{x}^{J_{0}}}\right\|_{L^{q}\left(\mathbb{R}^{\sim J_{0}}\right)} \tag{8.18}
\end{equation*}
$$

for every fixed $\mathbf{x}^{J_{0}}=\left(x_{i}\right)_{i \in J_{0}}$, where $f_{\mathbf{x}^{J}}\left(\mathbf{x}^{\sim J}\right):=f\left(\mathbf{x}^{J}, \mathbf{x}^{\sim J}\right)=f(x)$, and we see that $E\left(J_{0}, g\right)=\left(\sum_{j \in J_{0}}\left|\mu_{j}^{R} * \mathbf{T}^{J_{0}} \Psi\right|^{2}\right)^{1 / 2}$. Assume that there exists $b_{0}\left(q_{0}, R, n\right)$ such that for every subset $J$ of $\{1, \ldots, n\}$ and $f \in L^{q_{0}}\left(\mathbb{R}^{J}\right)$ we have

$$
\begin{equation*}
\left\|\left(\sum_{j \in J}\left|\left(\mu_{Q^{J}}\right)_{j}^{R} *\left(\prod_{i \in J} T_{i}\right) f\right|^{2}\right)^{1 / 2}\right\|_{L^{q_{0}}\left(\mathbb{R}^{J}\right)} \leqslant b_{0}\left(q_{0}, R, n\right)\|f\|_{L^{q_{0}}\left(\mathbb{R}^{J}\right)} \tag{8.19}
\end{equation*}
$$

with $\mu_{Q^{J}}$ uniform on $Q^{J}:=[-1 / 2,1 / 2]^{J}$ in $\mathbb{R}^{J}$. It follows from (8.18), by integrating in the $J_{0}$ variables, that

$$
\left\|E\left(J_{0}, g\right)\right\|_{L^{q_{0}}\left(\mathbb{R}^{n}\right)} \leqslant b_{0}\left(q_{0}, R, n\right) h_{q}^{k}\|g\|_{L^{q_{0}}\left(\mathbb{R}^{n}\right)}
$$

For $\mathbf{F}_{k}(R, n, g)$ we proceed similarly, writing each $S \in \Sigma_{k}^{j}$ as $S=\{j\} \cup S_{1}$, with $\left|S_{1}\right|=k-1$, and using now $\sigma_{S_{1}, j}=\gamma_{j} \prod_{i \in S_{1}}\left(1-\gamma_{i}\right)$ for which we have $\mathrm{E} \sigma_{S_{1}, j}=k^{k-1}(k+1)^{-k} \geqslant 1 /(\mathrm{e} k)>1 /(\mathrm{e} M)$. We obtain for some $\omega_{0} \in \Omega$ that

$$
\left\|\mathbf{F}_{k}(R, n, g)\right\|_{L^{q}\left(\mathbb{R}^{n}\right)} \leqslant \mathrm{e} M\left\|\left(\sum_{j=1}^{n}\left|\mu_{j}^{R} *\left(\sum_{S \in \Sigma_{k}^{j}} \sigma_{S_{1}, j}\left(\omega_{0}\right) \Gamma^{S} g\right)\right|^{2}\right)^{1 / 2}\right\|_{L^{q}\left(\mathbb{R}^{n}\right)}
$$

Considering again $J_{0}=\left\{j: \gamma_{j}\left(\omega_{0}\right)=1\right\}$, we get instead of $E\left(J_{0}, g\right)$ the expression

$$
F\left(J_{0}, g\right)=\left(\sum_{j \in J_{0}}\left|\mu_{j}^{R} *\left(\sum_{S_{1} \in \Sigma_{k-1}^{\sim J_{0}}} \Gamma^{\{j\} \cup S_{1}} g\right)\right|^{2}\right)^{1 / 2}
$$

When $k=1$, we have $S_{1}=\emptyset, S=\{j\}$ and $\sigma_{S_{1}, j}=\gamma_{j}$, the argument remains correct but becomes "inactive". Let now $\Psi=\sum_{S_{1} \in \Sigma_{k-1}^{\sim J_{0}}} \mathbf{T}^{\sim\left(J_{0} \cup S_{1}\right)}(\mathbf{I}-\mathbf{T})^{S_{1}} g$, satisfying by Proposition 8.3 applied to $L^{q}\left(\mathbb{R}^{\sim J_{0}}\right)$ the inequality

$$
\left\|\Psi_{\mathbf{x}^{J_{0}}}\right\|_{L^{q}\left(\mathbb{R}^{\sim J_{0}}\right)} \leqslant h_{q}^{k-1}\left\|g_{\mathbf{x}^{J_{0}}}\right\|_{L^{q}\left(\mathbb{R}^{\sim J_{0}}\right)}
$$

For each $j \in J_{0}$, let $B_{j}=\left(I-T_{j}\right) \mathbf{T}^{J_{0} \backslash\{j\}}$. Then $F\left(J_{0}, g\right)=\left(\sum_{j \in J_{0}} \mid \mu_{j}^{R} *\right.$ $\left.\left.B_{j} \Psi\right|^{2}\right)^{1 / 2}$. If there exists $b_{1}\left(q_{0}, R, n\right)$ such that for every subset $J$ of $\{1, \ldots, n\}$ and every function $f \in L^{q_{0}}\left(\mathbb{R}^{J}\right)$ we have an inequality

$$
\begin{align*}
\|\left(\sum_{j \in J}\left|\mu_{j}^{R} *\left(I-T_{j}\right)\left(\prod_{i \in J, i \neq j} T_{i}\right) f\right|^{2}\right)^{1 / 2} & \|_{L^{q_{0}\left(\mathbb{R}^{J}\right)}} \\
& \leqslant b_{1}\left(q_{0}, R, n\right)\|f\|_{L^{q_{0}}\left(\mathbb{R}^{J}\right)} \tag{8.20}
\end{align*}
$$

it implies that $F\left(J_{0}, g\right)$ may be bounded by $b_{1}\left(q_{0}, R, n\right) h_{q}^{k-1}\|g\|_{L^{q_{0}}\left(\mathbb{R}^{n}\right)}$ in $L^{q_{0}}\left(\mathbb{R}^{n}\right)$.

In view of (8.19) and (8.20), all we need to do in order to control in $L^{q_{0}}\left(\mathbb{R}^{n}\right)$ the expressions $\nabla \mu^{R} * H_{k} g$, when $1 \leqslant k<M$, is to establish in all lower dimensions $\ell \leqslant n$ and for every function $f \in L^{q_{0}}\left(\mathbb{R}^{\ell}\right)$ the inequalities

$$
\begin{align*}
&\left\|\nabla \mu^{R} * H_{0} f\right\|_{L^{q_{0}}\left(\mathbb{R}^{\ell}\right)}=\left\|\left(\sum_{j=1}^{\ell}\left|\mu_{j}^{R} * H_{0} f\right|^{2}\right)^{1 / 2}\right\|_{L^{q_{0}\left(\mathbb{R}^{\ell}\right)}} \\
& \leqslant b_{0}\left(q_{0}, R, n\right)\|f\|_{L^{q_{0}}\left(\mathbb{R}^{\ell}\right)} \tag{8.21}
\end{align*}
$$

and

$$
\begin{align*}
&\|\mathbf{F}(R, \ell, f)\|_{L^{q_{0}}\left(\mathbb{R}^{\ell}\right)}:=\left\|\left(\sum_{j=1}^{\ell}\left|\mu_{j}^{R} * \Gamma^{j} f\right|^{2}\right)^{1 / 2}\right\|_{L^{q_{0}}\left(\mathbb{R}^{\ell}\right)} \\
& \leqslant b_{1}\left(q_{0}, R, n\right)\|f\|_{L^{q_{0}}\left(\mathbb{R}^{\ell}\right)} \tag{8.22}
\end{align*}
$$

for suitable $b_{0}\left(q_{0}, R, n\right)$ and $b_{1}\left(q_{0}, R, n\right)$, with $\Gamma^{j}:=\Gamma^{\{j\}}=\left(I-T_{j}\right) \mathbf{T}^{\{1, \ldots, \ell\} \backslash\{j\}}$. Note that (8.21) controls the so far neglected term $k=0$ in (8.12). From (8.13) and the preceding, this will permit us to estimate

$$
\left\|\nabla \mu^{R} * g\right\|_{L^{q_{0}}\left(\mathbb{R}^{n}\right)}=\left\|\left(\sum_{j=1}^{n}\left|\mu_{j}^{R} * g\right|^{2}\right)^{1 / 2}\right\|_{L^{q_{0}}\left(\mathbb{R}^{n}\right)} \leqslant C\left(q_{0}, R, n\right)\|g\|_{L^{q_{0}}\left(\mathbb{R}^{n}\right)} .
$$

Recalling (8.15), (8.17) and that $M=M(\delta)$ depends on the fixed value $\delta>0$, we have when $R \geqslant R_{0}$ that

$$
\begin{equation*}
C\left(q_{0}, R, n\right) \leqslant \kappa_{q_{0}, \delta}+\mathrm{e} M(\delta)^{2} h_{q_{0}}^{M(\delta)}\left(b_{0}\left(q_{0}, R, n\right)+b_{1}\left(q_{0}, R, n\right)\right) \tag{8.23}
\end{equation*}
$$

where the three terms correspond to the decompositions (8.12) and (8.16). By definition, it will follow that the a priori bound $B\left(q_{0}, R, n\right)$ is less than
$C\left(q_{0}, R, n\right)$. Bounds on $b_{0}\left(q_{0}, R, n\right)$ and $b_{1}\left(q_{0}, R, n\right)$ will be obtained below, and will use the other quantities $B\left(q_{0}, R, \ell\right) \leqslant B\left(q_{0}, R, n\right)$, with $\ell \leqslant n$. We shall get a relation

$$
B\left(q_{0}, R, n\right) \leqslant c\left(q_{0}, \delta\right) R^{4 \delta}+B\left(q_{0}, R, n\right) / 2, \quad n \geqslant 1
$$

for $R$ larger than some $R_{1} \geqslant R_{0}$, and we shall be able to conclude.

### 8.3. Second reduction

Let $\tau>0$ be given. We say that a nonnegative function $f$ defined on $\mathbb{R}$ is $\tau$-stable with constant $C$ if whenever $|t| \leqslant \tau$, we have

$$
f(s+t) \leqslant C f(s), \quad s \in \mathbb{R}
$$

One sees that $C \geqslant 1$. Evident properties are to be observed about products, integrals, translations, convolutions... For example, if $f_{1}, \ldots, f_{k}$ are $\tau$-stable with respective constants $C_{i}$, then clearly the product $f_{1} \ldots f_{k}$ is $\tau$-stable with constant $C_{1} \ldots C_{k}$. If $f$ is $\tau$-stable with constant $C$ and if $g \geqslant 0$, then for $|t| \leqslant \tau$ we have

$$
\begin{align*}
(f * g)(s+t) & =\int_{\mathbb{R}} f(s+t-v) g(v) \mathrm{d} v  \tag{8.24}\\
& \leqslant C \int_{\mathbb{R}} f(s-v) g(v) \mathrm{d} v=C(f * g)(s)
\end{align*}
$$

hence $f * g$ is also $\tau$-stable with constant $C$. Suppose that $f, g, h$ are nonnegative on $\mathbb{R}$, and that $f$ is $\tau$-stable with constant $C$. If $|t| \leqslant \tau$ then

$$
\begin{aligned}
\int_{\mathbb{R}} f(s) g(s-t) h(t) \mathrm{d} s & \geqslant C^{-1} \int_{\mathbb{R}} f(s-t) g(s-t) h(t) \mathrm{d} s \\
& =C^{-1} h(t)\left(\int_{\mathbb{R}} f(v) g(v) \mathrm{d} v\right)
\end{aligned}
$$

therefore

$$
\begin{equation*}
\int_{\mathbb{R}} f(s)(g * h)(s) \mathrm{d} s \geqslant C^{-1}\left(\int_{|t| \leqslant \tau} h(t) \mathrm{d} t\right)\left(\int_{\mathbb{R}} f(v) g(v) \mathrm{d} v\right) \tag{8.25}
\end{equation*}
$$

We shall now move to $\mathbb{R}^{\ell}$ with $\ell \geqslant 1$. Let $\Phi$ be a probability density on $\mathbb{R}$ that is $\tau$-stable with constant $C$, for some $\tau>0$. This implies that $\Phi(s)>0$ for every $s \in \mathbb{R}$. Let us define $\beta \geqslant 1$ by

$$
\begin{equation*}
\beta^{-1}=\int_{|t| \leqslant \tau} \Phi(t) \mathrm{d} t \in(0,1) \tag{8.26}
\end{equation*}
$$

We denote by $\Phi_{j}$ the operator on $L^{q}\left(\mathbb{R}^{\ell}\right)$ of convolution with $\Phi$ in the variable $x_{j}$, for each $j \in L=\{1, \ldots, \ell\}$. For instance, when $j=1$ we let

$$
\left(\Phi_{1} f\right)\left(x_{1}, x_{2}, \ldots, x_{\ell}\right)=\int_{\mathbb{R}} f\left(x_{1}-s, x_{2}, \ldots, x_{\ell}\right) \Phi(s) \mathrm{d} s
$$

For $j=2, \ldots, \ell$ we let the transposition $\tau_{j}=(1 j)$ act on $x=\left(x_{1}, \ldots, x_{\ell}\right)$ in $\mathbb{R}^{\ell}$ by $\tau_{j} x=\left(x_{\tau_{j}(i)}\right)_{i=1}^{\ell}$ and on functions by $\tau_{j}(g)=g \circ \tau_{j}$. Letting $\tau_{1}=I$, we have

$$
\begin{equation*}
\Phi_{j} f=\tau_{j}\left(\Phi_{1}\left(f \circ \tau_{j}\right)\right), \quad j=1, \ldots, n \tag{8.27}
\end{equation*}
$$

For every subset $J \subset L$ we set $\Phi^{J}=\prod_{k \in J} \Phi_{k}$, and $\Phi^{\sim j}=\Phi^{L \backslash\{j\}}=$ $\prod_{k \neq j} \Phi_{k}$. We understand that $\Phi^{\emptyset}=I$. Each $\Phi^{J}$ is an operator acting on $L^{q}\left(\mathbb{R}^{\ell}\right)$ with norm equal to 1 , when $1 \leqslant q \leqslant+\infty$. The next Bourgain's lemma is not too difficult, but the details are long and painful to write down precisely. We have chosen to break it into two parts, the first one containing the serious work.

Lemma 8.7 (a first part of Bourgain's [13, Lemma 7]). - Let $\Phi$ be a probability density on $\mathbb{R}$ that is $\tau$-stable with constant $C$, let $\beta \geqslant 1$ be defined by (8.26). Let $\ell$ be an integer $\geqslant 1, L=\{1, \ldots, \ell\}$ and define $\Phi_{j}$ by (8.27), for $j=1, \ldots, \ell$. For all integers $q \geqslant 1$, for all nonnegative integrable functions $\left(f_{j}\right)_{j=1}^{\ell}$ on $\mathbb{R}^{\ell}$, one has

$$
\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}\right\|_{q} \leqslant \beta C^{q-1}\left\|\sum_{j \in L} \Phi^{L} f_{j}\right\|_{q}+\sqrt{q-1}\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}^{2}\right\|_{q / 2}^{1 / 2}
$$

Proof. - The fundamental remark compares

$$
\int_{\mathbb{R}}\left(\Phi_{1} g_{1}\right)(s)\left(\Phi_{1} g_{2}\right)(s) \ldots\left(\Phi_{1} g_{k-1}\right)(s) g_{k}(s) \mathrm{d} s
$$

and

$$
\int_{\mathbb{R}}\left(\Phi_{1} g_{1}\right)(s)\left(\Phi_{1} g_{2}\right)(s) \ldots\left(\Phi_{1} g_{k-1}\right)(s)\left(\Phi_{1} g_{k}\right)(s) \mathrm{d} s
$$

when $k \geqslant 2$ and when the functions $g_{j}$ s are nonnegative on $\mathbb{R}$. We know by (8.24) that $\Phi_{1} g=\Phi * g$ is $\tau$-stable with constant $C$ for every $g$ nonnegative, so the product $f=\left(\Phi_{1} g_{1}\right)\left(\Phi_{1} g_{2}\right) \ldots\left(\Phi_{1} g_{k-1}\right)$ is $\tau$-stable with constant $C^{k-1}$. Applying (8.25) and the definition of $\beta$ with $f, g=g_{k}, h=\Phi$ and $g * h=$ $\Phi_{1} g_{k}$, we get

$$
\begin{align*}
\int_{\mathbb{R}} & \left(\Phi_{1} g_{1}\right)(s)\left(\Phi_{1} g_{2}\right)(s) \ldots\left(\Phi_{1} g_{k-1}\right)(s) g_{k}(s) \mathrm{d} s \\
& \leqslant C^{k-1} \beta \int_{\mathbb{R}}\left(\Phi_{1} g_{1}\right)(s)\left(\Phi_{1} g_{2}\right)(s) \ldots\left(\Phi_{1} g_{k-1}\right)(s)\left(\Phi_{1} g_{k}\right)(s) \mathrm{d} s \tag{8.28}
\end{align*}
$$

The case $q=1$ of the lemma follows from $\beta \geqslant 1$ and $\int_{\mathbb{R}^{\ell}} g * f=\int_{\mathbb{R}^{\ell}} f$ for every probability density $g$. For the simplest non-trivial case, when $q=2$, we write

$$
\left(\sum_{j \in L} \Phi^{\sim j} f_{j}\right)^{2}=\sum_{i \neq j}\left(\Phi^{\sim i} f_{i}\right)\left(\Phi^{\sim j} f_{j}\right)+\sum_{j \in L}\left(\Phi^{\sim j} f_{j}\right)^{2}
$$

When $j \neq i$, the function $\Phi^{\sim i} f_{i}=\Phi^{L \backslash\{i\}} f_{i}=\Phi_{j} \Phi^{L \backslash\{i, j\}} f_{i}$ is of the form $\Phi_{j} g_{1}$, and letting $g_{2}=\Phi^{\sim j} f_{j}$ we get by (8.28) for the $x_{j}$ variable that

$$
\begin{aligned}
& \int_{\mathbb{R}}\left(\Phi^{\sim i} f_{i}\right)\left(\Phi^{\sim j} f_{j}\right) \mathrm{d} x_{j}=\int_{\mathbb{R}}\left(\Phi_{j} g_{1}\right) g_{2} \mathrm{~d} x_{j} \\
& \leqslant C \beta \int_{\mathbb{R}}\left(\Phi_{j} g_{1}\right)\left(\Phi_{j} g_{2}\right) \mathrm{d} x_{j}=C \beta \int_{\mathbb{R}}\left(\Phi^{\sim i} f_{i}\right)\left(\Phi^{L} f_{j}\right) \mathrm{d} x_{j}
\end{aligned}
$$

because $\Phi_{j} \Phi^{\sim j}=\Phi^{L}$. Integrating in the remaining variables, and since the functions are nonnegative, we obtain

$$
\begin{aligned}
& \int_{\mathbb{R}^{\ell}} \sum_{i \neq j}\left(\Phi^{\sim i} f_{i}\right)\left(\Phi^{\sim j} f_{j}\right) \mathrm{d} x \leqslant C \beta \int_{\mathbb{R}^{\ell}} \sum_{i \neq j}\left(\Phi^{\sim i} f_{i}\right)\left(\Phi^{L} f_{j}\right) \mathrm{d} x \\
& \leqslant C \beta \int_{\mathbb{R}^{\ell}}\left(\sum_{i \in L} \Phi^{\sim i} f_{i}\right)\left(\sum_{j \in L} \Phi^{L} f_{j}\right) \mathrm{d} x \leqslant C \beta\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}\right\|_{2}\left\|\sum_{j \in L} \Phi^{L} f_{j}\right\|_{2}
\end{aligned}
$$

When $j=i$, we use $\left(\Phi^{\sim j} * g\right)^{2} \leqslant \Phi^{\sim j} * g^{2}$ and get

$$
\int_{\mathbb{R}^{\ell}} \sum_{j \in L}\left(\Phi^{\sim j} f_{j}\right)^{2} \mathrm{~d} x \leqslant \int_{\mathbb{R}^{\ell}} \sum_{j \in L} \Phi^{\sim j} f_{j}^{2} \mathrm{~d} x=: B
$$

It follows that $E:=\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}\right\|_{2}$ satisfies an inequality $E^{2} \leqslant A E+B$, where we let $A:=C \beta\left\|\sum_{j \in L} \Phi^{L} f_{j}\right\|_{2}$. This yields $E \leqslant A+B^{1 / 2}$ and we have

$$
\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}\right\|_{2} \leqslant C \beta\left\|\sum_{j \in L} \Phi^{L} f_{j}\right\|_{2}+\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}^{2}\right\|_{1}^{1 / 2}
$$

This is Lemma 8.7 when $q=2$. In general, when $q \geqslant 3$, we expand

$$
\begin{equation*}
\int_{\mathbb{R}^{e}}\left(\sum_{j \in L} \Phi^{\sim j} f_{j}\right)^{q} \mathrm{~d} x=\sum_{j_{1}, j_{2}, \ldots, j_{q} \in L} \int_{\mathbb{R}^{\ell}}\left(\Phi^{\sim j_{1}} f_{j_{1}}\right) \ldots\left(\Phi^{\sim j_{q}} f_{j_{q}}\right) \mathrm{d} x \tag{8.29}
\end{equation*}
$$

Consider a multi-index $\left(j_{1}, j_{2}, \ldots, j_{q}\right) \in\{1, \ldots, \ell\}^{q}=L^{q}$ and suppose that $j_{q}$ is not equal to any of $j_{1}, \ldots, j_{q-1}$. Then we can write $\Phi^{\sim j_{k}} f_{j_{k}}=\Phi_{j_{q}} g_{k}$ for
each $k<q$, so as before, by (8.28) applied in the $x_{j_{q}}$ variable, we get that

$$
\begin{aligned}
& \int_{\mathbb{R}^{\ell}}\left(\Phi^{\sim j_{1}} f_{j_{1}}\right) \ldots\left(\Phi^{\sim j_{q}} f_{j_{q}}\right) \mathrm{d} x \\
& \leqslant C^{q-1} \beta \int_{\mathbb{R}^{\ell}}\left(\Phi^{\sim j_{1}} f_{j_{1}}\right) \ldots\left(\Phi^{\sim j_{q-1}} f_{j_{q-1}}\right)\left(\Phi^{L} f_{j_{q}}\right) \mathrm{d} x
\end{aligned}
$$

Let us denote by $\sum_{1}$ the part of the summation at the right-hand side of (8.29) that is extended to all $j_{1}, \ldots, j_{q}$ such that $j_{q} \notin\left\{j_{1}, \ldots, j_{q-1}\right\}$. We obtain that

$$
\begin{aligned}
& \sum_{1} \int_{\mathbb{R}^{\ell}}\left(\Phi^{\sim j_{1}} f_{j_{1}}\right) \ldots\left(\Phi^{\sim j_{q}} f_{j_{q}}\right) \mathrm{d} x \\
& \leqslant C^{q-1} \beta \int_{\mathbb{R}^{\ell}}\left(\sum_{j \in L} \Phi^{\sim j} f_{j}\right)^{q-1}\left(\sum_{j \in L} \Phi^{L} f_{j}\right) \mathrm{d} x \\
& \leqslant C^{q-1} \beta\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}\right\|_{q}^{q-1}\left\|\sum_{j \in L} \Phi^{L} f_{j}\right\|_{q}
\end{aligned}
$$

The remaining sum $\sum_{2}$ is less than the sum of $q-1$ terms corresponding to which index $j_{k}, k=1, \ldots, q-1$ is equal to $j_{q}$. Each of these $q-1$ terms is similar to

$$
\sum_{j_{1}, j_{2}, \ldots, j_{q-1} \in L} \int_{\mathbb{R}^{e}}\left(\Phi^{\sim j_{1}} f_{j_{1}}\right) \ldots\left(\Phi^{\sim j_{q-2}} f_{j_{q-2}}\right)\left(\Phi^{\sim j_{q-1}} f_{j_{q-1}}\right)^{2} \mathrm{~d} x
$$

which is bounded by

$$
\int_{\mathbb{R}^{\ell}}\left(\sum_{j \in L} \Phi^{\sim j} f_{j}\right)^{q-2}\left(\sum_{j \in L} \Phi^{\sim j} f_{j}^{2}\right) \leqslant\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}\right\|_{q}^{q-2}\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}^{2}\right\|_{q / 2}
$$

We obtain for $E_{q}=\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}\right\|_{q}^{q}$ a bound by $\Sigma_{1}+\Sigma_{2}$ of the form

$$
E_{q} \leqslant C^{q-1} \beta\left\|\sum_{j \in L} \Phi^{L} f_{j}\right\|_{q} E_{q}^{1-1 / q}+(q-1)\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}^{2}\right\|_{q / 2} E_{q}^{1-2 / q},
$$

which can be written also as

$$
E_{q}^{2 / q} \leqslant C^{q-1} \beta\left\|\sum_{j \in L} \Phi^{L} f_{j}\right\|_{q} E_{q}^{1 / q}+(q-1)\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}^{2}\right\|_{q / 2}
$$

This implies as before that

$$
\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}\right\|_{q}=E_{q}^{1 / q} \leqslant C^{q-1} \beta\left\|\sum_{j \in L} \Phi^{L} f_{j}\right\|_{q}+\sqrt{q-1}\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}^{2}\right\|_{q / 2}^{1 / 2}
$$

Lemma 8.8 ([13, Lemma 7$]$ ). - Let $\Phi$ be a probability density on $\mathbb{R}$ that is $\tau$-stable with constant $C$, and let $\beta \geqslant 1$ be defined by (8.26). Let $\ell \geqslant 1$ be an integer, $L=\{1, \ldots, \ell\}$ and define $\Phi_{j}$ by (8.27), for $j=1, \ldots, \ell$. For every integer $\nu \geqslant 1$ and for all nonnegative integrable functions $\left(f_{j}\right)_{j=1}^{\ell}$ on $\mathbb{R}^{\ell}$, one has

$$
\begin{equation*}
\kappa_{\nu}^{-1}\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}\right\|_{2^{\nu}} \leqslant \sum_{k=0}^{\nu}\left\|\sum_{j \in L} \Phi^{L} f_{j}^{2^{k}}\right\|_{2^{\nu-k}}^{2^{-k}} \leqslant(\nu+1)\left\|\sum_{j \in L} f_{j}\right\|_{2^{\nu}} \tag{8.30}
\end{equation*}
$$

with $\kappa_{\nu} \leqslant \max \left(2^{\nu}, \beta C^{2^{\nu}}\right)$. Each term $\left\|\sum_{j \in L} \Phi^{L} f_{j}^{2^{k}}\right\|_{2^{\nu-k}}^{2^{-k}}$, for $0 \leqslant k \leqslant \nu$, satisfies

$$
\left\|\sum_{j \in L} \Phi^{L} f_{j}^{2^{k}}\right\|_{2^{\nu-k}}^{2^{-k}}=\left\|\left(\sum_{j \in L} \Phi^{L} f_{j}^{2^{k}}\right)^{2^{-k}}\right\|_{2^{\nu}} \leqslant\left\|\sum_{j \in L} f_{j}\right\|_{2^{\nu}}
$$

Proof. - We begin with the easy last sentence. For $r=2^{k}$ and $k=$ $0, \ldots, \nu$, we have

$$
\left\|\sum_{j \in L} \Phi^{L} f_{j}^{r}\right\|_{2^{\nu-k}}^{1 / r} \leqslant\left\|\Phi^{L}\left(\sum_{j \in L} f_{j}\right)^{r}\right\|_{2^{\nu-k}}^{1 / r} \leqslant\left\|\left(\sum_{j \in L} f_{j}\right)^{r}\right\|_{2^{\nu-k}}^{1 / r}=\left\|\sum_{j \in L} f_{j}\right\|_{2^{\nu}}
$$

The constant in the right-hand inequality of (8.30) is therefore bounded by $\nu+1$.

We pass to the left-hand inequality. Let $q=2^{\nu}$. By Lemma 8.7, we can reduce the case $q=2^{\nu}$ to the case $q / 2$. We proceed by induction, with a number of steps bounded by $\nu$. Using $(a+b)^{\alpha} \leqslant a^{\alpha}+b^{\alpha}$ when $a, b \geqslant 0$ and $\alpha \in(0,1]$, we obtain

$$
\begin{aligned}
&\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}\right\|_{q} \leqslant \beta C^{q-1}\left\|\sum_{j \in L} \Phi^{L} f_{j}\right\|_{q}+\sqrt{q-1}\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}^{2}\right\|_{q / 2}^{1 / 2} \\
& \leqslant \beta C^{2^{\nu}}\left\|\sum_{j \in L} \Phi^{L} f_{j}\right\|_{2^{\nu}}+2^{\nu / 2} \beta^{1 / 2} C^{2^{\nu-2}}\left\|\sum_{j \in L} \Phi^{L} f_{j}^{2}\right\|_{2^{\nu-1}}^{1 / 2} \\
&+2^{\nu / 2} 2^{(\nu-1) / 4}\left\|\sum_{j \in L} \Phi^{\sim j} f_{j}^{4}\right\|_{2^{\nu-2}}^{1 / 4} \leqslant \ldots
\end{aligned}
$$

and the successive factors in front of $\left\|\sum_{j \in L} \Phi^{L} f_{j}^{2^{k}}\right\|_{2^{\nu-k}}^{2^{-k}}$, for $0 \leqslant k \leqslant \nu$, have the form $q(\beta / q)^{2^{-k}}\left(C^{q}\right)^{4^{-k}} \leqslant q\left(\beta C^{q} / q\right)^{2^{-k}}$, leading to $\kappa_{\nu} \leqslant$ $\max \left(q, \beta C^{q}\right)$.

We can try to optimize the constant $\kappa_{\nu}$ in the following way. Suppose that the function $\ln \Phi$ is Lipschitz on $\mathbb{R}$ with constant $\lambda$. Then we see that
$\Phi$ is $\tau$-stable with constant $C_{\tau}=\mathrm{e}^{\lambda \tau}$ for every $\tau>0$, and

$$
1 \geqslant \beta_{\tau}^{-1}:=\int_{|t| \leqslant \tau} \Phi(t) \mathrm{d} t \geqslant 2 \Phi(0) \int_{0}^{\tau} \mathrm{e}^{-\lambda t} \mathrm{~d} t=2 \Phi(0) \frac{1-\mathrm{e}^{-\lambda \tau}}{\lambda}
$$

Let $q=2^{\nu}$ and select $\tau=1 /(\lambda q)$. Then $C_{\tau} \leqslant \mathrm{e}^{1 / q}$ and

$$
\beta_{\tau} C_{\tau}^{q} \leqslant \frac{\mathrm{e} \lambda}{2 \Phi(0)\left(1-\mathrm{e}^{-1 / q}\right)} \leqslant \frac{\mathrm{e}^{2} \lambda}{2 \Phi(0)} q \leqslant \frac{4 \lambda}{\Phi(0)} q
$$

Coming back to Lemma 8.8 and noticing that $\lambda \geqslant 2 \Phi(0)$, we obtain

$$
\begin{equation*}
\kappa_{\nu} \leqslant \frac{4 \lambda}{\Phi(0)} 2^{\nu} \tag{8.31}
\end{equation*}
$$

We now introduce Bourgain's specific example $\varphi$ of a function $\Phi$, defined by

$$
\forall s \in \mathbb{R}, \quad \varphi(s)=\frac{c}{1+s^{4}},
$$

where $c=\sqrt{2} / \pi$ is chosen so that $\varphi$ is a probability density. This value $c$ is obtained by the residue theorem, which also gives the Fourier transform

$$
\widehat{\varphi}(t)=(\cos (\pi \sqrt{2}|t|)+\sin (\pi \sqrt{2}|t|)) \mathrm{e}^{-\pi \sqrt{2}|t|}, \quad t \in \mathbb{R}
$$

Notice that $(\cos u+\sin u) \mathrm{e}^{-u}=\sqrt{2} \cos (u-\pi / 4) \mathrm{e}^{-u} \geqslant \mathrm{e}^{-u^{2}}$ when $0 \leqslant u \leqslant$ $\pi / 2$, because $h(u)=\ln (\sqrt{2} \cos (u-\pi / 4))-u+u^{2} \geqslant 0$ on this interval. Indeed, we have $h(0)=h^{\prime}(0)=0$ and $h^{\prime \prime}(u)=1-\tan (u-\pi / 4)^{2} \geqslant 0$ on $[0, \pi / 2]$. It follows that

$$
\widehat{\varphi}(t) \geqslant \mathrm{e}^{-2 \pi^{2} t^{2}} \quad \text { when } \quad \pi \sqrt{2}|t| \leqslant \frac{\pi}{2}
$$

in particular when $\pi|t| \leqslant 1$. We shall need later the estimate given by Lemma 8.9.

Lemma 8.9. - For all $s \in \mathbb{R}, \ell$ integer $\geqslant 1$ and $\xi=\left(\xi_{1}, \ldots, \xi_{\ell}\right) \in \mathbb{R}^{\ell}$, one has that

$$
\left(1-\prod_{j=1}^{\ell} \widehat{\varphi}\left(s \xi_{j}\right)\right) \prod_{j=1}^{\ell} \widehat{\eta}\left(\xi_{j}\right) \leqslant 2 \pi^{2} s^{2}
$$

Proof. - Suppose that $\pi|s \xi| \leqslant 1$. Then $\pi\left|s \xi_{j}\right| \leqslant 1$ and $\widehat{\varphi}\left(s \xi_{j}\right) \geqslant$ $\mathrm{e}^{-2 \pi^{2} s^{2} \xi_{j}^{2}}$ for each index $j=1, \ldots, \ell$, thus by Lemma 8.6 we have

$$
\left(1-\prod_{j=1}^{\ell} \widehat{\varphi}\left(s \xi_{j}\right)\right) \prod_{j=1}^{\ell} \widehat{\eta}\left(\xi_{j}\right) \leqslant\left(1-\mathrm{e}^{-2 \pi^{2} s^{2}|\xi|^{2}}\right) \prod_{j=1}^{\ell} \widehat{\eta}\left(\xi_{j}\right) \leqslant 2 \pi^{2} s^{2}
$$

Otherwise, we have $\pi|s \xi| \geqslant 1$ and applying Lemma 8.6 with $r \rightarrow 0$ we get

$$
\left(1-\prod_{j=1}^{\ell} \widehat{\varphi}\left(s \xi_{j}\right)\right) \prod_{j=1}^{\ell} \widehat{\eta}\left(\xi_{j}\right) \leqslant 2 \prod_{j=1}^{\ell} \widehat{\eta}\left(\xi_{j}\right) \leqslant 2|\xi|^{-2} \leqslant 2 \pi^{2} s^{2}
$$

We know that $\varphi$ is 1-stable, because $F(x)=\ln \varphi(x)$ is Lipschitz. Indeed, its derivative $F^{\prime}(x)=-4 x^{3} /\left(1+x^{4}\right)$ is bounded on the real line. To be precise, the second derivative $F^{\prime \prime}$ vanishes when $x^{4}=3$, which implies that $\left|F^{\prime}(x)\right| \leqslant 3^{3 / 4}$ for every $x$. When $|t| \leqslant 1$, we have thus

$$
\varphi(s+t) \leqslant \mathrm{e}^{3^{3 / 4}} \varphi(s), \quad s \in \mathbb{R}
$$

with $\mathrm{e}^{3^{3 / 4}}<9,772<10$. This shows that $\varphi$ is 1 -stable with constant $\leqslant 10$. We shall need more than the 1-stability of the function $\varphi$, namely, we shall use the polynomial character of $1 / \varphi$. When $|t| \geqslant 1$ and $u \in \mathbb{R}$, we have

$$
\begin{equation*}
1+(u-t)^{4} \leqslant 1+8\left(u^{4}+t^{4}\right) \leqslant 8\left(1+t^{4}\right)\left(1+u^{4}\right) \leqslant 16 t^{4}\left(1+u^{4}\right) \tag{8.32}
\end{equation*}
$$

implying in this case, and with $u=s+t$, that $\varphi(s+t) \leqslant 16 t^{4} \varphi(s)$.
We introduce $w_{1}=w_{0}^{2}=R^{-\delta}<w_{0}$. The dilate $\varphi_{\left(w_{1}\right)}$ of $\varphi$ is $w_{1}$-stable with constant 10 and we shall consider from now on that $\Phi=\varphi_{\left(w_{1}\right)}$. We denote as before by $\Phi_{j}$ the operator on $L^{q}\left(\mathbb{R}^{\ell}\right)$ of convolution with $\varphi_{\left(w_{1}\right)}$ in the variable $x_{j}$, where $j \in L=\{1, \ldots, \ell\}$. For every subset $J \subset L$ we define $\Phi^{J}$ as before, as well as $\Phi^{\sim j}=\Phi^{L \backslash\{j\}}$. For $|t| \geqslant w_{1}$ we have by (8.32) the inequality

$$
\begin{equation*}
\varphi_{\left(w_{1}\right)}(s+t) \leqslant 16\left(t / w_{1}\right)^{4} \varphi_{\left(w_{1}\right)}(s)=16 R^{4 \delta} t^{4} \varphi_{\left(w_{1}\right)}(s), \quad s \in \mathbb{R} \tag{8.33}
\end{equation*}
$$

Here is perhaps the crux of the matter. The boundary measures $\mu_{j}$, partial derivatives of $\mu_{Q}$, will be swallowed and disappear as if by magic. The cube $Q$ here is the cube $Q_{\ell}$ in $\mathbb{R}^{\ell}$.

Lemma 8.10 ([13, Lemma 8]). - Let $\nu$ be an integer $\geqslant 1$, let $q=2^{\nu}$, and let $f_{1}, \ldots, f_{\ell}$ be functions in $L^{q}\left(\mathbb{R}^{\ell}\right)$. Let $\mu_{j}$ denote the partial derivative $\partial_{j} \mu_{Q}$ of the probability measure $\mu_{Q}$, for $j=1, \ldots, \ell$. With $\Phi_{j}$ defined as in (8.27) from $\Phi f=\varphi_{\left(w_{1}\right)} * f$ when $f \in L^{q}(\mathbb{R})$, one has that

$$
\left\|\left(\sum_{j=1}^{\ell}\left|\mu_{j} * \Phi^{\sim j} f_{j}\right|^{2}\right)^{1 / 2}\right\|_{L^{q}\left(\mathbb{R}^{\ell}\right)} \leqslant \kappa \sqrt{q \ln q} R^{4 \delta}\left\|\left(\sum_{j=1}^{\ell}\left|f_{j}\right|^{2}\right)^{1 / 2}\right\|_{L^{q}\left(\mathbb{R}^{\ell}\right)}
$$

Proof. - Let us write $L=\{1, \ldots, \ell\}$ and $\mathbb{R}^{L}$ for $\mathbb{R}^{\ell}$. For each $j \in L$, let $Q^{\sim j}$ denote the cube $\prod_{i \neq j}[-1 / 2,1 / 2]$ in $\mathbb{R}^{L \backslash\{j\}}$, let $\mathrm{d} \mathbf{x}^{\sim j}$ be the Lebesgue measure on $\mathbb{R}^{L \backslash\{j\}}$ and consider the probability measures $\tau_{j}, K^{\sim j}, K^{\{j\}}$ on
$\mathbb{R}^{\ell}$ defined by

$$
\begin{aligned}
\tau_{j} & =\frac{1}{2}\left(\delta_{1 / 2}\left(x_{j}\right)+\delta_{-1 / 2}\left(x_{j}\right)\right) \otimes\left(\otimes_{i \neq j} \delta_{0}\left(x_{i}\right)\right), \\
K^{\sim j} & =\delta_{0}\left(x_{j}\right) \otimes\left(\otimes_{i \neq j} \mathbf{1}_{[-1 / 2,1 / 2]}\left(x_{i}\right) \mathrm{d} x_{i}\right)=\delta_{0}\left(x_{j}\right) \otimes\left(\mathbf{1}_{Q^{\sim j}} \mathrm{~d} \mathbf{x}^{\sim j}\right), \\
K^{\{j\}} & =\left(\mathbf{1}_{[-1 / 2,1 / 2]}\left(x_{j}\right) \mathrm{d} x_{j}\right) \otimes\left(\otimes_{i \neq j} \delta_{0}\left(x_{i}\right)\right) .
\end{aligned}
$$

When convenient, we shall identify a kernel $K$ and the convolution operator with that kernel. Note that the signed measure $\mu_{j}=\partial_{j} \mathbf{1}_{Q}$ satisfies

$$
\left|\partial_{j} \mathbf{1}_{Q}\right|=\left(\delta_{1 / 2}\left(x_{j}\right)+\delta_{-1 / 2}\left(x_{j}\right)\right) \otimes\left(\mathbf{1}_{Q^{\sim j}} \mathrm{~d} \mathbf{x}^{\sim j}\right)=2 \tau_{j} * K^{\sim j}
$$

Using $|\mu * f|^{p} \leqslant \mu *|f|^{p}$ when $\mu$ is a probability measure and $p \geqslant 1$, we have

$$
\sum_{j=1}^{\ell}\left|\partial_{j} \mathbf{1}_{Q} * \Phi^{\sim j} f_{j}\right|^{2} \leqslant 4 \sum_{j=1}^{\ell} \Phi^{\sim j}\left(\tau_{j} * K^{\sim j} *\left|f_{j}\right|^{2}\right)
$$

We evaluate the $L^{q}$ norm applying Lemma 8.8, obtaining that

$$
\begin{aligned}
\left\|\left(\sum_{j=1}^{\ell}\left|\mu_{j} * \Phi^{\sim j} f_{j}\right|^{2}\right)^{1 / 2}\right\|_{q}^{2} & \leqslant 4\left\|\sum_{j=1}^{\ell} \Phi^{\sim j}\left(\tau_{j} * K^{\sim j} *\left|f_{j}\right|^{2}\right)\right\|_{q / 2} \\
& \leqslant 4 \kappa_{\nu-1} \sum_{k=0}^{\nu-1}\left(E_{k}\right)^{2^{-k}}
\end{aligned}
$$

where the expressions $E_{k}$ are given by

$$
E_{k}:=\left\|\sum_{j=1}^{\ell} \Phi^{L}\left[\tau_{j} * K^{\sim j} *\left|f_{j}\right|^{2}\right]^{2^{k}}\right\|_{q / 2^{k+1}}, \quad 1 \leqslant 2^{k} \leqslant q / 2=2^{\nu-1}
$$

Using again $|\mu * f|^{p} \leqslant \mu *|f|^{p}$ for $p \geqslant 1$, we get

$$
E_{k} \leqslant F_{k}:=\left\|\sum_{j=1}^{\ell} \Phi^{L}\left(\tau_{j} * K^{\sim j} *\left|f_{j}\right|^{2^{k+1}}\right)\right\|_{q / 2^{k+1}}
$$

Next, observe that $\varphi_{\left(w_{1}\right)}(s+t) \leqslant w_{1}^{-4} \varphi_{\left(w_{1}\right)}(s)=R^{4 \delta} \varphi_{\left(w_{1}\right)}(s)$ for $|t| \leqslant 1 / 2$. Indeed, when $w_{1} \leqslant|t| \leqslant 1 / 2$ we have $\varphi_{\left(w_{1}\right)}(s+t) \leqslant 16\left(t / w_{1}\right)^{4} \varphi_{\left(w_{1}\right)}(s) \leqslant$ $w_{1}^{-4} \varphi_{\left(w_{1}\right)}(s)$ by $(8.33)$, and $\varphi_{\left(w_{1}\right)}(s+t) \leqslant 10 \varphi_{\left(w_{1}\right)}(s) \leqslant R^{4 \delta} \varphi_{\left(w_{1}\right)}(s)$ when $|t| \leqslant w_{1}$, because we assumed that $R^{\delta} \geqslant R_{0}^{\delta}=16$. When $\mu$ is a probability measure supported on $[-1 / 2,1 / 2]$, it follows that $\mu * \varphi_{\left(w_{1}\right)} \leqslant R^{4 \delta} \varphi_{\left(w_{1}\right)}$ and $\varphi_{\left(w_{1}\right)} \leqslant R^{4 \delta} \mu * \varphi_{\left(w_{1}\right)}$. We have therefore $\tau_{j} \Phi_{j} \leqslant R^{4 \delta} \Phi_{j}$ and $\Phi_{j} \leqslant R^{4 \delta} \Phi_{j} K^{\{j\}}$. For $g$ nonnegative we obtain

$$
\Phi_{j} \tau_{j} g \leqslant R^{4 \delta} \Phi_{j} g \leqslant R^{8 \delta} \Phi_{j} K^{\{j\}} g
$$

Consequently, observing that $K^{\{j\}} * K^{\sim j}=\mathbf{1}_{Q}(x) \mathrm{d} x$, we have

$$
\Phi^{L}\left(\tau_{j} K^{\sim j} g\right)=\Phi^{\sim j} \Phi_{j} \tau_{j} K^{\sim j} g \leqslant R^{8 \delta} \Phi^{\sim j} \Phi_{j} K^{\{j\}} K^{\sim j} g=R^{8 \delta} \Phi^{L} * \mathbf{1}_{Q} * g
$$

and by the last assertion of Lemma 8.8, we obtain for $k=0, \ldots, \nu-1$ that

$$
\begin{aligned}
F_{k} & \leqslant R^{8 \delta}\left\|\left(\sum_{j=1}^{\ell} \Phi^{L}\left|f_{j}\right|^{2^{k+1}}\right) * \mathbf{1}_{Q}\right\|_{q / 2^{k+1}} \\
& \leqslant R^{8 \delta}\left\|\sum_{j=1}^{\ell} \Phi^{L}\left|f_{j}\right|^{2^{k+1}}\right\|_{q / 2^{k+1}} \leqslant R^{8 \delta}\left\|\left(\sum_{j=1}^{\ell}\left|f_{j}\right|^{2}\right)^{1 / 2}\right\|_{q}^{2^{k+1}}
\end{aligned}
$$

Finally, assuming $\left\|\left(\sum_{j=1}^{\ell}\left|f_{j}\right|^{2}\right)^{1 / 2}\right\|_{q} \leqslant 1$ we get

$$
\left\|\left(\sum_{j=1}^{\ell}\left|\mu_{j} * \Phi^{\sim j} f_{j}\right|^{2}\right)^{1 / 2}\right\|_{q}^{2} \leqslant 4 \kappa_{\nu-1} \sum_{k=0}^{\nu-1}\left(R^{8 \delta}\right)^{2^{-k}} \leqslant 4 \nu \kappa_{\nu-1} R^{8 \delta}
$$

Since $\ln \varphi$ is Lipschitz on $\mathbb{R}$, we can estimate $\kappa_{\nu}$ by (8.31) and conclude.
Recalling that $G^{R}$ is a probability density and $\mu_{j}^{R}=\mu_{j} * G^{R}$, we immediately deduce the result that we really need.

Lemma 8.11 ([13, Lemma 9]). - Assume that $q=2^{\nu}$, with $\nu \geqslant 1$ an integer. Let $f_{1}, \ldots, f_{\ell}$ be elements of $L^{q}\left(\mathbb{R}^{\ell}\right)$. With $\Phi_{j}$ as in Lemma 8.10, we have

$$
\left\|\left(\sum_{j=1}^{\ell}\left|\mu_{j}^{R} * \Phi^{\sim j} f_{j}\right|^{2}\right)^{1 / 2}\right\|_{L^{q}\left(\mathbb{R}^{\ell}\right)} \leqslant \kappa_{q} R^{4 \delta}\left\|\left(\sum_{j=1}^{\ell}\left|f_{j}\right|^{2}\right)^{1 / 2}\right\|_{L^{q}\left(\mathbb{R}^{\ell}\right)}
$$

### 8.4. Conclusion

It remains to estimate the two terms $\mathbf{E}(R, \ell, f):=\left|\nabla \mu^{R} * H_{0} f\right|$ and $\mathbf{F}(R, \ell, f)$ defined in (8.22), for $f \in L^{q_{0}}\left(\mathbb{R}^{\ell}\right), q_{0}=2^{\nu}$ and $\ell \leqslant n$. Each one will be cut into two pieces, one of order a power of $R^{\delta}$ and the second bounded by a "small" multiple of $B\left(q_{0}, R, n\right)$. Let us start with $\mathbf{E}(R, \ell, f)$, and cut $\left|\nabla \mu^{R} * H_{0} f\right|$ into
$\mathbf{E}^{\prime}(R, \ell, f):=\left|\nabla \mu^{R} * G_{\left(w_{1}\right)} * H_{0} f\right|, \mathbf{E}^{\prime \prime}(R, \ell, f):=\left|\nabla \mu^{R} *\left(\delta_{0}-G_{\left(w_{1}\right)}\right) * H_{0} f\right|$.
We begin with $\mathbf{E}^{\prime}(R, \ell, f)$. The mapping $f \mapsto \nabla \mu^{R} * G_{\left(w_{1}\right)} * H_{0} f$, equal to $U_{\mu^{R} * G_{\left(w_{1}\right)}} \circ H_{0}$ is studied by applying Lemma 7.11 to the log-concave probability density $\mu^{R} * G_{\left(w_{1}\right)}$. Using (7.3) and (8.1), we see that $V\left(\mu^{R} *\right.$ $\left.G_{\left(w_{1}\right)}\right) \leqslant V\left(G_{\left(w_{1}\right)}\right)=w_{1}^{-1}=R^{\delta}$. The variance of $\mu^{R} * G_{\left(w_{1}\right)}=\mu_{Q} * G_{(1 / R)} *$
$G_{\left(w_{1}\right)}$ is larger than that of $\mu_{Q}$, which is equal to (12) ${ }^{-1}$. By Lemma 7.11 and $q_{0} \geqslant 2$, we get that

$$
\left\|\mathbf{E}^{\prime}(R, \ell, f)\right\|_{q_{0}} \leqslant 24^{1 / q_{0}}\left(R^{\delta}\right)^{1-2 / q_{0}}\left\|H_{0} f\right\|_{q_{0}} \leqslant 5 R^{\delta}\|f\|_{q_{0}} .
$$

We study now $\mathbf{E}^{\prime \prime}(R, \ell, f)$ with the a priori estimate that involves the constant $B\left(q_{0}, R, \ell\right)$. By the definition (8.7), one writes

$$
\begin{aligned}
\left\|\mathbf{E}^{\prime \prime}(R, \ell, f)\right\|_{q_{0}} & =\left\|\nabla \mu^{R} *\left(\delta_{0}-G_{\left(w_{1}\right)}\right) * H_{0} f\right\|_{q_{0}} \\
& \leqslant B\left(q_{0}, R, \ell\right)\left\|\left(\delta_{0}-G_{\left(w_{1}\right)}\right) * H_{0} f\right\|_{q_{0}} .
\end{aligned}
$$

We continue by interpolation $\left(L^{\infty}, L^{2}\right)$ for $f \mapsto\left(\delta_{0}-G_{\left(w_{1}\right)}\right) * H_{0} f$. In $L^{\infty}\left(\mathbb{R}^{\ell}\right)$ one has simply $\left\|\left(\delta_{0}-G_{\left(w_{1}\right)}\right) * H_{0}\right\|_{\infty \rightarrow \infty} \leqslant 2$ by using the $L^{1}$ norm of the convolution kernel. Lemma 8.6 with $r=2 \sqrt{\pi} w_{1} / w_{0}$ gives for the Fourier transform a bound

$$
\left(1-\mathrm{e}^{-4 \pi w_{1}^{2}|\xi|^{2}}\right) \prod_{j=1}^{\ell} \widehat{\eta}\left(w_{0} \xi_{j}\right) \leqslant 4 \pi\left(w_{1} / w_{0}\right)^{2}=4 \pi w_{0}^{2}=4 \pi R^{-\delta}, \quad \xi \in \mathbb{R}^{\ell}
$$

implying $\left\|\left(\delta-G_{\left(w_{1}\right)}\right) * H_{0}\right\|_{2 \rightarrow 2} \leqslant 4 \pi R^{-\delta}$. We get in this way that

$$
\left\|\left(\delta_{0}-G_{\left(w_{1}\right)}\right) * H_{0}\right\|_{q_{0} \rightarrow q_{0}} \leqslant 2^{1-2 / q_{0}}\left(4 \pi R^{-\delta}\right)^{2 / q_{0}} \leqslant 4 \pi R^{-2 \delta / q_{0}}
$$

thus $\left\|\mathbf{E}^{\prime \prime}(R, \ell, f)\right\|_{q_{0}} \leqslant \kappa B\left(q_{0}, R, \ell\right) R^{-2 \delta / q_{0}}\|f\|_{q_{0}}$ and we obtain

$$
\|\mathbf{E}(R, \ell, f)\|_{q_{0}} \leqslant \kappa\left(R^{\delta}+B\left(q_{0}, R, \ell\right) R^{-2 \delta / q_{0}}\right)\|f\|_{q_{0}}
$$

Now we consider $\mathbf{F}(R, \ell, f)$ and we cut it into

$$
\begin{aligned}
\mathbf{F}^{\prime}(R, \ell, f) & :=\left(\sum_{j=1}^{\ell}\left|\mu_{j}^{R} * \Gamma^{j} \Phi^{\sim j} f\right|^{2}\right)^{1 / 2} \\
\mathbf{F}^{\prime \prime}(R, \ell, f) & :=\left(\sum_{j=1}^{\ell}\left|\mu_{j}^{R} * \Gamma^{j}\left(I-\Phi^{\sim j}\right) f\right|^{2}\right)^{1 / 2}
\end{aligned}
$$

By Lemma 8.11, we have that

$$
\left\|\mathbf{F}^{\prime}(R, \ell, f)\right\|_{q_{0}} \leqslant \kappa_{q_{0}} R^{4 \delta}\left\|\left(\sum_{j=1}^{\ell}\left|\Gamma^{j} f\right|^{2}\right)^{1 / 2}\right\|_{q_{0}}
$$

Using Khinchin's (1.22.K) and (1.27) we reduce to $\left\|\sum_{j=1}^{\ell} \pm \Gamma^{j} f\right\|_{q_{0}}$, and dividing according to the sign $\pm$, we further reduce to $\left\|\sum_{j \in J_{1}} \Gamma^{j} f\right\|_{q_{0}}$ and $\left\|\sum_{j \notin J_{1}} \Gamma^{j} f\right\|_{q_{0}}$, where $J_{1} \subset\{1, \ldots, \ell\}$. The first sum corresponds to the operator $H_{1}$ relative to $J_{1}$, the second is the one for $\sim J_{1}:=\{1, \ldots, \ell\} \backslash J_{1}$.

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By Proposition 8.3 for the set $J_{1}$ of variables, writing $x=\left(\mathbf{x}^{J_{1}}, \mathbf{x}^{\sim J_{1}}\right) \in \mathbb{R}^{\ell}$, we have for $1<q<+\infty$ that

$$
\left\|\left(\sum_{j \in J_{1}} \Gamma^{j} f\right)_{\mathbf{x}^{\sim J_{1}}}\right\|_{L^{q}\left(\mathbb{R}^{J_{1}}\right)} \leqslant h_{q}\left\|f_{\mathbf{x}^{\sim J_{1}}}\right\|_{L^{q}\left(\mathbb{R}^{J_{1}}\right)}, \quad \mathbf{x}^{\sim J_{1}} \in \mathbb{R}^{\sim J_{1}}
$$

and integrating in the variables in $\sim J_{1}$ we get with $A_{q}$ from (1.22.K) that

$$
\begin{equation*}
\left\|\left(\sum_{j=1}^{\ell}\left|\Gamma^{j} f\right|^{2}\right)^{1 / 2}\right\|_{q} \leqslant 2 A_{q}^{-1} h_{q}\|f\|_{q} \tag{8.34}
\end{equation*}
$$

for $1<q<+\infty$. It follows that $\left\|\mathbf{F}^{\prime}(R, \ell, f)\right\|_{q_{0}} \leqslant \kappa_{q_{0}}^{\prime} R^{4 \delta}\|f\|_{q_{0}}$.
For the second term $\mathbf{F}^{\prime \prime}(R, \ell, f)$ we first obtain an $L^{2}$ bound for the nonlinear operator $V: f \mapsto\left(\sum_{j=1}^{\ell}\left|\Gamma^{j}\left(I-\Phi^{\sim j}\right) f\right|^{2}\right)^{1 / 2}$, by estimating the Fourier transform

$$
\chi(\xi):=\sum_{j=1}^{\ell}\left(1-\widehat{\eta}\left(w_{0} \xi_{j}\right)\right)^{2}\left(\prod_{i \neq j} \widehat{\eta}\left(w_{0} \xi_{i}\right)\right)^{2}\left(1-\prod_{i \neq j} \widehat{\varphi}\left(w_{1} \xi_{i}\right)\right)^{2} \leqslant 4 \pi^{2} R^{-\delta}
$$

Indeed, we know that $0 \leqslant \widehat{\eta}(t) \leqslant 1$ and $-1 \leqslant \widehat{\varphi}(t) \leqslant 1$, therefore

$$
\begin{aligned}
& \sum_{j=1}^{\ell}\left(1-\widehat{\eta}\left(w_{0} \xi_{j}\right)\right)^{2}\left(\prod_{i \neq j} \widehat{\eta}\left(w_{0} \xi_{i}\right)\right)\left(1-\prod_{i \neq j} \widehat{\varphi}\left(w_{1} \xi_{i}\right)\right) \\
& \quad \leqslant 2 \sum_{j=1}^{\ell}\left(1-\widehat{\eta}\left(w_{0} \xi_{j}\right)\right) \prod_{i \neq j} \widehat{\eta}\left(w_{0} \xi_{i}\right) \leqslant 2 \prod_{j=1}^{\ell}\left(\left(1-\widehat{\eta}\left(w_{0} \xi_{j}\right)\right)+\widehat{\eta}\left(w_{0} \xi_{j}\right)\right)=2
\end{aligned}
$$

and by Lemma 8.9 applied to $\mathbb{R}^{L \backslash\{j\}}$ with $s=w_{1} / w_{0}=w_{0}=R^{-\delta / 2}$, it follows that

$$
\chi(\xi) \leqslant 2 \max _{1 \leqslant j \leqslant \ell}\left(\prod_{i \neq j} \widehat{\eta}\left(w_{0} \xi_{i}\right)\right)\left(1-\prod_{i \neq j} \widehat{\varphi}\left(w_{1} \xi_{i}\right)\right) \leqslant 4 \pi^{2} R^{-\delta}, \quad \xi \in \mathbb{R}^{\ell}
$$

We get $\|V f\|_{2}^{2} \leqslant 4 \pi^{2} R^{-\delta}\|f\|_{2}^{2}$ and $\|V\|_{2 \rightarrow 2} \leqslant 2 \pi R^{-\delta / 2}$. On the other hand, given functions $\left(g_{j}\right)_{j=1}^{\ell}$ and independent Bernoulli random variables $\left(\varepsilon_{j}\right)_{j=1}^{\ell}$, we have

$$
\begin{aligned}
& 2\left(\sum_{j=1}^{\ell}\left|\mu_{j}^{R} * g_{j}\right|^{2}\right)^{1 / 2}=2\left(\sum_{j=1}^{\ell}\left|\mu_{j}^{R} * \varepsilon_{j} g_{j}\right|^{2}\right)^{1 / 2} \\
& \quad \leqslant\left(\sum_{j=1}^{\ell}\left|\mu_{j}^{R} *\left(\varepsilon_{j} g_{j}+\sum_{i \neq j} \varepsilon_{i} g_{i}\right)\right|^{2}\right)^{1 / 2}+\left(\sum_{j=1}^{\ell}\left|\mu_{j}^{R} *\left(\varepsilon_{j} g_{j}-\sum_{i \neq j} \varepsilon_{i} g_{i}\right)\right|^{2}\right)^{1 / 2}
\end{aligned}
$$

hence with $g_{j}=\Gamma^{j}\left(I-\Phi^{\sim j}\right) f$ and $F_{\varepsilon}=\sum_{i=1}^{\ell} \varepsilon_{i} \Gamma^{i}\left(I-\Phi^{\sim i}\right) f$ we see that

$$
\left\|\mathbf{F}^{\prime \prime}(R, \ell, f)\right\|_{q_{0}} \leqslant \mathrm{E}_{\varepsilon}\left\|\left(\sum_{j=1}^{\ell}\left|\mu_{j}^{R} * F_{\varepsilon}\right|^{2}\right)^{1 / 2}\right\|_{q_{0}}=\mathrm{E}_{\varepsilon}\left\|\nabla \mu^{R} * F_{\varepsilon}\right\|_{q_{0}}=: D
$$

With Khinchin (1.27) and the a priori bound (8.7) we obtain

$$
D \leqslant B\left(q_{0}, R, \ell\right) \mathrm{E}_{\varepsilon}\left\|F_{\varepsilon}\right\|_{q_{0}} \leqslant B_{q_{0}} B\left(q_{0}, R, \ell\right)\left\|\left(\sum_{i=1}^{\ell}\left|\Gamma^{i}\left(I-\Phi^{\sim i}\right) f\right|^{2}\right)^{1 / 2}\right\|_{q_{0}}
$$

In $L^{q_{1}}\left(\mathbb{R}^{\ell}\right)$ with $q_{1}=2 q_{0}=2^{\nu+1}$ we have by (8.34) and Lemma 8.8 that

$$
\begin{aligned}
& \left\|\left(\sum_{j=1}^{\ell}\left|\Gamma^{j}\left(I-\Phi^{\sim j}\right) f\right|^{2}\right)^{1 / 2}\right\|_{q_{1}} \\
& \leqslant\left\|\left(\sum_{j=1}^{\ell}\left|\Gamma^{j} f\right|^{2}\right)^{1 / 2}\right\|_{q_{1}}+\left\|\sum_{j=1}^{\ell} \Phi^{\sim j}\left(\Gamma^{j} f\right)^{2}\right\|_{q_{1} / 2}^{1 / 2} \leqslant \kappa_{q_{0}}\|f\|_{q_{1}}
\end{aligned}
$$

Interpolating with the $L^{2}$ bound, and with $\kappa_{q_{0}}$ changing from line to line, we get

$$
\left\|\left(\sum_{j=1}^{\ell}\left|\Gamma^{j}\left(I-\Phi^{\sim j}\right) f\right|^{2}\right)^{1 / 2}\right\|_{q_{0}} \leqslant \kappa_{q_{0}} R^{-\delta /\left(2 q_{0}-2\right)}\|f\|_{q_{0}} \leqslant \kappa_{q_{0}} R^{-\delta /\left(2 q_{0}\right)}\|f\|_{q_{0}}
$$

$$
\text { therefore }\left\|\mathbf{F}^{\prime \prime}(R, \ell, f)\right\|_{q_{0}} \leqslant \kappa_{q_{0}} B\left(q_{0}, R, \ell\right) R^{-\delta /\left(2 q_{0}\right)}\|f\|_{q_{0}} \text { and }
$$

$$
\|\mathbf{F}(R, \ell, f)\|_{q_{0}} \leqslant \kappa_{q_{0}}\left(R^{4 \delta}+B\left(q_{0}, R, \ell\right) R^{-\delta /\left(2 q_{0}\right)}\right)\|f\|_{q_{0}}
$$

The estimates are proved in every dimension $\ell \leqslant n$, we have thus realized our objectives (8.21) and (8.22). Noticing that $R \geqslant 1$, we have consequently

$$
b_{0}\left(q_{0}, R, n\right)+b_{1}\left(q_{0}, R, n\right) \leqslant \kappa_{q_{0}}\left(R^{4 \delta}+B\left(q_{0}, R, n\right) R^{-\delta /\left(2 q_{0}\right)}\right)
$$

At last, we put all parts of (8.23) together. We may assume that $4 \delta<1$. We use again $R \geqslant 1$ in order to absorb the constant bound from (8.15), thus obtaining

$$
\left\|\nabla \mu^{R} * g\right\|_{q_{0}} \leqslant c\left(q_{0}, \delta\right)\left(R^{4 \delta}+B\left(q_{0}, R, n\right) R^{-\delta /\left(2 q_{0}\right)}\right)\|g\|_{q_{0}}
$$

for $g \in L^{q_{0}}\left(\mathbb{R}^{n}\right)$ and $R \geqslant R_{0}$. Since $B\left(q_{0}, R, n\right)$ is the maximum of $\| \nabla \mu^{R} *$ $g \|_{q_{0}}$ for $g$ of norm $\leqslant 1$ in $L^{q_{0}}\left(\mathbb{R}^{n}\right)$, we deduce that $B\left(q_{0}, R, n\right) \leqslant c\left(q_{0}, \delta\right) R^{4 \delta}+$ $B\left(q_{0}, R, n\right) / 2$ for $R \geqslant R_{1}$, if $R_{1} \geqslant R_{0}$ is such that $c\left(q_{0}, \delta\right) R_{1}^{-\delta /\left(2 q_{0}\right)} \leqslant$ $1 / 2$, thus $B\left(q_{0}, R, n\right) \leqslant 2 c\left(q_{0}, \delta\right) R^{4 \delta}$ for $R \geqslant R_{1}$. The value of $R_{1}$ depends on $\delta$ and $q_{0}$ that are fixed. For $R \leqslant R_{1}$, we may estimate directly
$\left\|\nabla \mu^{R} * g\right\|_{q_{0}} \leqslant \kappa R\|g\|_{q_{0}} \leqslant \kappa R_{1}^{1-4 \delta} R^{4 \delta}\|g\|_{q_{0}}$ by Lemma 7.11. It follows finally that $B\left(q_{0}, R, n\right) \leqslant c^{\prime}\left(q_{0}, \delta\right) R^{4 \delta}$, and $\delta$ being arbitrarily small, we have proved Proposition 8.2.

## 9. The Aldaz weak type result for cubes, and improvements

We work again in this section with the symmetric cube $Q_{n}$ of volume 1 in $\mathbb{R}^{n}$, that is to say, with $Q_{1}=[-1 / 2,1 / 2]$ when $n=1$ and $Q_{n}=\left(Q_{1}\right)^{n}$. We first present, following Aubrun [3], a rather soft argument proving the result of Aldaz [1] that the weak type $(1,1)$ constant $\kappa_{Q, n}$ associated to the cubes $Q_{n}$ is not bounded when $n$ tends to infinity. We shall indicate and comment the quantitative improvement obtained by Aubrun [3], who gave a lower bound $\kappa_{Q, n} \geqslant \kappa_{\varepsilon}(\log n)^{1-\varepsilon}$ for every $\varepsilon>0$. We then give a version of the proof of Iakovlev and Strömberg [46] who considerably improved this lower bound, showing that $\kappa_{Q, n} \geqslant \kappa n^{1 / 4}$. All the arguments though are based on the same initial principle that we now recall.

We begin with a few simple reflections. If we want to contradict the uniform boundedness of the weak type $(1,1)$ constant $\kappa_{Q, n}$ we must, in view of Bourgain's Theorem 8.1, look for functions $f_{n}$ on $\mathbb{R}^{n}$ that do not stay bounded, as $n \rightarrow \infty$, in any $L^{p}\left(\mathbb{R}^{n}\right)$ with $p>1$. Also, we may easily obtain by mollifying techniques that the weak type inequality for $L^{1}$ functions, stating that

$$
c\left|\left\{x \in \mathbb{R}^{n}:\left(\mathrm{M}_{Q} f\right)(x)>c\right\}\right| \leqslant \kappa_{Q, n}\|f\|_{L^{1}\left(\mathbb{R}^{n}\right)}, \quad c>0, f \in L^{1}\left(\mathbb{R}^{n}\right)
$$

where we let $\mathrm{M}_{Q}=\mathrm{M}_{Q_{n}}$, extends to bounded nonnegative measures $\mu$ on $\mathbb{R}^{n}$ : if for every $x \in \mathbb{R}^{n}$ we define $\left(\mathrm{M}_{Q} \mu\right)(x)$ to be the supremum over $r>0$ of all quotients $\mu(x+r Q) /|x+r Q|$, then (9.1) extends with the same constant $\kappa_{Q, n}$ as

$$
c\left|\left\{x \in \mathbb{R}^{n}:\left(\mathrm{M}_{Q} \mu\right)(x)>c\right\}\right| \leqslant \kappa_{Q, n} \mu\left(\mathbb{R}^{n}\right), \quad c>0
$$

These two remarks lead naturally to consider measures on $\mathbb{R}^{n}$ that are sums of Dirac measures, in order to contradict the boundedness of $\kappa_{Q, n}$ when $n \rightarrow \infty$.

Let $\mu_{N}=\sum_{j=1}^{N}\left(\delta_{j-1 / 2}+\delta_{-j+1 / 2}\right)$ stand for an "approximation" of the Lebesgue measure $\lambda$ on a large segment $S_{N}=[-N, N]$. The measure $\mu_{N}$ has a unit mass at the middle of each interval $(j, j+1), j$ integer and $-N \leqslant j<$ $N$. Every interval $[u, u+h)$ contained in $S_{N}$, with length an integer $h>0$, has the same measure $h$ for $\mu_{N}$ or for $\lambda$. However, if $I$ is a segment of length $1+\alpha, 0<\alpha=1-\varepsilon<1$, centered at $s=0$ or at any $s=j$, integer with $|j|<N$, then $I$ contains $j \pm 1 / 2$ and

$$
\mu_{N}(I)=2 \quad \text { but } \quad \lambda(I)=1+\alpha=2-\varepsilon<2
$$

so that $\left(\mathrm{M}_{Q} \mu_{N}\right)(s) \geqslant \mu_{N}(I) / \lambda(I)=2 /(2-\varepsilon)$. The same observation is valid if $s$ is not too far from an integer $j$ in $(-N, N)$, precisely, if $|s-j|<\alpha / 2$. If we pass to $\mathbb{R}^{n}$ and to the tensor product measure $\mu_{N}^{(n)}:=\otimes^{n} \mu_{N}$, we obtain a huge magnification due to dimension, which reads as

$$
\left(\mathrm{M}_{Q} \mu_{N}^{(n)}\right)(x) \geqslant\left(\frac{2}{2-\varepsilon}\right)^{n}
$$

when all coordinates $x_{i}, i=1, \ldots, n$, of the point $x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n}$ belong to the subset $C_{\alpha}$ of $[-N, N]$ defined by

$$
\begin{equation*}
C_{\alpha}=\bigcup_{-N<j<N}(j-\alpha / 2, j+\alpha / 2) \tag{9.2}
\end{equation*}
$$

If $\ell=2 h+1 \geqslant 1$ is an odd integer, if $J=(-h-1 / 2-\alpha / 2, h+1 / 2+\alpha / 2)=$ $(\ell+\alpha) Q$ and if $s+J$ is contained in $S_{N}$, we see in the same way, when $s \in C_{\alpha}$, that the segment $s+J$ contains $\ell+1=2 h+2$ of the unit masses forming $\mu_{N}$. Consequently, we have $\left(\mathrm{M}_{Q} \mu_{N}\right)(s) \geqslant(\ell+1) /|s+J|=(\ell+1) /(\ell+\alpha)$ and

$$
\left(\mathrm{M}_{Q} \mu_{N}^{(n)}\right)(x) \geqslant\left(\frac{\ell+1}{\ell+1-\varepsilon}\right)^{n}
$$

when $x=\left(x_{1}, \ldots, x_{n}\right)$ has all coordinates $x_{i}$ in $C_{\alpha}$ and $x+J^{n} \subset S_{N}^{n}=$ $2 N Q_{n}$.

This case is much too particular, since the set of such points $x$ represents only a tiny proportion $\alpha^{n}$ of the cube $S_{N}^{n}$. One has actually to consider that some coordinates $x_{i}$ of $x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n}$ are in $C_{\alpha}$, say $m \leqslant n$ of them. For the other coordinates $x_{i}$, observe that any interval of length $\ell+\alpha$ contained in $S_{N}$ contains at least $\ell$ points of the support of $\mu_{N}$. Assuming that $x+(\ell+\alpha) Q \subset S_{N}^{n}$, we get for this point $x$ with $m$ coordinates in $C_{\alpha}$ the lower bound

$$
\begin{equation*}
\left(\mathrm{M}_{Q} \mu_{N}^{(n)}\right)(x) \geqslant \frac{\mu_{N}^{(n)}(x+(\ell+\alpha) Q)}{|x+(\ell+\alpha) Q|} \geqslant\left(\frac{\ell+1}{\ell+\alpha}\right)^{m}\left(\frac{\ell}{\ell+\alpha}\right)^{n-m} \tag{9.3}
\end{equation*}
$$

We want the cardinality $m$ of the "good", "centered" coordinates $x_{i}$ to be as big as possible. Since they are chosen out of subsets of length $\alpha$ in unit intervals $(j-1 / 2, j+1 / 2)$, it is likely that the proportion of "good coordinates" among $n$ coordinates be around $\alpha$, with a plausible deviation of order $\sqrt{n}$ from the expected number $\alpha n$. We shall thus think henceforth that $m=\alpha n+\delta \sqrt{n}$ for some $\delta>0$.

We try to make the lower bound (9.3) as large as possible, by a suitable choice of $\ell$. Setting $\beta=1-\alpha$, we rewrite the right-hand side of (9.3) under
the form

$$
\left(\frac{\ell+\alpha+\beta}{\ell+\alpha}\right)^{m}\left(\frac{\ell+\alpha-\alpha}{\ell+\alpha}\right)^{n-m}=\left(1+\frac{\beta}{\ell+\alpha}\right)^{m}\left(1-\frac{\alpha}{\ell+\alpha}\right)^{n-m}
$$

Considering now $y=(\ell+\alpha)^{-1}$ as a real parameter, we will study

$$
V(y):=(1+\beta y)^{m}(1-\alpha y)^{n-m}, \quad-1 / \beta \leqslant y \leqslant 1 / \alpha
$$

and find the maximal value $V(\bar{y})$. Equivalently, we let $f$ denote the fraction $m / n$ of coordinates of $x$ that are in $C_{\alpha}$, and we maximize $v_{f, \alpha}(s)=V(s)^{1 / n}$ defined by

$$
v_{f, \alpha}(s)=(1+\beta s)^{f}(1-\alpha s)^{1-f}, \quad s \in[-1 / \beta, 1 / \alpha] .
$$

We have to remember though that the lower bound $V(y)$ for $\mathrm{M}_{Q} \mu_{N}^{(n)}(x)$ given in (9.3) is only valid when $1 / y-\alpha$ is an odd integer $\ell$. We shall replace $\bar{y}$ by a value $y=y_{\mathbb{N}}>0$ close to $\bar{y}$, such that $1 / y_{\mathbb{N}}-\alpha$ is an odd integer, thus obtaining that $\mathrm{M}_{Q} \mu_{N}^{(n)}(x) \geqslant V\left(y_{\mathbb{N}}\right)$. We must ensure that the value of $V(y)$ does not decrease too much when moving from $\bar{y}$ to $y_{\mathbb{N}}$. We would like to have

$$
\begin{equation*}
V\left(y_{\mathbb{N}}\right) \geqslant \mathrm{e}^{-c} V(\bar{y}) \text { or } v_{f, \alpha}\left(y_{\mathbb{N}}\right) \geqslant \mathrm{e}^{-c / n} v_{f, \alpha}(\bar{y}), \text { for some } c \geqslant 0 \tag{9.4}
\end{equation*}
$$

The maximal argument $\bar{y}$ is produced from $f$ and a choice of $\alpha<f$. We shall say that the couple $(f, \alpha)$ is $c$-allowable if the above condition (9.4) is satisfied.

Lemma 9.1. - Let $0<\alpha<f<1, \sigma_{\alpha}^{2}=\alpha(1-\alpha)$ and let us define $\tau>0$ by writing $f=\alpha+\sigma_{\alpha} \tau$. The function $v_{f, \alpha}$ reaches its maximum at

$$
\begin{equation*}
\bar{y}=\bar{y}_{f, \alpha}=\frac{\tau}{\sigma_{\alpha}}=\frac{f-\alpha}{\sigma_{\alpha}^{2}}>0 \tag{9.5}
\end{equation*}
$$

If $0<y, \bar{y} \leqslant 1 / 2$ then

$$
\begin{equation*}
\mathrm{e}^{(y-\bar{y})^{2} / 2} v_{f, \alpha}(y) \geqslant v_{f, \alpha}(\bar{y})=\left(\frac{f}{\alpha}\right)^{f}\left(\frac{1-f}{1-\alpha}\right)^{1-f} \tag{9.6}
\end{equation*}
$$

If $0<\bar{y} \leqslant 1 / 4$ and $\bar{y}^{4} \leqslant c / n$, then the couple $(f, \alpha)$ is $c$-allowable.
Proof. - Let $w(s)=\ln v_{f, \alpha}(s)=f \ln (1+\beta s)+(1-f) \ln (1-\alpha s)$. We have

$$
w^{\prime}(s)=\frac{\beta f}{1+\beta s}-\frac{\alpha(1-f)}{1-\alpha s}, \quad w^{\prime \prime}(s)=-\frac{\beta^{2} f}{(1+\beta s)^{2}}-\frac{\alpha^{2}(1-f)}{(1-\alpha s)^{2}}
$$

The maximal argument $\bar{y}$ is found by solving $w^{\prime}(\bar{y})=0$, yielding

$$
\bar{y}=\frac{f-\alpha}{\sigma_{\alpha}^{2}}, \quad 1+\beta \bar{y}=\frac{f}{\alpha}, \quad 1-\alpha \bar{y}=\frac{1-f}{1-\alpha} .
$$

This gives us the maximal value $v_{f, \alpha}(\bar{y})$ at the right-hand side of (9.6). Suppose now that we have $0<y, \bar{y} \leqslant 1 / 2$. Using Taylor-Lagrange at $\bar{y}$, we get

$$
w(y)-w(\bar{y})=w^{\prime \prime}(\xi) \frac{(y-\bar{y})^{2}}{2}
$$

for some $\xi$ between $y$ and $\bar{y}$, hence $0<\xi \leqslant 1 / 2$. We have $1-\alpha \xi \geqslant 1 / 2$ and

$$
-w^{\prime \prime}(\xi) \leqslant \beta^{2} f+4 \alpha^{2}(1-f) \leqslant \beta+4 \alpha^{2}(1-\alpha) \leqslant \beta+\alpha=1
$$

because $\alpha(1-\alpha) \leqslant 1 / 4$. This implies the left-hand side of (9.6).
Suppose that $0<\bar{y} \leqslant 1 / 4$. Moving around $\bar{y}$, we can find $y_{\mathbb{N}}>0$ satisfying

$$
\frac{\left|y_{\mathbb{N}}-\bar{y}\right|}{y_{\mathbb{N}} \bar{y}}=\left|\frac{1}{\bar{y}}-\frac{1}{y_{\mathbb{N}}}\right| \leqslant 1
$$

and such that $1 / y_{\mathbb{N}}-\alpha$ is an odd integer. From $\left|y_{\mathbb{N}}-\bar{y}\right| \leqslant y_{\mathbb{N}} \bar{y}$ and $\bar{y} \leqslant 1 / 4$ follows that $y_{\mathbb{N}} \leqslant 4 \bar{y} / 3 \leqslant 1 / 3<1 / 2$. Also, $\left|y_{\mathbb{N}}-\bar{y}\right| \leqslant 4 \bar{y}^{2} / 3<\sqrt{2} \bar{y}^{2}$. By (9.6), we deduce that $v_{f, \alpha}\left(y_{\mathbb{N}}\right) \geqslant \mathrm{e}^{-\bar{y}^{4}} v_{f, \alpha}(\bar{y})$ and the conclusion is reached.

Given $f$ and $\alpha$ such that $0<\alpha<f<1$, let us now examine the optimal value

$$
\begin{equation*}
E_{f, \alpha}:=v_{f, \alpha}(\bar{y})=(1+\beta \bar{y})^{f}(1-\alpha \bar{y})^{1-f}=\left(\frac{f}{\alpha}\right)^{f}\left(\frac{1-f}{1-\alpha}\right)^{1-f} \tag{9.7}
\end{equation*}
$$

Consider the function $\phi_{\alpha}$ defined on $(0,1)$ by

$$
\begin{equation*}
\phi_{\alpha}(s)=s \ln \left(\frac{s}{\alpha}\right)+(1-s) \ln \left(\frac{1-s}{1-\alpha}\right), \quad s \in(0,1) . \tag{9.8}
\end{equation*}
$$

We see that $\phi_{\alpha}^{\prime}(s)=\ln (s / \alpha)-\ln ((1-s) /(1-\alpha)), \phi_{\alpha}^{\prime \prime}(s)=1 / s+1 /(1-s)=$ $1 /(s(1-s))$, and $\phi_{\alpha}^{(3)}(s)=-s^{-2}+(1-s)^{-2}$. Note that $\phi_{\alpha}(\alpha)=\phi_{\alpha}^{\prime}(\alpha)=0$, and that $\phi_{\alpha}^{\prime \prime}(\alpha)=\sigma_{\alpha}^{-2}$.

Lemma 9.2. - If $0<\alpha<f=\alpha+\sigma_{\alpha} \tau<1$, the maximal value $v_{f, \alpha}(\bar{y})$ satisfies

$$
\begin{equation*}
\ln v_{f, \alpha}(\bar{y})=\phi_{\alpha}(f) \geqslant \frac{\tau^{2}}{2}-\frac{1-2 \alpha}{\sigma_{\alpha}} \frac{\tau^{3}}{6} . \tag{9.9}
\end{equation*}
$$

Proof. - By Taylor-Lagrange for $\phi_{\alpha}$ at the point $\alpha$, we have

$$
\phi_{\alpha}(f)=\phi_{\alpha}^{\prime \prime}(\alpha) \frac{(f-\alpha)^{2}}{2}+\phi_{\alpha}^{(3)}(\xi) \frac{(f-\alpha)^{3}}{6}=\frac{\tau^{2}}{2}+\phi_{\alpha}^{(3)}(\xi) \frac{\left(\sigma_{\alpha} \tau\right)^{3}}{6}
$$

for some $\xi \in(\alpha, f)$. Since $\phi_{\alpha}^{(3)}$ is increasing, we get that

$$
\phi_{\alpha}(f)-\frac{\tau^{2}}{2} \geqslant \phi_{\alpha}^{(3)}(\alpha) \frac{\left(\sigma_{\alpha} \tau\right)^{3}}{6}=\frac{2 \alpha-1}{\alpha^{2}(1-\alpha)^{2}} \frac{\sigma_{\alpha}^{3} \tau^{3}}{6}=\frac{2 \alpha-1}{\sigma_{\alpha}} \frac{\tau^{3}}{6}
$$

In all that follows, we see $\Omega_{1}:=[-N, N]$ as a probability space equipped with the uniform probability measure, denoted here by $P_{1}$, and we shall consider the cube $S_{N}^{n}=2 N Q_{n}$, equipped with the product measure $P=$ $P_{1}^{\otimes n}$, also the uniform probability measure, as being our main probability space $(\Omega, \mathcal{F}, P)$. On this space, the random variables $\left(\mathbf{1}_{C_{\alpha}}\left(x_{i}\right)\right)_{i=1}^{n}$, where $x=\left(x_{1}, \ldots, x_{n}\right) \in \Omega$, are independent and equal to 0 or 1 with respective probabilities $1-\alpha$ and $\alpha$. Their expectation is $\alpha$ and their variance is equal to $\sigma_{\alpha}^{2}=\alpha(1-\alpha) \leqslant 1 / 4$. For every $\alpha \in(0,1)$, we introduce the centered and variance 1 Bernoulli variable $X_{1, \alpha}$ defined on $\Omega_{1}$ by

$$
\begin{equation*}
X_{1, \alpha}=\frac{\mathbf{1}_{C_{\alpha}}-\alpha}{\sigma_{\alpha}}=\sqrt{\frac{1-\alpha}{\alpha}} \mathbf{1}_{C_{\alpha}}-\sqrt{\frac{\alpha}{1-\alpha}} \mathbf{1}_{\Omega_{1} \backslash C_{\alpha}} \tag{9.10}
\end{equation*}
$$

and we let

$$
X_{n, \alpha}(x)=\frac{\sum_{i=1}^{n} X_{1, \alpha}\left(x_{i}\right)}{\sqrt{n}}=\sum_{i=1}^{n} \frac{\mathbf{1}_{C_{\alpha}}\left(x_{i}\right)-\alpha}{\sigma_{\alpha} \sqrt{n}}, \quad x=\left(x_{1}, \ldots, x_{n}\right) \in \Omega .
$$

We also let $N_{n, \alpha}(x)=\sum_{i=1}^{n} \mathbf{1}_{C_{\alpha}}\left(x_{i}\right)$ denote the number of coordinates of $x$ that are in $C_{\alpha}$. We are ready for a first explicit estimate of the maximal function $\mathrm{M}_{Q} \mu_{N}^{(n)}$.

Lemma 9.3. - Let $0<\alpha<1$ and $\sigma_{\alpha}^{2}=\alpha(1-\alpha)$. Let $n \in \mathbb{N}^{*}, t>0$ and $0<\theta<1$ be such that $\sqrt{n} \geqslant 2 t \sigma_{\alpha}^{-2}(1-\theta)^{-1}$. We have $\mathrm{M}_{Q} \mu_{N}^{(n)}>\mathrm{e}^{\theta t^{2} / 2}$ on the set

$$
A_{\alpha, t}^{(n)}=\left\{x \in 2\left(N-t^{-1} \sqrt{n}\right) Q_{n}: N_{n, \alpha}(x)=\sum_{i=1}^{n} \mathbf{1}_{C_{\alpha}}\left(x_{i}\right)>\alpha n+t \sigma_{\alpha} \sqrt{n}\right\}
$$

where $C_{\alpha}$ is defined at (9.2). When the dimension $n$ is large, and assuming the size $N$ large enough compared to $n$, it follows that

$$
\frac{\left|\left\{\mathrm{M}_{Q} \mu_{N}^{(n)}>\mathrm{e}^{\theta t^{2} / 2}\right\}\right|}{\left|S_{N}^{n}\right|} \geqslant \frac{\left|A_{\alpha, t}^{(n)}\right|}{\left|2 N Q_{n}\right|}>\frac{1}{2} \gamma_{1}((t,+\infty))
$$

Proof. - By the central limit theorem (see [32] for instance), we know that the distribution of $X_{n, \alpha}$ tends to the distribution of a $N(0,1)$ Gaussian random variable $G$ when $n$ tends to infinity. This yields

$$
P\left(N_{n, \alpha}>\alpha n+t \sigma_{\alpha} \sqrt{n}\right)=P\left(X_{n, \alpha}>t\right) \underset{n}{\longrightarrow} P(G>t)=\gamma_{1}((t,+\infty))
$$

Let $A_{\alpha, t}^{(n, 0)}$ be the set of points $x \in \Omega$ where $N_{n, \alpha}(x)>\alpha n+t \sigma_{\alpha} \sqrt{n}$. Fix $x \in A_{\alpha, t}^{(n, 0)}$ and let $m=N_{n, \alpha}(x)$. We shall apply Lemma 9.1 with $f=m / n$ and $\tau=t / \sqrt{n}$. By assumption, the optimal argument $\bar{y}$ satisfies

$$
\bar{y}=\frac{t}{\sigma_{\alpha} \sqrt{n}} \leqslant \frac{\sigma_{\alpha}(1-\theta)}{2}<1 / 4
$$

At (9.3), we used a cube centered at $x$, with side length $\ell+\alpha, \ell$ an odd integer. We can choose $\ell+\alpha<1 / \bar{y}+2<2 / \bar{y}<t^{-1} \sqrt{n}$. This cube must be contained in $\Omega=S_{N}^{n}$, so we have to give up a small part of $A_{\alpha, t}^{(n, 0)}$, close to the boundary of $\Omega$. We thus introduce the subset $A_{\alpha, t}^{(n)}=A_{\alpha, t}^{(n, 0)} \cap 2\left(N-t^{-1} \sqrt{n}\right) Q_{n}$. The difference $A_{\alpha, t}^{(n, 0)} \backslash A_{\alpha, t}^{(n)}$ gets negligible when the side $2 N$ of $S_{N}^{n}$ tends to infinity since $\left(1-t^{-1} \sqrt{n} / N\right)^{n} \rightarrow_{N} 1$, so the set $A_{\alpha, t}^{(n)}$ has essentially the same probability as $A_{\alpha, t}^{(n, 0)}$ when $N=N(n)>\kappa(t) n^{3 / 2}$ is large enough. When $n$ tends to infinity, the probability of $A_{\alpha, t}^{(n)}$ is therefore, say, larger than $\gamma_{1}((t,+\infty)) / 2$.

We first show that the couple $(f, \alpha)$ is $c$-allowable with $c=(1-\theta) t^{2} / 4$. We know that $\bar{y}<1 / 4$ and on the other hand, we have

$$
\bar{y}^{4}=\frac{t^{4}}{\sigma_{\alpha}^{4} n^{2}}=\frac{c}{n} \frac{4 t^{2}}{(1-\theta) \sigma_{\alpha}^{4} n}<\frac{c}{n}\left(\frac{2 t}{(1-\theta) \sigma_{\alpha}^{2} \sqrt{n}}\right)^{2} \leqslant \frac{c}{n}
$$

It follows from Lemma 9.1 that $\mathrm{M}_{Q} \mu_{N}^{(n)}(x) \geqslant \mathrm{e}^{-(1-\theta) t^{2} / 4} V(\bar{y})$ for every $x \in$ $A_{\alpha, t}^{(n)}$. It remains to estimate the optimal value $V(\bar{y})$. For this we apply (9.9). It implies that $V(\bar{y}) \geqslant \mathrm{e}^{t^{2} / 2}$ when $\alpha \geqslant 1 / 2$, and when $\alpha \leqslant 1 / 2$, we see that

$$
\frac{1-2 \alpha}{\sigma_{\alpha}} \frac{\tau^{3}}{6}=\frac{(1-2 \alpha) \tau}{3 \sigma_{\alpha}} \frac{\tau^{2}}{2}<\frac{t}{3 \sigma_{\alpha} \sqrt{n}} \frac{\tau^{2}}{2}<\frac{\sigma_{\alpha}(1-\theta)}{6} \frac{\tau^{2}}{2}<\frac{(1-\theta) \tau^{2}}{4}
$$

so that $V(\bar{y}) \geqslant \mathrm{e}^{t^{2} / 2-(1-\theta) t^{2} / 4}$ and $\mathrm{M}_{Q} \mu_{N}^{(n)}(x) \geqslant \mathrm{e}^{t^{2} / 2} \mathrm{e}^{-(1-\theta) t^{2} / 2}=$ $\mathrm{e}^{\theta t^{2} / 2}$.

Given $\alpha \in(0,1)$, we have identified a subset $A_{\alpha, t}^{(n)}$ of $S_{N}^{n}$ where $\mathrm{M}_{Q} \mu_{N}^{(n)}$ is large. We shall have to use several values of $\alpha$, and show that the union of the corresponding sets provides a fair amount of the total volume of $S_{N}^{n}$. We thus introduce $0<\alpha_{0}<\alpha_{2}<\ldots<\alpha_{K}<1$ and we will prove that the probability of the union of sets $\left(A_{\alpha_{j}, t}^{(n)}\right)_{j=0}^{K}$ gets $>1 / 4$, say, when $K$ is large but fixed and when $n$ tends to infinity. Rather than relying, as Aubrun does, on the law of iterated logarithm, we apply easy facts behind the proof of that "law". In a simple qualitative approach, we shall analyze the Gaussian limit of the joint distribution of $\left(X_{n, \alpha_{j}}\right)_{j=0}^{K}$, which is the distribution of a Gaussian vector $\left(G_{j}\right)_{j=0}^{K}$ whose covariance matrix $C$ is the same as that of $\left(X_{n, \alpha_{j}}\right)_{j=0}^{K}$. Letting $\sigma_{j}^{2}=\alpha_{j}\left(1-\alpha_{j}\right)$, the entries of $C$ are $C_{j, k}=\mathrm{E}\left(X_{1, \alpha_{j}} X_{1, \alpha_{k}}\right)=\sigma_{j}^{-1} \sigma_{k}^{-1}\left(\alpha_{j} \wedge \alpha_{k}-2 \alpha_{j} \alpha_{k}+\alpha_{j} \alpha_{k}\right), \quad 0 \leqslant j, k \leqslant K$.
Note that $C_{j, j}=1$. Assuming $\alpha_{j} \leqslant \alpha_{k}$, that is to say, assuming $j \leqslant k$, we get

$$
C_{j, k}=\sigma_{j}^{-1} \sigma_{k}^{-1} \alpha_{j}\left(1-\alpha_{k}\right)=\sqrt{\frac{\alpha_{j}}{1-\alpha_{j}}} \sqrt{\frac{1-\alpha_{k}}{\alpha_{k}}} .
$$

We fix $v \in(0,1)$ and set $w=\sqrt{1-v^{2}}$. We define $\alpha_{j}=\left(1+v^{2 j}\right)^{-1}, j=$ $0, \ldots, K$, and obtain $C_{j, k}=v^{|k-j|}$. We can realize the distribution of $\left(G_{j}\right)_{j=0}^{K}$ by considering the larger Gaussian sequence indexed by $\mathbb{Z}$, which is defined by the sums of the series $G_{j}=w \sum_{i \leqslant j} v^{j-i} U_{i}$, for every $j \in \mathbb{Z}$, where the $\left(U_{i}\right)_{i \in \mathbb{Z}}$ are independent $N(0,1)$ Gaussian variables. Indeed, if $j \leqslant k$ we have that

$$
\mathrm{E}\left(G_{j} G_{k}\right)=\left(1-v^{2}\right) \sum_{i \leqslant j} v^{j+k-2 i}=v^{k-j}=v^{|k-j|}
$$

We see that $G_{j}-v G_{j-1}=w U_{j}$ and it follows that

$$
\begin{equation*}
\max _{1 \leqslant j \leqslant J}\left|U_{j}\right|=w^{-1} \max _{1 \leqslant j \leqslant J}\left|G_{j}-v G_{j-1}\right| \leqslant w^{-1}(1+v) \max _{0 \leqslant j \leqslant J}\left|G_{j}\right| \tag{9.11}
\end{equation*}
$$

We now recall an extremely classical estimate.
Lemma 9.4. - Let $J \geqslant 21$ be an integer and set

$$
s_{J}:=\sqrt{2 \ln J-\ln (16 \pi \ln J)} .
$$

If $U_{1}, \ldots, U_{J}$ are independent $N(0,1)$ Gaussian variables, one has that

$$
P\left(\max _{1 \leqslant j \leqslant J} U_{j}>s_{J}\right)>1 / 2
$$

Proof. - We have for $s>0$ that

$$
\begin{equation*}
\int_{s}^{+\infty} \mathrm{d} \gamma_{1}(s)>\frac{s}{\sqrt{2 \pi}\left(1+s^{2}\right)} \mathrm{e}^{-s^{2} / 2} \tag{9.12}
\end{equation*}
$$

consequence of

$$
\mathrm{e}^{-s^{2} / 2} / s=\int_{s}^{+\infty}\left(1+u^{-2}\right) \mathrm{e}^{-u^{2} / 2} \mathrm{~d} u<\left(1+s^{-2}\right) \int_{s}^{+\infty} \mathrm{e}^{-u^{2} / 2} \mathrm{~d} u
$$

When $J \geqslant 21$, one has $\mathrm{e}^{-1} J^{2}>16 \pi \ln J>1$, hence $1<s_{J}<\sqrt{2 \ln J}$. Therefore, we see by (9.12) for each $j=1, \ldots, J$ that

$$
P\left(U_{j}>s_{J}\right) \geqslant \frac{s_{J}}{\sqrt{2 \pi}\left(1+s_{J}^{2}\right)} \frac{\sqrt{16 \pi \ln J}}{J} \geqslant \frac{2 s_{J}^{2}}{\left(1+s_{J}^{2}\right) J} \geqslant \frac{1}{J} .
$$

It follows that

$$
P\left(\max _{1 \leqslant j \leqslant J} U_{j} \leqslant s_{J}\right) \leqslant\left(1-\frac{1}{J}\right)^{J}<\mathrm{e}^{-1}<\frac{1}{2}
$$

Theorem 9.5 (Aldaz [1]). - The weak type $(1,1)$ constant $\kappa_{Q, n}$ in (9.1) does not stay bounded when the dimension $n$ tends to infinity.

Proof. - Given an arbitrary $t>1$, we let $t_{1}:=t w^{-1}(1+v)>t$ and choose an integer $K \geqslant 21$ such that $s_{K}>t_{1}$. Applying Lemma 9.4, we obtain that the event $\left\{\max _{1 \leqslant j \leqslant K}\left|U_{j}\right|>t_{1}\right\}$ has probability $>1 / 2$, and by (9.11), it follows that the event $\left\{\max _{0 \leqslant j \leqslant K}\left|G_{j}\right|>t\right\}$ also has probability $>1 / 2$.

We see that $\sup \left|G_{j}\right|$ is the maximum of $\sup G_{j}$ and $\sup \left(-G_{j}\right)$ that have the same distribution, hence $\left\{\max _{0 \leqslant j \leqslant K} G_{j}>t\right\}$ has probability $>1 / 4$. Consequently, given any $t>1$, we obtain by the central limit theorem that the union of sets $A_{\alpha_{j}, t}^{(n)}$, for $j=0, \ldots, K$, has a probability close to that of $\left\{\max _{0 \leqslant j \leqslant K} G_{j}>t\right\}$, hence $>1 / 4$ when $n$ is large. By Lemma 9.3, given $\theta \in(0,1)$ and if $\sqrt{n}(1-\theta) \sigma_{K}^{2}>2 t$, the maximal function $\mathrm{M}_{Q} \mu_{N}^{(n)}$ is larger than $\mathrm{e}^{\theta t^{2} / 2}$ on the union $\bigcup_{j=0}^{K} A_{\alpha_{j}, t}^{(n)}$, i.e., on a subset of $\Omega=S_{N}^{n}$ having probability $>1 / 4$, hence $\kappa_{Q, n} \geqslant \mathrm{e}^{\theta t^{2} / 2} / 4$ when $n$ is large enough.

Aubrun [3] gives a lower bound $\kappa_{Q, n} \geqslant \kappa_{\varepsilon}(\ln n)^{1-\varepsilon}$ for every $\varepsilon>0$ by making quantitative the proof above. He applies to this end results proved years before (by Bretagnolle-Massart [14] in 1989 and previously, by Komlós-Major-Tusnády [51] in 1975) on the approximation of Brownian bridges, when $n \rightarrow+\infty$ and with explicit bounds, by binomial processes

$$
Z_{t}^{(n)}=\sum_{i=1}^{n} \frac{\mathbf{1}_{\left\{Y_{i} \leqslant t\right\}}-t}{\sqrt{n}}, \quad t \in[0,1]
$$

where the $\left(Y_{i}\right)_{i=1}^{n}$ are independent and uniform on $[0,1]$. One can see that the distribution of the process $\left(Z_{t}^{(n)}\right)_{t \in(0,1)}$ is equal to that of $\left(\sigma_{t} X_{n, t}\right)_{t \in(0,1)}$.

Iakovlev and Strömberg [46] begin with the same observations, in particular introducing the measure $\mu_{N}^{(n)}$, using the fundamental estimate (9.3) and, in a less apparent manner, the value $\mathrm{e}^{\theta t^{2} / 2}$ from Lemma 9.3. But instead of working in a probabilistic setting, they proceed to a finer combinatorial analysis. Contrary to Aubrun, they do not use values $\alpha$ close to 1 , nor close to 0 . In our exposition of their arguments, we shall work towards simplicity rather than optimality.

Let us digress a little with some comments on the Gaussian process viewpoint, and express in terms of stochastic maximal function the lower bound for $\mathrm{M}_{Q} \mu_{N}^{(n)}$ given in (9.3). Let $x \in \Omega$ and $m=N_{n, \alpha}(x), \sigma_{\alpha}^{2}=\alpha(1-\alpha)$ and write $m=\alpha n+\sigma_{\alpha} t \sqrt{n}$. Notice that $t=(m-\alpha n) /\left(\sigma_{\alpha} \sqrt{n}\right)=X_{n, \alpha}(x)$. We let $f$ be the fraction $m / n$, and rewrite the preceding formula for $m$ as $f=\alpha+\sigma_{\alpha} \tau$, with $\tau=t / \sqrt{n}$. We know the optimal argument $\bar{y}$ for $V(y)$, given in (9.5) by

$$
\bar{y}=\frac{t}{\sigma_{\alpha} \sqrt{n}}=\frac{\tau}{\sigma_{\alpha}}, \quad \text { and } \quad \frac{\ln V(\bar{y})}{n}=f \ln \left(\frac{f}{\alpha}\right)+(1-f) \ln \left(\frac{1-f}{1-\alpha}\right) .
$$

By Lemma 9.2 we have $\ln E_{f, \alpha}=\phi_{\alpha}(f) \geqslant \tau^{2} / 2=t^{2} /(2 n)$ if $\tau>0$ and $\alpha \geqslant 1 / 2$. Let $1 / 2 \leqslant \alpha \leqslant 3 / 4$ and assume that $0<t=X_{n, \alpha}(x) \leqslant n^{1 / 4} / 2$.

We see then that $\bar{y}=t /\left(\sigma_{\alpha} \sqrt{n}\right)<2 n^{-1 / 4} / \sqrt{3}$, thus $n \bar{y}^{4} \leqslant 16 / 9, \bar{y} \leqslant 1 / 4$ for $n>455$ and by Lemma 9.1 we are then in the allowable case with $c \leqslant 16 / 9$. This yields

$$
\begin{align*}
\mathrm{M}_{Q} \mu_{N}^{(n)}(x) \geqslant \kappa^{-1} E_{f, \alpha}^{n} \geqslant \kappa^{-1} \exp & \left(\frac{t^{2}}{2}\right) \\
& \text { with } \kappa<\mathrm{e}^{16 / 9}<6, n>455 . \tag{9.13}
\end{align*}
$$

Let us define a maximal function $X^{*}(x)=\sup _{1 / 2 \leqslant \alpha \leqslant 3 / 4} X_{n, \alpha}^{(1)}(x)$, where $X_{n, \alpha}^{(1)}(x)=X_{n, \alpha}(x)$ when $0 \leqslant 2 X_{n, \alpha}(x) \leqslant n^{1 / 4}$ and $X_{n, \alpha}^{(1)}(x)=0$ otherwise. We get

$$
6 \mathrm{M}_{Q} \mu_{N}^{(n)}(x) \geqslant \exp \left(\frac{X^{*}(x)^{2}}{2}\right)
$$

and the weak type $(1,1)$ constant $\kappa_{Q, n}$ must satisfy the condition

$$
P\left(\left\{X^{*}>s\right\}\right) \leqslant P\left(\left\{\mathrm{M}_{Q} \mu_{N}^{(n)}>\mathrm{e}^{s^{2} / 2} / 6\right\}\right) \leqslant 6 \kappa_{Q, n} \mathrm{e}^{-s^{2} / 2}, \quad s>0 .
$$

This explains how delicate the question can be. Indeed, given a subgaussian process $\left(Y_{t}\right)_{t \in T}$ satisfying tail estimates of the form $P(Y>s) \leqslant$ $\kappa \mathrm{e}^{-s^{2} /\left(2 d^{2}\right)}$ for every $s>0$, for each difference $Y=Y_{t_{2}}-Y_{t_{1}}$ and with $d=d\left(t_{1}, t_{2}\right)=\left\|Y_{t_{1}}-Y_{t_{2}}\right\|_{2}$, the well known chaining technique of Dudley [28] does not allow one to prove for the maximal process $\sup _{t \in T} Y_{t}$ a subgaussian inequality with the same bounding function $\mathrm{e}^{-s^{2} / 2}$, but rather with $\mathrm{e}^{-C s^{2} / 2}$ for some $C<1$, which is inoperative here.

Theorem 9.6 (Iakovlev and Strömberg [46]). - One has that

$$
\kappa_{Q, n} \geqslant \kappa n^{1 / 4} .
$$

Rather than exploiting the exponential asymptotics (9.13) of $E_{f, \alpha}^{n}$, we shall observe some more nice features of the expression $E_{f, \alpha}$ defined in (9.7), where $f=m / n=\alpha+t \sigma_{\alpha} / \sqrt{n}=\alpha+\sigma_{\alpha} \tau$. We replace the value $\mathrm{e}^{\theta t^{2} / 2}$ seen in Lemma 9.3 by a fixed large value $V>1$ and we try to keep the (conditional on allowability) lower bound $E_{f, \alpha}^{n}$ for $\mathrm{M}_{Q} \mu_{N}^{(n)}$ constantly equal to $V$. Equivalently, we keep

$$
\begin{equation*}
E_{f, \alpha}=\mathrm{e}^{\phi_{\alpha}(f)}=V^{1 / n}>1 \tag{9.14}
\end{equation*}
$$

for all values of $f$ (or of $m$ ) that will be handled. The possibility of finding $\alpha$ satisfying (9.14) comes from the fact that for every given $f \in(0,1)$, the function

$$
\begin{equation*}
\psi_{f}: s \mapsto\left(\frac{f}{s}\right)^{f}\left(\frac{1-f}{1-s}\right)^{1-f}=\mathrm{e}^{\phi_{s}(f)}, \quad s \in(0,1) \tag{9.15}
\end{equation*}
$$

is convex on $(0,1)$ (actually, log-convex), tends to infinity at 0 and at 1 , and assumes its minimal value $\psi_{f}(f)=1$ at $s=f$. Consequently, there
are exactly two values $\alpha_{0}<f<\alpha_{1}$ of $\alpha \in(0,1)$ solving (9.14), we shall consider the smallest one and set $\alpha(f)=\alpha_{0}$. Notice that $\left(\ln \psi_{f}\right)^{\prime}(s)=$ $-f / s+(1-f) /(1-s)$ vanishes at $s=f$, and

$$
\begin{equation*}
\left(\ln \psi_{f}\right)^{\prime \prime}(s)=\frac{f}{s^{2}}+\frac{1-f}{(1-s)^{2}}>f+(1-f)=1 \tag{9.16}
\end{equation*}
$$

We have therefore for every $s \in(0,1)$ that

$$
\begin{align*}
& \ln \psi_{f}(s) \geqslant(s-f)^{2} / 2 \\
& \quad \text { thus }(f-\alpha(f))^{2} / 2 \leqslant \ln \psi_{f}(\alpha(f))=(\ln V) / n \tag{9.17}
\end{align*}
$$

From now on, we fix two values $0<f_{*}<f^{*} \leqslant 1 / 2$, independent of the dimension $n$. For every integer $m$ in the range $\left[f_{*} n, f^{*} n\right.$ ], we shall consider the set

$$
F_{m}=\left\{x \in \Omega: N_{n, \alpha(f)}(x)=m\right\}, \quad \text { with } f=m / n
$$

Let us write $\alpha=\alpha(f)$ for brevity. We have that $E_{f, \alpha}=V^{1 / n}$ and if we assume $c$-allowability for $(f, \alpha)$ we get $\mathrm{M}_{Q} \mu_{N}^{(n)}(x) \geqslant \mathrm{e}^{-c} V$ for every $x \in F_{m}$, by (9.4). The probability of $F_{m}$ is $\alpha^{m}(1-\alpha)^{n-m}\binom{n}{m}$ and we see that

$$
\begin{aligned}
V P\left(F_{m}\right) & =\left(\frac{f}{\alpha}\right)^{m}\left(\frac{1-f}{1-\alpha}\right)^{n-m} \alpha^{m}(1-\alpha)^{n-m}\binom{n}{m} \\
& =\frac{m^{m}(n-m)^{n-m}}{n^{n}}\binom{n}{m} .
\end{aligned}
$$

Stirling's formula in the form $\mathrm{e}^{-1 /(12 p)} p!\leqslant p^{p} \mathrm{e}^{-p} \sqrt{2 \pi p} \leqslant p!$ (see [66]) gives

$$
\begin{equation*}
\mathrm{e}^{-1 /(12 n)} V P\left(F_{m}\right) \leqslant \frac{\sqrt{n}}{\sqrt{2 \pi m(n-m)}} \leqslant \mathrm{e}^{n /(12 m(n-m))} V P\left(F_{m}\right) \tag{9.18}
\end{equation*}
$$

With $s_{*}=\sqrt{f_{*}\left(1-f_{*}\right)}$ and $s^{*}=\sqrt{f^{*}\left(1-f^{*}\right)}$, it follows that

$$
\begin{equation*}
V P\left(F_{m}\right) \geqslant \frac{\mathrm{e}^{-1 /(12 f(1-f) n)}}{\sqrt{2 \pi f(1-f) n}} \geqslant \frac{\mathrm{e}^{-1 /\left(12 s_{*}^{2} n\right)}}{s^{*} \sqrt{2 \pi}} \frac{1}{\sqrt{n}} . \tag{9.19}
\end{equation*}
$$

If the sets $F_{m}$ were disjoint (and the couples $(f, \alpha(f)) c$-allowable) we would get immediately, by summing on $m$ between $f_{*} n$ and $f^{*} n$, a lower bound of

$$
\kappa_{Q, n} \geqslant \mathrm{e}^{-c} V P\left(\left\{\mathrm{M}_{Q} \mu_{N}^{(n)} \geqslant \mathrm{e}^{-c} V\right\}\right) \text { by } \kappa\left[\mathrm{e}^{-c}\left(f^{*}-f_{*}\right) /\left(s^{*} \sqrt{2 \pi}\right)\right] \sqrt{n},
$$

but this disjointness property is clearly not true. We shall specify a suitable large $V$ such that the probability of the intersection of two events $F_{m_{1}}$ and $F_{m_{2}}$ will be small compared to the probability of $F_{m_{1}}$, when $m_{1}<m_{2}$ are not too close. We shall find a subset $M \subset\left[f_{*} n, f^{*} n\right]$, as large as possible,
consisting of "well spaced" values $m_{j}$ giving rise to $c$-allowable couples. The final estimate has the form

$$
\begin{equation*}
\kappa_{Q, n} \geqslant \mathrm{e}^{-c} V P\left(\bigcup_{m \in M} F_{m}\right) \tag{9.20}
\end{equation*}
$$

where the probability of the union will be larger than half of the sum of probabilities. The seemingly harmless allowability restriction that $y^{-1}-\alpha=$ $\sigma_{\alpha} / \tau-\alpha$ must be an odd integer $\ell$ will actually cause a heavy loss at the end.

We fix $\varepsilon \in\left(0, f_{*}\right]$ and introduce $\eta:=\sqrt{1-\varepsilon / f_{*}}$. We define the "big" value $V$ as $V=\mathrm{e}^{\varepsilon^{2} n / 2}$. By (9.17), we have that

$$
\begin{equation*}
0<f-\alpha(f) \leqslant \varepsilon \tag{9.21}
\end{equation*}
$$

Lemma 9.7. - Suppose that $0<\alpha<\xi \leqslant f \leqslant \alpha+\varepsilon$ and $f_{*} \leqslant f \leqslant 1 / 2$. One has

$$
\begin{equation*}
\eta^{2} \leqslant \frac{\alpha}{\xi}<\frac{\alpha(1-\alpha)}{\xi(1-\xi)}<1, \quad \text { in particular } \eta \sigma_{f}=\eta \sqrt{f(1-f)}<\sigma_{\alpha} \tag{9.22}
\end{equation*}
$$

Assuming $V=\mathrm{e}^{\varepsilon^{2} n / 2}, \alpha=\alpha(f)$ and writing $\sigma_{\alpha} \tau=f-\alpha$, one has that

$$
\begin{equation*}
\eta \tau \leqslant \varepsilon \leqslant \tau \tag{9.23}
\end{equation*}
$$

Proof. - We see that $\alpha(1-\alpha)<\xi(1-\xi)$ because $0<\alpha<\xi \leqslant 1 / 2$. Next, we get

$$
\frac{\alpha(1-\alpha)}{\xi(1-\xi)}>\frac{\alpha}{\xi} \geqslant \frac{f-\varepsilon}{f} \geqslant 1-\frac{\varepsilon}{f_{*}}=\eta^{2}
$$

By Taylor-Lagrange at $\alpha$ for the function $\phi_{\alpha}$ defined in (9.8), we have

$$
\phi_{\alpha}(f)=\phi_{\alpha}^{\prime \prime}\left(\xi_{0}\right) \frac{(f-\alpha)^{2}}{2}=\frac{\sigma_{\alpha}^{2}}{\xi_{0}\left(1-\xi_{0}\right)} \frac{\tau^{2}}{2}=\frac{\alpha(1-\alpha)}{\xi_{0}\left(1-\xi_{0}\right)} \frac{\tau^{2}}{2}
$$

for some $\xi_{0} \in(\alpha, f)$, and $\phi_{\alpha}(f)=\phi_{\alpha(f)}(f)=(\ln V) / n=\varepsilon^{2} / 2$ by assumption. The inequalities in (9.23) follow then from (9.21) and (9.22).

We have to understand how the values $\alpha(f)$ are distributed when $f$ varies in $\left[f_{*}, f^{*}\right]$. To this end, we estimate the derivative $\alpha^{\prime}(f)$.

Lemma 9.8. - Let $0<\varepsilon \leqslant f_{*}$ and $V=\mathrm{e}^{\varepsilon^{2} n / 2}$. The mapping $(0,1) \ni$ $f \mapsto \alpha(f)$ implicitly defined at (9.14) is increasing, and when $f \in\left[f_{*}, f^{*}\right]$ we have that

$$
\eta^{2}<\alpha^{\prime}(f)<1
$$

Proof. - We express the derivative $\alpha^{\prime}(f)$ by differentiating with respect to $f$ the equality $\phi_{\alpha(f)}(f)=(\ln V) / n$. Writing $\phi_{\alpha}$ for $\phi_{\alpha(f)}$, we obtain

$$
\phi_{\alpha}^{\prime}(f)+\left(\frac{\partial}{\partial \alpha} \phi_{\alpha}(f)\right) \alpha^{\prime}(f)=\phi_{\alpha}^{\prime}(f)-\frac{f-\alpha}{\alpha(1-\alpha)} \alpha^{\prime}(f)=0 .
$$

By Taylor-Lagrange at $\alpha$ for $s \mapsto \phi_{\alpha}^{\prime}(s)$, there is $\xi \in(\alpha, f)$ such that

$$
\phi_{\alpha}^{\prime \prime}(\xi)(f-\alpha)=\phi_{\alpha}^{\prime}(f)=\frac{f-\alpha}{\alpha(1-\alpha)} \alpha^{\prime}(f), \quad \text { hence } \quad \alpha^{\prime}(f)=\frac{\alpha(1-\alpha)}{\xi(1-\xi)}>0
$$

because $\phi_{\alpha}^{\prime \prime}(\xi)=\sigma_{\xi}^{-2}$. We have that $\alpha<\xi<f \leqslant \alpha+\varepsilon$ by (9.21), and when we further assume $f_{*} \leqslant f \leqslant f^{*} \leqslant 1 / 2$ the conclusion follows by (9.22).

We need to study the intersections $F_{m_{1}} \cap F_{m_{2}}$, when $m_{1}, m_{2} \in\left[f_{*} n, f^{*} n\right]$.
Lemma 9.9. - Suppose that $f_{*} n \leqslant m_{1}<m_{2} \leqslant f^{*} n$. One has that

$$
\mathrm{e}^{-1 /(6 n)} P\left(F_{m_{1}} \cap F_{m_{2}}\right) / P\left(F_{m_{1}}\right)<\lambda \mathrm{e}^{-\delta^{2} \varepsilon^{2}\left(m_{2}-m_{1}\right) / 2} / \sqrt{2 \pi\left(m_{2}-m_{1}\right)},
$$

with $\delta=\eta^{3} s_{*} /\left(1-f_{*}\right)$ and $\lambda=\sqrt{1-f_{*}} / \sqrt{1-f^{*}}$.
Proof. - Let $f_{j}=m_{j} / n, f_{*} \leqslant f_{j} \leqslant f^{*}$, and $\alpha_{j}=\alpha\left(f_{j}\right)$, for $j=1,2$. By Lemma 9.8, we have that $\alpha_{1}<\alpha_{2}$ since $f_{1}<f_{2}$. Let $J$ be an arbitrary subset of $\{1, \ldots, n\}$ satisfying $|J|=m_{1}$, and let $A(J)$ be the subset of $\Omega=S_{N}^{n}$ defined by

$$
A=A(J)=\left\{x=\left(x_{1}, \ldots, x_{n}\right) \in \Omega: J=\left\{i: x_{i} \in C_{\alpha_{1}}\right\}\right\} .
$$

One has thus $N_{n, \alpha_{1}}(x)=m_{1}$ when $x \in A$. The conditional probability $p_{A}$ that $N_{n, \alpha_{2}}(x)=m_{2}$ knowing that $x \in A$ is equal to the probability that $m_{A}:=m_{2}-m_{1}$ of the remaining $n_{A}:=n-m_{1}=\left(1-f_{1}\right) n \geqslant n / 2$ coordinates of $x$ (those coordinates that are in $\Omega_{1} \backslash C_{\alpha_{1}}$ ) fall in $C_{\alpha_{2}} \backslash C_{\alpha_{1}}$. This is given by the binomial distribution corresponding to $n_{A}$ and to $\alpha_{A}=\left(\alpha_{2}-\alpha_{1}\right) /\left(1-\alpha_{1}\right)$, and we know therefore that

$$
\begin{aligned}
p_{A}:=\frac{P\left(\left\{N_{n, \alpha_{2}}=m_{2}\right\} \cap A\right)}{P(A)} & =P\left(\left\{N_{n_{A}, \alpha_{A}}=m_{A}\right\}\right) \\
& =\alpha_{A}^{m_{A}}\left(1-\alpha_{A}\right)^{n_{A}-m_{A}}\binom{n_{A}}{m_{A}} .
\end{aligned}
$$

Let $f_{A}=m_{A} / n_{A}=\left(f_{2}-f_{1}\right) /\left(1-f_{1}\right)$. Since $\alpha^{\prime}(f)<1$ on $\left[f_{*}, f^{*}\right]$, we get

$$
\begin{aligned}
f_{A} & =\frac{f_{2}-f_{1}}{1-f_{1}}=\frac{f_{2}-f_{1}}{1-\alpha_{1}}\left(1+\frac{f_{1}-\alpha_{1}}{1-f_{1}}\right) \\
& >\frac{\alpha_{2}-\alpha_{1}}{1-\alpha_{1}}+\frac{\left(f_{1}-\alpha_{1}\right)\left(f_{2}-f_{1}\right)}{\left(1-\alpha_{1}\right)\left(1-f_{1}\right)}=\alpha_{A}+\frac{f_{1}-\alpha_{1}}{\left(1-\alpha_{1}\right)\left(1-f_{1}\right)}\left(f_{2}-f_{1}\right)
\end{aligned}
$$

Let $f_{1}-\alpha_{1}=\sigma_{\alpha_{1}} \tau_{1}$. We have $\tau_{1} \geqslant \varepsilon$ by (9.23), $\sigma_{f_{1}}>\sigma_{\alpha_{1}}>\eta \sigma_{f_{1}} \geqslant \eta s_{*}$ by (9.21) and (9.22), and $f_{*} \leqslant f_{1} \leqslant 1 / 2$. By the leftmost inequality in (9.22), we obtain

$$
\frac{1}{1-\alpha_{1}}>\frac{\alpha_{1}}{f_{1}\left(1-f_{1}\right)} \geqslant \frac{\eta^{2}}{1-f_{1}} \geqslant \frac{\eta^{2}}{1-f_{*}}
$$

therefore

$$
\begin{equation*}
f_{A}-\alpha_{A}>\frac{\eta^{3} s_{*} \varepsilon}{\left(1-f_{*}\right)\left(1-f_{1}\right)}\left(f_{2}-f_{1}\right)=\frac{\delta \varepsilon}{1-f_{1}}\left(f_{2}-f_{1}\right) . \tag{9.24}
\end{equation*}
$$

Recalling the function $\psi_{f}$ from (9.15), we see that

$$
p_{A}=\psi_{f_{A}}\left(\alpha_{A}\right)^{-n_{A}} f_{A}^{m_{A}}\left(1-f_{A}\right)^{n_{A}-m_{A}}\binom{n_{A}}{m_{A}}
$$

Applying Stirling as before in (9.18), and because we have that $n_{A} /\left(n_{A}-\right.$ $\left.m_{A}\right)=\left(1-f_{1}\right) /\left(1-f_{2}\right) \leqslant\left(1-f_{*}\right) /\left(1-f^{*}\right)$, we obtain

$$
\begin{aligned}
\mathrm{e}^{-1 /\left(12 n_{A}\right)} p_{A} & <\psi_{f_{A}}\left(\alpha_{A}\right)^{-n_{A}} \sqrt{\frac{n_{A}}{2 \pi m_{A}\left(n_{A}-m_{A}\right)}} \\
& \leqslant \psi_{f_{A}}\left(\alpha_{A}\right)^{-n_{A}} \sqrt{\frac{1-f_{*}}{2 \pi\left(1-f^{*}\right) m_{A}}}
\end{aligned}
$$

For some $\xi \in\left(\alpha_{A}, f_{A}\right)$, and since $\left(\ln \psi_{f_{A}}\right)^{\prime \prime}(\xi)>f_{A} / \xi^{2}>1 / f_{A}$ by (9.16), we get

$$
\ln \psi_{f_{A}}\left(\alpha_{A}\right)=\left(\ln \psi_{f_{A}}\right)^{\prime \prime}(\xi) \frac{\left(f_{A}-\alpha_{A}\right)^{2}}{2}>\frac{\left(f_{A}-\alpha_{A}\right)^{2}}{2 f_{A}}
$$

Consequently, we can write

$$
p_{A}<\mathrm{e}^{1 /\left(12 n_{A}\right)} \exp \left(-n_{A} \frac{\left(f_{A}-\alpha_{A}\right)^{2}}{2 f_{A}}\right) \frac{\lambda}{\sqrt{2 \pi m_{A}}}, \quad \text { with } \lambda=\frac{\sqrt{1-f_{*}}}{\sqrt{1-f^{*}}}
$$

We see that $n_{A} / f_{A}=n_{A}^{2} /\left(m_{2}-m_{1}\right)$. By (9.24) we have

$$
\frac{n_{A}}{f_{A}}\left(f_{A}-\alpha_{A}\right)^{2} \geqslant \frac{n^{2}\left(1-f_{1}\right)^{2}}{m_{2}-m_{1}} \frac{\delta^{2} \varepsilon^{2}\left(f_{2}-f_{1}\right)^{2}}{\left(1-f_{1}\right)^{2}}=\delta^{2} \varepsilon^{2}\left(m_{2}-m_{1}\right)
$$

Using also $n<2 n_{A}$ and the definition of $p_{A}$, we obtain for $A=A(J)$ that

$$
\begin{aligned}
& P\left(A(J) \cap\left\{N_{n, \alpha_{2}}=m_{2}\right\}\right) \\
& \quad<\left(\mathrm{e}^{1 /(6 n)} \lambda \mathrm{e}^{-\delta^{2} \varepsilon^{2}\left(m_{2}-m_{1}\right) / 2} / \sqrt{2 \pi\left(m_{2}-m_{1}\right)}\right) P(A(J)) .
\end{aligned}
$$

Summing on all subsets $J$ of $\{1, \ldots, n\}$ with $|J|=m_{1}$, and because $\bigcup_{|J|=m_{1}} A(J)$ is equal to $\left\{N_{n, \alpha_{1}}=m_{1}\right\}=F_{m_{1}}$, we get

$$
P\left(F_{m_{1}} \cap F_{m_{2}}\right)<\left(\mathrm{e}^{1 /(6 n)} \lambda \mathrm{e}^{-\delta^{2} \varepsilon^{2}\left(m_{2}-m_{1}\right) / 2} / \sqrt{2 \pi\left(m_{2}-m_{1}\right)}\right) P\left(F_{m_{1}}\right)
$$

End of proof of Theorem 9.6. - Let $H$ be a sufficiently large integer, and let us now define $M=\{j H: j \in \mathbb{N}\} \cap\left[f_{*} n, f^{*} n\right]$ to be the set of multiples of $H$ located in the segment $\left[f_{*} n, f^{*} n\right]$. We fix $m_{1} \in M$ and let $m_{2}>m_{1}$ be any other element of $M$. Then $m_{2}-m_{1}=k H$ with $k$ integer $\geqslant 1$. Summing on $k \geqslant 1$ we see that

$$
\sum_{k=1}^{+\infty} \frac{\mathrm{e}^{-\delta^{2} \varepsilon^{2} k H / 2}}{\sqrt{k H}}<\int_{0}^{+\infty} \mathrm{e}^{-\delta^{2} \varepsilon^{2} H s / 2} \frac{\mathrm{~d} s}{\sqrt{H s}}=\frac{\sqrt{2} \Gamma(1 / 2)}{\varepsilon H \delta}=\frac{\sqrt{2 \pi}}{\varepsilon H \delta}
$$

By Lemma 9.9 , we get $\sum_{m_{2} \in M, m_{2}>m_{1}} P\left(F_{m_{1}} \cap F_{m_{2}}\right)<P\left(F_{m_{1}}\right) / 2$ when $\varepsilon H$ is larger than $2 \lambda \mathrm{e}^{1 /(6 n)} / \delta$. It follows then that at least one half of the set $F_{m_{1}}$ is not covered by the other sets $F_{m_{2}}$ for $m_{2}>m_{1}$ and $m_{2} \in M$, therefore $P\left(\bigcup_{m \in M, m \geqslant m_{1}} F_{m}\right) \geqslant P\left(\bigcup_{m \in M, m>m_{1}} F_{m}\right)+P\left(F_{m_{1}}\right) / 2$ for $m_{1} \in$ $M$. The probability of $\bigcup_{m \in M} F_{m}$ is thus at least equal to half of the sum of probabilities. By (9.19) and (9.20) we get

$$
\begin{equation*}
\kappa_{Q, n} \geqslant \mathrm{e}^{-c} V P\left(\bigcup_{m \in M} F_{m}\right) \geqslant \frac{\mathrm{e}^{-c}}{2} \sum_{m \in M} V P\left(F_{m}\right) \geqslant \frac{\mathrm{e}^{-c}}{2} \frac{\mathrm{e}^{-1 /\left(12 s_{*}^{2} n\right)}}{s^{*} \sqrt{2 \pi}} \frac{|M|}{\sqrt{n}} \tag{9.25}
\end{equation*}
$$

So far we could hope for a lower bound of order $\sqrt{n}$ for the weak type constant. But we have to comply with the allowability restriction, and we must estimate the number of couples $(f, \alpha(f))$ that are $c$-allowable. We let

$$
\varepsilon=\frac{s_{*} n^{-1 / 4}}{1+s_{*} n^{-1 / 4} / f_{*}}, \quad \text { so that } \frac{\varepsilon}{\eta^{2}}=\frac{\varepsilon}{1-\varepsilon / f_{*}}=s_{*} n^{-1 / 4}
$$

and $\varepsilon<f_{*}$. We choose a spacing $H \sim n^{1 / 4}$. For every $m \in M$, for $f=m / n$, $\alpha=\alpha(f)$ and $f=\alpha+\sigma_{\alpha} \tau$ we have by (9.5), (9.22) and (9.23) that

$$
\bar{y}=\frac{\tau}{\sigma_{\alpha}} \leqslant \frac{\varepsilon}{\eta} \frac{1}{\eta \sigma_{f}} \leqslant \frac{\varepsilon}{\eta^{2} s_{*}}=n^{-1 / 4} .
$$

For $n>256$ we see that $\bar{y}<1 / 4$ and $\bar{y}^{4}<1 / n$, thus $(f, \alpha(f))$ is allowable with constant $c=1$ according to Lemma 9.1. We choose the spacing integer $H$ such that $H>2 \lambda \mathrm{e}^{1 /(6 n)} /(\delta \varepsilon)$. Since $\varepsilon=\eta^{2} s_{*} n^{-1 / 4}$, we arrive to the condition

$$
H>\left(2 \lambda / \delta \eta^{2} s_{*}\right) \mathrm{e}^{1 /(6 n)} n^{1 / 4}
$$

We obtain a set $M \subset\left[f_{*} n, f^{*} n\right]$ of multiples of $H$ with cardinality at least equal to $\left\lfloor\left(f^{*} n-f_{*} n\right) / H\right\rfloor>\left[\eta^{2} \delta s_{*}\left(f^{*}-f_{*}\right) /(2 \lambda)\right] \mathrm{e}^{-1 /(6 n)} n^{3 / 4}-1$ where each element $m$ produces a 1-allowable couple $(f, \alpha(f))$. By (9.25), we get that

$$
\kappa_{Q, n} \geqslant \frac{1}{2 \mathrm{e}} \frac{\mathrm{e}^{-1 /\left(12 s_{*}^{2} n\right)}}{s^{*} \sqrt{2 \pi}} \frac{|M|}{\sqrt{n}} \geqslant \frac{\eta^{2} \delta s_{*}\left(f^{*}-f_{*}\right)}{4 \mathrm{e} \sqrt{2 \pi} \lambda s^{*}} n^{1 / 4}-O\left(n^{-1 / 2}\right) .
$$

Our version of the Iakovlev-Strömberg proof is not optimal, we shall however try to figure out a numerical value for the constant that we get in front of $n^{1 / 4}$. We have for $n$ large that $\varepsilon=o(1)$, thus $\eta \simeq 1$. Let us introduce

$$
z:=\frac{\delta s_{*}\left(f^{*}-f_{*}\right)}{\eta^{3} \lambda s^{*}}=\frac{s_{*}^{2}}{1-f_{*}} \frac{f^{*}-f_{*}}{\sqrt{f^{*}\left(1-f_{*}\right)}}=\frac{f_{*}\left(f^{*}-f_{*}\right)}{\sqrt{f^{*}\left(1-f_{*}\right)}}
$$

This expression increases with $f^{*}$, so we set $f^{*}=1 / 2$, the maximal possibility. Then, the resulting value of $z$ is maximal for $f_{*}=3 / 4-\sqrt{11 / 48} \sim 0.271$,
yielding $z>0.102$. When $n$ is large, we have

$$
\kappa_{Q, n}>\frac{z}{4 \mathrm{e} \sqrt{2 \pi}} n^{1 / 4}-o\left(n^{1 / 4}\right)>0.0037 n^{1 / 4}>\frac{n^{1 / 4}}{271}
$$

Notice that we have set the constant value $V$ as $V=V_{n} \sim \mathrm{e}^{\kappa \sqrt{n}}$ in dimension $n$. The corresponding sequence of values $t_{n}=\sqrt{2 \ln V_{n}} \sim n^{1 / 4}$ for the "test sets" $\left\{X_{n, \alpha}>t_{n}\right\}$ is "invisible" to the Gaussian limit argument of Theorem 9.5.

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| stability under translation (of a function on $\mathbb{R}$ ) . . . . . . . . . . . . . . . . . 8.3, p. 164 |  |  |  |
| stochastic integral . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . 1.4.2, p. 28 |  |  |  |
| stopping time . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . 1.4.3, p. 30 |  |  |  |
| strong type inequality . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . Intro., p. 3 |  |  |  |
| subordination principle . . . . . . . . . . . . . . . . . . . . . . . . . bf. eq. (1.30), p. 34 |  |  |  |
| Thorin's method (for interpolation) . . . . . . . . . . . . . . . . . . . . . . . . . 3.2.2, p. 60 |  |  |  |
| three lines lemma . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . Lem. 3.1, p. 51 |  |  |  |
| transition matrix (of a Markov chain) . . . . . . . . . . . . . . . . . . . . . . . . 1.3, p. 17 |  |  |  |
| uncentered maximal function |  |  |  |
| unconditionality (of martingale differences) . . . . . . . . . . . . . . . Rem. 1.8, p. 29 |  |  |  |
| variance . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . 1.4.1, p. 22 |  |  |  |
| Vitali covering lemma . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . Intro., p. 5 weak type inequality . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . Intro., p. 3 |  |  |  |
|  |  |  |  |
| $\sigma$-field . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . 1 , p. 8 |  |  |  |
| $\sigma$-finite measure . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . . Rem. 1.2, p. 11 |  |  |  |
| -stable (- function on $\mathbb{R}$ ) . . . . . . . . . . . . . . . . . . . . . . . . . . . . . 8.3, p. 164 |  |  |  |

## Notation

| $\mathbf{1}_{A}$ | indicator function of the set $A$. . . . . . . . . . . . . . . . . . 1, p. 9 |
| :---: | :---: |
| $A_{p}, B_{p}$ | constants in Khinchin's inequalities . . . . . . . . . . eq. (1.22.K), p. 26 |
| $\left(A_{\mathbf{0}}\right) \ldots\left(A_{\mathbf{3}}\right)$ | assumptions for Carbery's Proposition 6.6 . . . . . . bf. Prop. 6.6, p. 96 |
| $A_{\alpha, t}^{(n)}$ | set where the maximal function $\mathrm{M}_{Q} \mu_{N}^{(n)}$ is large . . . . Lem. 9.3, p. 180 |
| $a \wedge b, a \vee b$ | minimum, maximum of two real numbers $a$ and $b$ |
| $\mathcal{B}_{Y}$ | Borel $\sigma$-field of a topological space $Y$ |
| $\left(B_{t}\right)_{t \geqslant 0}$ | Brownian motion in $\mathbb{R}^{n}$. . . . . . . . . . . . . . . . . . . . 1.4.1, p. 22 |
| $B_{\tau}$ | Brownian value $\omega \mapsto B_{\tau(\omega)}(\omega)$ at a stopping time $\tau \quad \ldots . .$. 1.4.3, p. 31 |
| $B\left(q_{0}, R, n\right)$ | a priori bound in Bourgain's cube proof . . . . . . . . . . . . 8.1.1, p. 154 |
| bar $\mu$ | barycenter of a probability measure $\mu$ on $\mathbb{R}^{n}$. . . . . . . . 1.4.1, p. 21 |
| $C_{p}, C_{p}^{\prime}, C_{p}^{\prime \prime}$ | constants for Carbery's Proposition 6.6 . . . . . eq. $\left(A_{\mathbf{0}}\right)-\left(A_{\mathbf{2}}\right)$, p. 96 |
| $C_{\alpha}(m)$ | Carbery's constant for a Fourier multiplier m . . . . . Prop. 6.14, p. 111 |
| $C_{\alpha}$ | a subset of $[-N, N]$ in proof of Aldaz-Aubrun weak type theorem |
|  |  |
| $D^{z} h, D^{\alpha} h$ | fractional derivative of $h$. . . . . . . eq. (6.7), p. 101 \& (6.10), p. 102 |
| $D_{t}^{\alpha} h(\lambda t)$ | fractional derivative of $t \mapsto h(\lambda t) \ldots \ldots$ |
| $\left.D_{t}^{\alpha} h(\lambda t)\right\|_{t}$ | fractional derivative evaluated at $t_{0}$. . . . . . . . af. eq. (6.8), p. 101 |
| $\left(d_{k}\right)_{k=0}^{N}$ | martingale difference sequence . . . . . . . . . . . . . . . . . 1.4.2, p. 25 |
| $d^{z} h, d^{\alpha} h$ | another fractional derivative of $h$. . . . . . . . . . . . eq. (7.9), p. 131 |
| $\mathrm{E} f$ | expectation of the random variable $f$. . . . . . . . . . . . . 1, p. 9 |
| $\underset{\mathrm{E}}{(f \mid \mathcal{G})}$ | conditional expectation of $f$ on the $\sigma$-field $\mathcal{G}$. . . . . . . . 1, p. 9 |
| $\left(\mathbf{e}_{j}\right)_{j=1}^{n}$ | standard unit vector basis of $\mathbb{R}^{n}$ |
| $\mathcal{F} f, \widehat{f}, \widehat{\mu}$ | Fourier transform of a function $f$, of a measure $\mu$. . . . . . . 2, p. 35 |
| $F_{m}$ | set of $x \in \mathbb{R}^{n}$ with $N_{n, \alpha(f)}(x)=m$ in Section 9 . . . af. eq. (9.17), p. 185 |
| $\\|f\\|_{p},\\|f\\|_{L^{p}}$ | norm of a function $f$ in $L^{p}$. . . . . . . . . . . . . . . . Th. 0.1, p. 3 |
| $f^{*}$ | uncentered classical maximal function of $f$. . . . . . . . . Intro., p. 3 |
| $f_{\#} \mu$ | pushforward image of the finite measure $\mu$ by the mapping $f$. . . . p. 9 |
| $f_{r}^{*}, f_{\ell}^{*}$ | right, left maximal function of a function $f$ on $\mathbb{R}$. . . . . eq. (5.5), p. 79 |
| $f_{*}, f^{*}$ | lower, upper bound for $f=m / n$, proof of Iakovlev-Strömberg |
|  | af. eq. (9.17), p. 185 |
| $G_{s} f$ | Gaussian semi-group acting on $f$. . . . . . . . . . . eq. (1.19), p. 23 |
| $G$ | Gaussian kernel in Bourgain's cube proof, $\widehat{G}(\xi)=\mathrm{e}^{-4 \pi \mid \xi}$ |
|  | af. eq. (A), p. 151 |
|  | absolute $p$-th moment of the Gaussian measure $\gamma_{1}$ on $\mathbb{R}$. eq. (1.18), p. 22 |
|  | inverse Fourier transform of a function $g$ on $\mathbb{R}^{n} \ldots . . . . . . . .2, ~ p . ~ 36 ~$ |
| $g(f), g_{k}(f)$ | Littlewood-Paley functions for $f$ on $\mathbb{R}^{n}, k \geqslant 1$. . . . . . . 2.1, p. 36 |
| $g_{(\lambda)}, g_{[\lambda]}$ | dilates of a function $g$ on $\mathbb{R}^{n}$. . . . . . . . . . . . . . . eq. (2.7), p. 40 |
| $H f, H_{\mathbb{T}} f$ | Hilbert transform of $f$ on $\mathbb{R}$ or $\mathbb{T}$. . . . . . . . . . . . . . . . 2.3, p. 44 |
| $H_{k}$ | homogeneous parts in Pisier's semi-group theorem . . bf. eq. (8.8), p. 155 |
| $\\|h\\|_{L_{\alpha}^{2}}$ | Carbery's multiplier norm of a function $h$ on $(0,+\infty)$. . eq. (6.21), p. 109 |
| $h_{q}$ | $L^{q}$ bound for $H_{1}$ in Pisier's semi-group theorem . . . . Prop. 8.3, p. 155 |
| $I$ | identity operator on $L^{q}\left(\mathbb{R}^{n}\right)$. . . . . . . . . . . . . . . . 8.2, p. 155 |
| $\mathrm{I}_{n}$ | identity matrix of size $n \times n$ |
| $I^{w} f, I^{\beta} f$ | fractional integration . . . . . . . . . . . . . . . . . . . eq. (6.9), p. 102 |
| $i^{w} f, i^{\beta} f$ | another fractional integration . . . . . . . . . . . . . . . . . 7.1, p. 131 |
|  | Bessel function of order $\nu$ |
| $\mathcal{K}\left(\mathbb{R}^{n}\right)$ | space of compactly supported continuous functions on $\mathbb{R}^{n}$ |
| $K_{C}$ | uniform probability density on a convex set $C$. . . . . . . . 5.1, p. 74 |
| $K_{l c}$ | a symmetric log-concave probability density on $\mathbb{R}^{n} \ldots$. Prop. 5.10, p. 86 |
| $K_{g}$ | a probability density or kernel on $\mathbb{R}^{n}$ satisfying (6.1.H) |
|  | bf. eq. (6.1.H), p. 93 |
| $K^{R}$ | Bourgain's kernels for the cube . . . . . . . . . . . . . . . . . . . 8, p. 150 |
| $L^{p}\left(\mathbb{R}^{n}\right)$ | Lebesgue spaces, $1 \leqslant p \leqslant+\infty$ |
| $L(C)$ | isotropy constant of the convex set C . . . . . . . . . . . . eq. (5.2), p. 76 |
| $\mathrm{M} f$ | classical Hardy-Littlewood maximal function of $f$. . . eq. (0.1), p. 2 |
| $\mathrm{M}_{C} f$ | maximal function of $f$ associated to $C$. . . . . . . . . eq. (0.3.M), p. 5 |


| $M_{N}^{*}$ | maximal function of a martingale ( $\left.M_{k}\right)_{k=0}^{N}$. . . . . . . . . 1.1, p. 10 |
| :---: | :---: |
| $\mathcal{M} f$ | radial maximal function of $f$. . . . . . . . . . . . . . . . . 4.1, p. 63 |
| $\mathrm{M}_{K} f, M_{K} f$ | maximal function of $f$ associated to a kernel $K$. . 3.3, p. $62 \& 5.1$, p. 75 |
| $\mathrm{M}_{C}^{(d)} f$ | dyadic maximal function of $f$ associated to C . . . . . . . . . . . 6, p. 92 |
| $\\|m\\|_{p \rightarrow p}$ | norm on $L^{p}$ of a multiplier m . . . . . . . . . . . . . . . . 2.2, p. 41 |
| $m_{\sigma}$ | Fourier transform of the uniform probability measure $\sigma$ on $S^{n-1}$ |
|  |  |
| $m^{*}$ | $m^{*}(\xi)=\xi \cdot \nabla m(\xi)$, for a multiplier $m$ on $\mathbb{R}^{n}$. . . . . . . . . . . p. 69 |
| $m_{C}, m_{C}(\xi)$ | Fourier transform of $K_{C}$. . . . . . . . . . . . . . . . . . . . . 5.1, p. 74 |
| $m_{l c}, m_{g}$ | Fourier transform of $K_{l c}$, of $K_{g}$ |
| $m_{z}^{\varepsilon}$ | Müller's holomorphic family of multipliers . . . . . . . . eq. (7.12), p. 132 |
| $m^{\#}$ | Müller's "crucial" multiplier $m^{\#}(\xi)=\|\xi\| m(\xi)$. . . . . . . eq. (7.19), p. 142 |
| $m^{R}$ | Bourgain's cube multiplier, Fourier transform of $K^{R}$. . . . . . 8, p. 150 |
| $\mathbb{N}^{*}$ | set of integers $n>0$ |
| $N\left(0, \mathrm{I}_{n}\right)$ | centered Gaussian distribution with covariance matrix $\mathrm{I}_{n}$. . . 1.4.1, p. 22 |
| $N_{n, \alpha}(x)$ | number of coordinates of $x \in \mathbb{R}^{n}$ that are in $C_{\alpha}$. . bf. Lem. 9.3, p. 180 |
| $\mathcal{O}(n)$ | orthogonal group |
| $P_{t}, P_{t} f$ | Poisson measure, Poisson semi-group acting on $f$. . . . . . . 1.5, p. 33 |
| $P_{t}^{(n)}$ | Poisson kernel on $\mathbb{R}^{n}$. . . . . . . . . . . . . . . . . . . eq. (1.32), p. 35 |
| $p^{*}$ | $p^{*}:=\max (p, p /(p-1))$, in Burkholder's unconditional constant $p^{*}-1$ |
| $Q_{\mu}$ | covariance quadratic form of a measure $\mu$. . . . . . . . . . . . 1.4 .1, p. 21 |
| $Q(C)$ | Müller's constant for a convex set C . . . . . . . . . . . . . . 7, p. 127 |
| $Q, Q_{n}$ | symmetric cube of volume one in $\mathbb{R}^{n}$. . . . . . . . . . . . . . . 8, p. 149 |
| $\mathrm{q}_{p}$ | Littlewood-Paley constant, for $1<p<+\infty$. . . . . . . eq. (2.4), p. 37 |
| $q(C)$ | modified Müller's constant for a symmetric convex set $C$. eq. (7.1), p. 127 |
| $R_{j} f$ |  |
| $\mathcal{R} f$ | vector form of the Riesz transform in $\mathbb{R}^{n}$. . . . . . . . . . . 2.3, p. 46 |
| $R_{0}$ | $R_{0}^{\delta / 2}=4$, in Bourgain's cube proof . . . . . . . . . . . . . . . . p. 159 |
| $\|S\|$ | measure of a set $S$ |
| $\|S\|_{n}$ | Lebesgue's $n$-dimensional measure of a set $S$ in $\mathbb{R}^{n}$ |
| $S_{N}$ | square function of a martingale . . . . . . . . . . . . . . . . . . 1.4.2, p. 25 |
| $\mathcal{S}\left(\mathbb{R}^{n}\right)$ | Schwartz function space on $\mathbb{R}^{n}$. . . . . . . . . . . . . . . . . 1.5, p. 33 |
| $S^{n-}$ | unit sphere in $\mathbb{R}^{n}$ |
| $S_{p}(\tau)$ | one-side moments of a log-concave probability density on $[\tau,+\infty)$ |
| $S_{N}, S_{N}^{n}$ | sets $[-N, N],[-N, N]^{n}$ in Section 9 . . . . . . . . . . af. eq. (9.1), p. 176 |
| $s_{n-}$ | Lebesgue measure of the unit sphere in $\mathbb{R}^{n}$. . . . . . . eq. (1.34), p. 35 |
|  | positive part of a real number $s, s^{+}=\max (s, 0)$ |
| $\lfloor s\rfloor,\lceil s\rceil$ | floor, ceiling of s, integers such that $s-1 \leqslant\lfloor s\rfloor \leqslant s \leqslant\lceil s\rceil<s+1$ |
| $s_{*}, s^{*}$ | $s_{*}=\sqrt{f_{*}\left(1-f_{*}\right)}, s^{*}=\sqrt{f^{*}\left(1-f^{*}\right)}$, in Iakovlev-Strömberg |
| $\mathbb{T}$ | $\cdots$ unit circle in $\mathbb{R}^{2}$ or $\mathbb{C}$. . . . . . . . . . . af. eq. (9.18), p. 185 |
| $\\|T\\|_{p \rightarrow p}$ | norm of an operator $T: L^{p} \rightarrow L^{p}$ |
| $T_{m}$ | linear operator associated to the multiplier $m$ on $\mathbb{R}^{n}$. . bf. eq. (2.9), p. 40 |
| $T_{j, v}, T_{j}, T^{*}$ | operators for Carbery's maximal theorem . . . . . . . . . . . 6.1, p. 94 |
| $\mathbf{T}^{J}$ | product $\prod_{j \in J} T_{j}$ of linear operators $\left(T_{j}\right)_{j \in J} \ldots . . . . . .88 .2, ~ p . ~ 155 ~$ |
| $t C$ | dilate by $t>0$ of the convex set $C$ |
| $U_{K}$ | operator $f \mapsto \nabla K * f$. . . . . . . . . . . . . . . . . eq. (7.20), p. 144 |
| $u(x, t)$ | harmonic extension of $f(x), x \in \mathbb{R}^{n}$, to the upper half-space in $\mathbb{R}^{n+1}$ |
| $V(K)$ |  |
| V | fixed large value $V$ in Iakovlev-Strömberg . . . . . . . af. Th. 9.6, p. 184 |
| $w_{0}$ | $w_{0}=R^{\delta / 2}$ in Bourgain's cube proof . . . . . . . . . af. Lem. 8.4, p. 157 |
| $w_{1}$ | $w_{1}=w_{0}^{2}=R^{\delta}$ in Bourgain's cube proof . . . . . . . af. eq. (8.32), p. 170 |
| $X_{1, \alpha}, X_{n, \alpha}$ | Bernoulli, binomial variable in Aldaz-Aubrun |
| $\begin{aligned} & \|x\| \\ & \\|x\\|_{C} \end{aligned}$ | . . . . . . . . . . . . . . . . . . . . . . . . eq. (9.10) \& bf. Lem. 9.3, p. 180 <br> norm of a vector $x$, usually Euclidean norm on $\mathbb{R}^{n}$ <br> norm of a vector $x$ relative to a symmetric convex set $C$. . . . 1.4.3, p. 32 |


| $x_{\Gamma}$ | point where $\Gamma$ reaches its minimum on ( $0, \infty$ ) . . . . . . eq. (3.7), p. 50 |
| :---: | :---: |
| $\bar{y}$ | maximal argument $\bar{y}=\tau / \sigma_{\alpha}$. . . . . . . . . . . . Lem. 9.1, p. 178 |
| $\alpha_{j}(m)$ | constituent of Bourgain's constant $\Gamma_{B}(K)$ for a kernel $K$ Lem. 5.14, p. 89 |
| $\alpha(f)$ | value associated to $f=m / n$ in Iakovlev-Strömberg . af. eq. (9.15), p. 185 |
| $\beta_{a}$ | in a bound for $\|\Gamma(z)\|^{-1}$. . . . . . . . . . . . . . . . . . eq. (3.12. $\Gamma$ ), p. 50 |
| $\beta_{j}(m)$ | constituent of Bourgain's constant $\Gamma_{B}(K)$ for a kernel $K$ Lem. 5.14, p. 89 |
| $\Gamma_{B}(K)$ | Bourgain's constant for a kernel K . . . . . . . . . . . . Lem. 5.14, p. 89 |
| $\Gamma^{S}$ | operator $\mathbf{T}^{\sim S}(\mathbf{I}-\mathbf{T})^{S}$, in Bourgain's cube proof . . . . eq. (8.11), p. 157 |
| $\gamma_{n}, \gamma_{F}$ | Gaussian probability measure on $\mathbb{R}^{n}$, on a Euclidean space $F$ eq. (1.17), p. 22 |
| $\Delta_{k}, \Delta_{k, c}$ |  |
| $\delta_{x}$ | Dirac probability measure at the point $x$. . . . . . . . . . . . . p. 22 |
| $\delta_{j, c}$ | bounds for $m_{C}, m_{l c}$ and their derivatives . . . . . . Lem. 5.11, p. 88 |
| $\delta_{j, g}$ | bounds for $m_{g}$ and its derivatives |
|  | . . . . . . . . . . eq. (6.1.H), p. 94 \& eq. (7.5.H ${ }_{\infty}$ ), p. 128 |
| $\delta$ | $\delta=\eta^{3} s_{*} /\left(1-f_{*}\right)$, in proof of Theorem 9.6 . . . . . . Lem. 9.9, p. 187 |
| $\partial S$ | boundary of a set $S \subset \mathbb{R}^{n}$ |
| $\partial_{i} f$ | $i$ th partial derivative of a function $f$ on $\mathbb{R}^{n}, i=1, \ldots, n$ |
| $\nabla f$ | gradient of the function $f$ on $\mathbb{R}^{n}$ |
| $\left(\varepsilon_{j}\right)_{j=1}^{+\infty}$ | independent Bernoulli random variables . . . . . . . bf. eq. (1.19), p. 23 |
| $\varepsilon$ | value $\varepsilon \in\left(0, f_{*}\right]$, in the proof of Theorem 9.6 . . . . bf. eq. (9.21), p. 186 |
| $\eta$ | $\eta=\sqrt{1-\varepsilon / f_{*}}$, in the proof of Theorem 9.6 . . . bf. eq. (9.21), p. 186 |
| $\theta^{\perp}$ | hyperplane orthogonal to $\theta \in S^{n-1}$ |
| $\kappa_{Q, n}$ | weak type ( 1,1 ) constant for the cube in $\mathbb{R}^{n}$. . . Intro., p. 6 \& 9, p. 176 |
| $\lambda_{p}$ | bound on $L^{p}$ for Laplace-type multipliers . . . . . . . . Prop. 2.2, p. 42 |
| $\widehat{\mu}$ | Fourier transform of the measure $\mu$. . . . . . . . . . . . . . . . 2, p. 35 |
| $\mu_{C}$ | uniform probability measure on a convex set C . . . . . . . . 5.1, p. 74 |
| $\mu^{+}, \mu^{-},\|\mu\|$ | positive, negative and absolute value of a measure . . . bf. Lem. 5.8, p. 84 |
| $\\|\mu\\|_{1}$ | mass of a real-valued measure $\mu$. . . . . . . . . . . . bf. Lem. 5.8, p. 84 |
| $\mu^{R}$ | Bourgain's cube measure . . . . . . . . . . . . . . . af. eq. (8.2), p. 151 |
| $\mu_{N}, \mu_{N}^{(n)}$ | discrete measure in the proof of Aldaz-Aubrun weak type theorem |
| $(\xi \cdot \nabla)^{\alpha}$ | Carbery's fractional operator . . . . . . . . . . . . . . eq. $\left(6.18 . \nabla^{\alpha}\right.$ ), p. 108 |
| $\rho_{p}$ | "collective" norm on $L^{p}$ for Riesz transforms . . . . . . eq. (2.22), p. 45 |
| $\Sigma_{k}$ | subsets of $\{1, \ldots, n\}$ having cardinality $k$. . . . . bf. eq. (8.8), p. 155 |
| $\sigma^{2}$ | variance of a probability measure or density . . . . . . . . . 1.4.1, p. 22 |
| $\sigma, \sigma_{n-1}$ | invariant probability measure on the unit sphere $S^{n-1}$ |
| $\sigma_{S, j}$ |  |
| $\sigma_{\alpha}^{2}$ | variance $\sigma_{\alpha}^{2}=\alpha(1-\alpha)$ of a Bernoulli random variable . . Lem. 9.1, p. 178 |
| $\Phi_{j}$ | convolution operator in Bourgain's cube proof . . . . bf. eq. (8.27), p. 165 |
| $\varphi_{\theta}, \varphi_{\theta, K}$ | marginal of a kernel $K$ on $\mathbb{R}^{n}$, image by $\theta \in S^{n-1} \ldots .$. eq. (2.14), p. 42 |
| $\varphi_{q}$ | lower bound of angle for Pisier's analytic semi-group theorem . . . . . . . |
| $\phi_{\alpha}$ |  |
| $\Omega_{1}, \Omega$ | sets $[-N, N],[-N, N]^{n}$ in Section 9 . . . . . . . . bf. eq. (9.10), p. 180 |
| $\omega_{n}$ | Lebesgue volume of the unit ball in $\mathbb{R}^{n}$. . . . . . . . eq. (1.34), p. 35 |


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