# SHIFTED MOMENT PROBLEM FOR GAUSSIAN MEASURES IN SOME ORLICZ SPACES 

## BY

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Abstract. Suppose that two Gaussian measures $\mu_{1}$ and $\mu_{2}$ on Orlicz space $\left(L_{M}(T, F, m),\| \|_{M}\right)$ fulfill the condition

$$
\begin{equation*}
\int\|x+y\|_{M}\left(\mu_{1}-\mu_{2}\right)(d y)=0 \tag{*}
\end{equation*}
$$

for each $x$ from $L_{M}$.
It is proved that, under some assumptions on modular $M$, measure $m$ and $q$, condition $\left(^{*}\right.$ ) implies $\mu_{1}=\mu_{2}$.

Introduction. Suppose that $\mu_{1}$ and $\mu_{2}$ are two probability measures on a separable Banach space ( $E,\|\cdot\|$ ) satisfying the condition

$$
\begin{equation*}
\int\|x+y\|^{q} \mu_{1}(d y)=\int\|x+y\|^{q} \mu_{2}(d y) \tag{+}
\end{equation*}
$$

for every $x \in E$ and some fixed $q>0$.
Linde [7] has recently proved that if $E=L^{r}$ and $q \neq k r$ for positive integers $k$, then $\mu_{1}=\mu_{2}$ (cf. also [4]). The purpose of this paper is to prove a similar theorem for some Orlicz spaces $L_{M}$, under the additional assumption that $\mu_{i}$ are Gaussian.

The technique employed here differs somewhat from that used in [7]; in particular, we do not rely on Hoffman-Jørgensen's result [5] and, even in the case of $L^{r}[0,1]$ spaces, our theorem is a little stronger, namely, $\mu_{1}=\mu_{2}$ if we only assume that $q \neq r$.

Preliminaries. We recall here briefly basic notions concerning Orlicz spaces and, after restricting our attention to the class of spaces we will work with, we state and prove some simple facts implied by our axioms.

Throughout the paper ( $T, F, m$ ) will stand for a finite separable measure space. By $M$ we denote a fixed Young function, that is a convex, strictly increasing function such that $M(t)=0$ iff $t=0$ and $M(-t)=M(t)$. We
assume further that $M$ satisfies $\left(\Delta_{2}\right)$ condition,

$$
M(2 t) \leqslant M(t) k \quad \text { for } t \geqslant t_{0}
$$

for some $k>0$ and $t_{0} \geqslant 0$. Observe that convexity of $M$ implies $k \geqslant 2$.
Now, by $L_{M}$ we denote the space of all real-valued measurable functions $f$ such that $\int M(f) d m<\infty$.

It is easy to see that $L_{M}$ is a linear space. Moreover, if

$$
\|f\|_{M}=\inf \left\{u>0 ; \int M(f / u) d m \leqslant 1\right\}, \quad f \in L_{M},
$$

then $\|\cdot\|_{M}$ is a norm on $L_{M}$ and $\left(L_{M},\|\cdot\|_{M}\right)$ becomes a (separable) Banach space [6].

Now, we restrict our attention to a rather special class of $L_{M}$-spaces which are a generalization of $L^{r}$-spaces ( $r \geqslant 2$ ).

First of all, we impose some smoothness properties on $M$. Namely, we assume that
(I) $M \in C^{2}$, that is, $M$ has continuous derivatives up to the order two.

For the sake of convenience we write: $p=d M / d t$ and $p^{\prime}=d^{2} M / d t^{2}$.
Next condition means that $p^{\prime}$, roughly speaking, cannot decrease "too fast":
(II) For all $t, s$ such that $0 \leqslant t \leqslant s$, we have $p^{\prime}(t) \leqslant p^{\prime}(s)+A$ for some $A \geqslant 0$.

The third condition,
(III) $M(t) \geqslant B t^{2}$ for $t \geqslant t_{1}$ and a $B>0$, says that $L_{M} \subset L^{2}$ and that the natural embedding of $L_{M}$ into $L^{2}$ is continuous. The $L^{2}$-norm and the corresponding inner product on $L_{M}$ will be denoted by $\|\cdot\|_{2}$ and $\langle\cdot, \cdot\rangle$, respectively.

In the sequel we assume that our $L_{M}$-space satisfies properties (I)-(III). We now draw some simple conclusions from these properties.

Proposition 1. Under the assumptions as above the following properties hold:

$$
\begin{equation*}
M(t) \leqslant t p(t) \leqslant M(2 t) \quad \text { for } t \geqslant 0 \tag{i}
\end{equation*}
$$

(ii)

$$
t p^{\prime}(t) \leqslant p(2 t)+A t \quad \text { for } t \geqslant 0
$$

(iii)

$$
t^{2} p^{\prime}(t) \leqslant\left(k^{2} / 2+A / B\right) M(t) \quad \text { for } t \geqslant \max \left(t_{0}, t\right)
$$

$$
\begin{equation*}
\sup _{\|f\|_{M} \leqslant c} \int M(f) d m<\infty \tag{iv}
\end{equation*}
$$

for every positive constant c.

Proof. (i) By convexity of $M$ the function $p$ is nondecreasing, so

$$
M(t)=\int_{0}^{t} p(s) d s \geqslant \frac{t}{2} p\left(\frac{t}{2}\right)
$$

whereas $M(t) / t \leqslant p(t)$, again by convexity of $M$.
(ii) Since $p^{\prime} \geqslant 0$, this, together with (II), yields

$$
p(t)=\int_{0}^{t} p^{\prime}(s) d s \geqslant \min _{t / 2 \leqslant s \leqslant t} p^{\prime}(s)(t / 2) \geqslant(t / 2)\left(p^{\prime}(t / 2)-A\right) .
$$

(iii) This follows easily by (i), (ii), and ( $\Delta_{2}$ )-condition.
(iv) Observe that the $\left(\Delta_{2}\right)$-condition implies, for $t \geqslant t_{0}$ and all $s \geqslant 0$, $M(s t) \leqslant K(s) M(t)$, where $K(s)=k(s \vee 1)^{\log _{2} k}$. Since $k \geqslant 2, K(s)$ is non-decreasing, so

$$
\begin{aligned}
\int M(f) d m & =\int_{\left.v /\|f\|_{M} \leqslant t_{0}\right\}} M(f) d m+\int_{\left.\forall f\|f\|_{M}>t_{0}\right\}} M(f) d m \\
& \leqslant M\left(t_{0}\|f\|_{M}\right) m(T)+K\left(\|f\|_{M}\right) \int M\left(f /\|f\|_{M}\right) d m \\
& =M\left(t_{0}\|f\|_{M}\right) m(T)+K\left(\|f\|_{M}\right),
\end{aligned}
$$

because $\int M\left(f /\|f\|_{M}\right) d m=1$.
Now, for fixed $u, h \in L_{M}$, let $G(t)=\|u+t h\|_{M}, t \in[-1,1]$. Further, let $u_{t}$ $=u+t h$ and $w_{t}=u_{t} /\left\|u_{t}\right\|_{M}$. The next proposition contains basic facts concerning smoothness of $G$ as well as some estimates essential in the sequel.

Proposition 2. Let $u \notin \operatorname{lin}\{h\}$. Then $G \in C^{2}$ and $\left|G^{\prime}(t)\right| \leqslant C\|h\|_{M}$, $\left|G^{\prime \prime}(t)\right| \leqslant C\|h\|_{M}^{2} /\left\|u_{t}\right\|_{M}$, where $C$ does not depend on $h, u$, $t$. Moreover,

$$
\begin{equation*}
G^{\prime}(t)=I_{1} / I_{2}, \tag{i}
\end{equation*}
$$

$$
\begin{equation*}
G^{\prime \prime}(t)=\frac{I_{3} I_{2}-I_{1} I_{4}}{I_{2}^{2}} \tag{ii}
\end{equation*}
$$

where

$$
\begin{aligned}
& I_{1}=\int p\left(w_{t}\right) h d m, \quad I_{2}=\int p\left(w_{t}\right) w_{t} d m, \quad I_{3}=\int p^{\prime}\left(w_{t}\right) w_{t}^{\prime} h d m \\
& I_{4}=\int\left(p^{\prime}\left(w_{t}\right) w_{t}+p\left(w_{t}\right)\right) w_{t}^{\prime} d m \quad \text { with } w_{t}^{\prime}=\frac{h\left\|u_{t}\right\|_{M}-u_{t} G^{\prime}(t)}{\|u\|_{M}^{2}}
\end{aligned}
$$

Proof. The existence of $G^{\prime}$ is standard and easily follows from our assumptions [6]. Formula (i) can be obtained by differentiating the equation $\int M\left(w_{t}\right) d m=1$.

We show that $\left|G^{\prime}(t)\right| \leqslant C\|h\|_{M}$. To do this assume that $\|h\|_{M} \leqslant 1$. Then
(i) and (iv) of Proposition 1 yield that

$$
\begin{aligned}
&\left|I_{1}\right|=\left|\int p\left(w_{t}\right) h d m\right| \leqslant \int p\left(\left|w_{t}\right|+|h|\right)\left(\left|w_{t}\right|+|h|\right) d m \\
& \leqslant \int M\left(2\left|w_{t}\right|+2|h|\right) d m \leqslant \sup _{\|f\|_{M} \leqslant 4} \int M(f) d m<\infty
\end{aligned}
$$

so $\left|I_{1}\right| \leqslant C\|h\|_{M}$.
To complete this part of the proof, observe that property (i) of Proposition 1 gives

$$
I_{2}=\int p\left(w_{t}\right) w_{t} d m=\int p\left(\left|w_{t}\right|\right)\left(\left|w_{t}\right|\right) d m \geqslant \int M\left(w_{t}\right) d m \geqslant 1
$$

- Next, observe that the formula for $G^{\prime \prime}$ is obtained by formal differentiation of expressions under the integral sign in $I_{1}$ and $I_{2}$, respectively. To justify this procedure it is enough, by virtue of the Mean Value Theorem, to show that
(a)

$$
\int \sup _{-1 \leqslant t \leqslant 1}\left|p^{\prime}\left(w_{t}\right) w_{t}^{\prime} h\right| d m<\infty
$$

and
(b)

$$
\int \sup _{-1 \leqslant t \leqslant 1}\left(\left|p^{\prime}\left(w_{t}\right) w_{t}^{\prime} w_{t}\right|+\left|p\left(w_{t}\right) w_{t}^{\prime}\right|\right) d m<\infty
$$

This can be done in a similar way as in estimating $I_{1}$ and is left to the reader. The above arguments also show that $G^{\prime \prime}$ is continuous.

Finally, to show that $\left|G^{\prime \prime}(t)\right| \leqslant C\|h\|_{M}^{2} /\left\|u_{t}\right\|_{M}$, we have to estimate $I_{2}, I_{3}$ and $I_{4}$. It is easy to see that

$$
\left|I_{2}\right| \leqslant \int M\left(2 w_{t}\right) d m \leqslant \sup _{\|f\|_{M} \leqslant 2} \int M(f) d m<\infty
$$

Next, using the formula for $w_{t}^{\prime}$ and the fact that $\left|G^{\prime}(t)\right| \leqslant C\|h\|_{M}$, writing $z_{t}=|h|+\left|w_{t}\right|, C_{1}=C \vee 1$, we get, for $\|h\|_{M} \leqslant 1$,

$$
\begin{aligned}
\left|I_{3}\right| & \leqslant \frac{C_{1}}{\left\|u_{t}\right\|_{M}} \int p^{\prime}\left(\left|w_{t}\right|\right)\left(|h|+\left|w_{t}\right|\right)|h| \cdot d m \leqslant \frac{C_{1}}{\left\|u_{t}\right\|_{M}} \int\left(A+p^{\prime}\left(z_{t}\right)\right) z_{t}^{2} d m \\
& \leqslant \frac{C_{1}}{\left\|u_{t}\right\|_{M}}\left[A t_{1}^{2} m(T)+\left(2 A / B+k^{2} / 2\right) \sup _{\|f\|_{M} \leqslant 2} \int M(f) d m\right] \leqslant \frac{C_{2}}{\left\|u_{t}\right\|_{M}}
\end{aligned}
$$

for a constant $C_{2}$ independent from $u, h, t$. Similarly, we get $\left|I_{4}\right| \leqslant C_{3}\|h\|_{M} /\left\|u_{t}\right\|_{M}$, which completes the proof.

Corollary 1. Let $u$ be a function on $T$ with values in $\{-1,0,1\}$ such that $\|u\|_{M} \neq 0$ and let $q$ be a fixed positive number. Write $S_{u}=\operatorname{supp} u$. Then

$$
\begin{equation*}
\frac{d G^{q}}{d t}(0)=q \frac{\|u\|_{M}^{q}}{m\left(S_{u}\right)}\langle u, h\rangle \tag{i}
\end{equation*}
$$

(ii) $\frac{d^{2} G^{q}}{d t^{2}}(0)=q\|u\|_{M}^{q-1}\left[\Phi_{1}(u)\left\|1_{s_{u}^{c}} h\right\|_{2}^{2}+\Phi_{2}(u)\left\|1_{s_{u}} h\right\|_{2}^{2}+\dot{\Phi}_{3}(u)\langle u, h\rangle^{2}\right]$,
where

$$
\begin{gathered}
\Phi_{1}(u)=\frac{p^{\prime}(0)}{p\left(1 /\|u\|_{M}\right) m\left(S_{u}\right)}, \quad \Phi_{2}(u)=\frac{p^{\prime}\left(1 /\|u\|_{M}\right)}{p\left(1 /\|u\|_{M}\right) m\left(S_{u}\right)} \\
\Phi_{3}(u)=\frac{\|u\|_{M}}{m^{2}\left(S_{u}\right)}\left[(q-1)-\frac{p^{\prime}\left(1 /\|u\|_{M}\right)}{\|u\|_{M} p\left(1 /\|u\|_{M}\right)}\right]
\end{gathered}
$$

and $G^{q}(t)=\|u+t h\|_{M}^{q}$.
The main result. We begin with one technical lemma.
Lemma 1. Let $(Y, F, v)$ be a finite measure space and let $h$ be a real-valued function defined on $(-1,1) \times Y$ with the properties:
(i)

$$
t \rightarrow h(t, y) \text { belongs to } C^{1} v \text {-a.e., }
$$

(ii)

$$
y \rightarrow h(t, y) \text { belongs to } L^{2}(v) \text { for all } t \in(-1,1)
$$

(iii) $\sup _{-1<t<1} \int\left(h^{\prime}(t, y)\right)^{2} v(d y)<\infty$, where $h^{\prime}=d h / d t$.

Put $g(t)=\int h(t, y) v(d y)$.
Then for every $\varepsilon, 0<\varepsilon<1$, we have

$$
\frac{d g}{d t}(t)=\int h^{\prime}(t, y) v(d y) \quad \text { for } t \in(-1+\varepsilon, 1-\varepsilon)
$$

Proof. Let

$$
A_{n}=\left\{y \in Y ; \sup _{-1+\varepsilon \leqslant t \leqslant 1-\varepsilon}\left|h^{\prime}(t, y)\right|<n\right\} .
$$

Then the functions $g_{n}(t)=\int \mathbb{1}_{A_{n}}(y) h(t, y) v(d y)$ are differentiable for $t \in(-1+\varepsilon, 1-\varepsilon)$. By the Mean Value Theorem we get, for $t$ and $t_{0}$ belonging to $(-1+\varepsilon, 1-\varepsilon)$,

$$
\frac{g_{n}(t)-g_{n}\left(t_{0}\right)}{t-t_{0}}=\int 1_{A_{n}}(y) h^{\prime}\left(t_{n}, y\right) v(d y)
$$

Choosing subsequence $t_{n_{k}} \rightarrow t^{*}$ and applying (i) and (ii) we get

$$
\frac{g(t)-g\left(t_{0}\right)}{t-t_{0}}=\int h^{\prime}\left(t^{*}, y\right) v(d y)
$$

Since $t^{*} \rightarrow t_{0}$ if $t \rightarrow t_{0}$, applying once more (i) and (ii) we finally obtain

$$
\frac{d g}{d t}\left(t_{0}\right)=\int h^{\prime}\left(t_{0}, y\right) v(d y)
$$

Now we recall some standard facts about Gaussian measures on Banach spaces [2].

Suppose that $(E,\|\cdot\|)$ is a separable Banach space and $E^{\prime}$ is its dual space. A Borel probability measure $\mu$ on $E$ is called centered Gaussian if every $\xi \in E^{\prime}$ is a real (symmetric) Gaussian random variable on the probability space $\left(E, \mathscr{B}_{E}, \mu\right)$, with $\mathscr{B}_{E}$ being the Borel $\sigma$-field of $E$. It is well-known that if $\mu$ is Gaussian, then there exists a (unique) $x \in E$ (called the barycenter of $\mu$ ) such that $\mu_{0}(\cdot)=\mu(\cdot+x)$ is centered.

Now, let $E^{\prime}(\mu)$ be the $L^{2}\left(\mu_{0}\right)$-closure of $E^{\prime}$, endowed with the usual $L^{2}\left(\mu_{0}\right)$ inner product. For every $\xi \in E^{\prime}(\mu)$ there exists a unique $\Lambda \xi \in E$ such that, for every $\eta \in E^{\prime}$, we have

$$
\eta(\Lambda \xi)=\int \eta \xi d \mu_{0}=\langle\eta, \xi\rangle_{\boldsymbol{H}(\mu)} .
$$

$\Lambda: E^{\prime}(\mu) \rightarrow E$ is a linear injective mapping.
The image of $\Lambda$, endowed with the inner product induced from $E^{\prime}(\mu)$, is denoted by $H(\mu)$ and is called the reproducing kernel Hilbert space (RKHS) of $(E, \mu)$. It is well-known that $H(\mu)$ is the space of all admissible translates of $\mu$.

The following theorem describes the density of the measure $\mu$ translated by some $x=\Lambda \xi \in H(\mu), \xi \in E^{\prime}(\mu)$ :

Cameron-Martins formula. Let $x=\Lambda \xi \in H(\mu)$, where $\xi \in E^{\prime}(\mu)$, with $\mu$ centered Gaussian. Then

$$
d \mu(\cdot-x)=\exp \left(\xi-\frac{1}{2}\|\xi\|_{H(\mu)}^{2}\right) d \mu
$$

The next theorem generalizes an immediate observation that one-dimensional symmetric Gaussian measure of translates of a fixed interval takes on the greatest value when this (translated) interval is symmetric with respect to the origin.

Andersons inequality. Let $\mu$ be a centered Gaussian measure and $V$ a Borel symmetric with respect to the origin convex set.

Then, for' every $x \in E, \mu(x+V) \leqslant \mu(V)$.
Lemma 2. Let $\mu$ be a Gaussian measure on $E$ with $\operatorname{dim}(\operatorname{supp} \mu)=\infty$. Assume that the barycenter of $\mu$ belongs to $H(\mu)$. Then, for every $x \neq 0$, $x \in H(\mu)$ and $r \in R$,

$$
\begin{equation*}
\sup _{-1 \leqslant t \leqslant 1} \int\|x+t y\|^{r} \mu(d y)<\infty \tag{1}
\end{equation*}
$$

Proof. For $r \geqslant 0$ the conclusion easily follows from the finiteness of all moments of any Gaussian measure.

It remains to prove our lemma for $r<0$. Write $V=\{y ;\|y\| \leqslant 1\}$. Then

$$
\begin{equation*}
\int\|x+t y\|^{r} \mu(d y)=(-r) \int_{0}^{\infty} u^{r-1} \mu_{0}\left(\frac{x}{t}+x_{0}+\frac{u}{t} V\right) d u \tag{2}
\end{equation*}
$$

where $\mu_{0}\left(\cdot+x_{0}\right)=\mu(\cdot)$.
Put $a=x / t+x_{0}, u / t=s$. Take a fixed $\xi \in E^{\prime}$ and let $b=\Lambda \xi, \eta=\Lambda^{-1} a$.

## By Cameron-Martin's Formula we get (see also [3])

$$
\begin{aligned}
\mu_{0}(a+s V) & =\exp \left(-\frac{1}{2}\|a\|_{H(\mu)}^{2}\right) \int_{s V} \exp (-\eta) d \mu_{0} \\
& \leqslant \exp \left(-\frac{1}{2}\|a\|_{H(\mu)}^{2}\right) \int_{s V} \exp \left(\xi-\inf _{s V} \xi-\eta\right) d \mu_{0} \\
& =\exp \left(-\frac{1}{2}\|a\|_{H(\mu)}^{2}-\inf _{s V} \xi\right) \int_{s V} \exp (\xi-\eta) d \mu_{0} .
\end{aligned}
$$

Since

$$
\begin{aligned}
& \int_{s V} \exp (\xi-\eta) d \mu_{0} \\
& =\exp \left\{\frac{1}{2}\|a-b\|_{H(\mu)}^{2}\right\} \int_{s V} \exp \left(\xi-\eta-\frac{1}{2}\|a-b\|_{H(\mu)}^{2}\right) d \mu_{0} \\
& =\exp \left\{\frac{1}{2}\|a-b\|_{H(\mu)}^{2}\right\} \mu_{0}(a-b+s V) \leqslant \exp \left\{\frac{1}{2}\|a-b\|_{H(\mu)}^{2}\right\} \mu_{0}(s V)
\end{aligned}
$$

we get

$$
\begin{equation*}
\mu_{0}(a+s V) \leqslant \exp \left(\frac{1}{2}\|a-b\|_{H(\mu)}^{2}-\frac{1}{2}\|a\|_{H(\mu)}^{2}-\inf _{s V} \xi\right) \mu_{0}(s V) \tag{3}
\end{equation*}
$$

We now choose a $w \in H(\mu)$ such that $\Lambda^{-1} w \in E^{\prime}$ and
(4)

$$
\|w-x\|_{H(\mu)}^{2}<\frac{1}{16}\|x\|_{H(\mu)}^{2}
$$

For $t_{0}>0$ which is small enough we have

$$
\begin{equation*}
\frac{1}{2}\left\|\frac{x}{t}+x_{0}\right\|_{H(\mu)}^{2}>\frac{1}{4 t^{2}}\|x\|_{H(\mu)}^{2} \quad \text { for }|t| \leqslant t_{0} \tag{5}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\|\frac{w-x}{t}+x_{0}\right\|_{H(\mu)}^{2} \leqslant \frac{2}{t^{2}}\|w-x\|_{H(\mu)}^{2} \tag{6}
\end{equation*}
$$

Moreover, if $\xi=\Lambda^{-1} w / t$, then, for some $u_{0}>0$, we get

$$
\begin{equation*}
\left|\inf _{(u / t) V} \xi\right| \leqslant \frac{1}{16 t^{2}}\|x\|_{H(\mu)}^{2} \quad \text { for } 0<u \leqslant u_{0} \tag{7}
\end{equation*}
$$

Combining (3)-(7) we have, for $|t| \leqslant t_{0}$ and $0<u \leqslant u_{0}$,

$$
\begin{equation*}
\mu_{0}\left(\frac{x}{t}+x_{0}+\frac{u}{t} V\right)<C \mu_{0}\left(\frac{u}{t} V\right) \exp \left(-\frac{D}{t^{2}}\right) \tag{8}
\end{equation*}
$$

for some positive constants $C$ and $D$. Now, (8) yields the following estimate for (2):

$$
\begin{aligned}
& \int_{0}^{\infty} u^{r-1} \mu_{0}\left(\frac{x}{t}+x_{0}+\frac{u}{t} V\right) d u \\
\leqslant & C \exp \left(-\frac{D}{t^{2}}\right) \int_{0}^{u_{0}} u^{r-1} \mu_{0}\left(\frac{u}{t} V\right) d u+\int_{u_{0}}^{\infty} u^{r-1} d u \\
\leqslant & \frac{C}{t^{r}} \exp \left(-\frac{D}{t^{2}}\right) \int_{0}^{\infty} u^{r-1} \mu_{0}(u V) d u-\frac{1}{r} u_{0}^{r} \quad \text { for }|t| \leqslant t_{0} .
\end{aligned}
$$

Since for every positive integer $n$ there exists a constant $C_{n}$ such that $\mu(u V) \leqslant C_{n} u^{n}$ for $0<u \leqslant 1 \quad[1]$, we obtain

$$
\int_{0}^{\infty} u^{r-1} \mu_{0}(u V) d u<\infty
$$

Finally, if $|t|>t_{0}$, then

$$
\int_{0}^{\infty} u^{r-1} \mu_{0}\left(\frac{x}{t}+x_{0}+\frac{u}{t} V\right) d u \leqslant \int_{0}^{\infty} u^{r-1} \mu_{0}\left(\frac{u}{t} V\right) d u=|t|^{r} \int_{0}^{\infty} u^{r-1} \mu_{0}(u V) d u
$$

which completes the proof of the lemma.
The last lemma and Corollary 1 give the following
Corollary 2. Let v be a Gaussian measure on $E=L_{M}$. Assume that the barycenter of $v$ belongs to $H(v)$ and let $u \neq 0, u \in H(v)$.

If $\operatorname{dim}(\operatorname{supp} v)=\infty$, then the functions $G^{q}(t)=\|u+t h\|_{M}^{q}$ and $d G^{q} / d t$ satisfy all assumptions of Lemma 1 for every $q>0$.

Before formulating our theorem, we again introduce some notation.
Write $T_{d}=\{\dot{x} ; m\{x\}>0\}$ and $T_{c}=T \backslash T_{d}$. Let

$$
R_{c}=\left\{\frac{t p^{\prime}(t)}{p(t)}+1 ; M(t) \in\left[\frac{1}{m(T)}, \infty\right)\right\}
$$

$$
R_{d}=\left\{\frac{t p^{\prime}(0)}{p(t)}+1 ; M(t)=\frac{1}{m\left(\left\{x_{\}}\right)\right.} \text {with } x \in T_{d}\right\}
$$

Theorem. Let $L_{M} \equiv L_{M}(T, F, m)$ be a separable Orlicz space, where $M$ satisifies (I)-(III). Assume that $\operatorname{dim} L_{M}=\infty, m \cdot$ is not purely atomic and that $q \notin R_{c} \cup R_{d}, q>0$.

If $v_{1}$ and $v_{2}$ are Gaussian measures on $L_{M}$ such that, for every $u \in L_{M}$,

$$
\begin{equation*}
\int\|u+h\|_{M}^{q} v_{1}(d h)=\int\|u+h\|_{M}^{q} v_{2}(d h), \tag{+}
\end{equation*}
$$

then $v_{1}=v_{2}$.
Proof. Denote $v_{1}-v_{2}$ by $v$. We show that for disjoint sets $A, B \in F$ the following conditions hold:

$$
\begin{equation*}
\int\left\langle\mathbf{1}_{A}, h\right\rangle v(d h)=0, \tag{a}
\end{equation*}
$$

(b)

$$
\int\left\langle\mathbb{1}_{A}, h\right\rangle\left\langle\mathbb{1}_{B}, h\right\rangle v(d h)=0,
$$

$$
\begin{equation*}
\int\left\langle\mathbf{1}_{A}, h\right\rangle^{2} v(d h)=0 \tag{c}
\end{equation*}
$$

This clearly implies that, for every continuous linear functional $\psi$ on $L_{M}$, the one-dimensional Gaussian measures $\psi\left(v_{1}\right)$ and $\psi\left(v_{2}\right)$ are identical, which, of course, yields the equality $v_{1}=v_{2}$.

Let $g_{i}$ be the barycenters of $v_{i}(i=1,2)$ and let $u$ be a fixed element of $L_{M}$ with values in $\{-1,0,1\}$. In order to apply Corollary 2 we construct a centered Gaussian random vector $Z$ with infinite-dimensional support satisfying $g_{1}, g_{2}, u \in H(\gamma)$, where $\gamma$ is the distribution of $Z$. Furthermore, if $\gamma_{s}$ is the distribution of $s Z, s>0$, then $\gamma_{s}$ as well as $v_{i} * \gamma_{s}(i=1,2)$ have all the above-mentioned properties.

Moreover, the property $(+)$ still holds with $v_{i}$ replaced by $v_{i} * \gamma_{s}$.
Applying Lemma 1 to $G^{q}$ and $d G^{q} / d t$ with $v$ replaced by $v * \gamma_{s}$, we get

$$
\int \frac{d G^{q}}{d t}(0) v * \gamma_{s}(d h)=0 \quad \text { and } \quad \int \frac{d^{2} G^{q}}{d t^{2}}(0) v * \gamma_{s}(d h)=0 .
$$

It is easy to see that these equations yield, as $s \rightarrow 0$,

$$
\begin{align*}
& \int \frac{d G^{q}}{d t}(0) v(d h)=0  \tag{9}\\
& \int \frac{d^{2} G^{q}}{d t^{2}}(0) v(d h)=0 \tag{10}
\end{align*}
$$

To conclude this part of the proof we have to construct a random vector $Z$ with the above listed properties. To do this, let $\left\{f_{i}\right\}_{i=1}^{\infty}$ be a sequence of linearly independent functions belonging to the unit sphere of $L_{M}$
and let $\xi_{-2}, \xi_{-1}, \xi_{0}, \xi_{1}, \ldots$ be a sequence of standard Gaussian random variables. It is easy to see that

$$
Z=g_{2} \xi_{-2}+g_{1} \xi_{-1}+u \xi_{0}+\sum_{i=1}^{\infty} f_{i} \xi_{i} \frac{1}{2^{i}}
$$

satisfies all the requirements.
Next, applying Corollary 1 to formulas (9) and (10) we get

$$
\begin{equation*}
\int\langle u, h\rangle v(d h)=0 \tag{11}
\end{equation*}
$$

and

$$
\begin{align*}
& \Phi_{1}(u) \int\|(1-|u|) \cdot h\|_{2}^{2} v(d h)+  \tag{12}\\
& \quad+\Phi_{2}(u) \int\||u| \cdot h\|_{2}^{2} v(d h)+\Phi_{3}(u) \int\langle u ; h\rangle^{2} v(d h)=0 .
\end{align*}
$$

Property (a) now follows from (11) if we put $u=\mathbf{1}_{A}, A \in F$. Further, if $A$ and $B$ are disjoint and of positive measure, then writing (12) for $u_{1}=1_{A}+\mathbf{1}_{B}$ and $u_{2}=\mathbb{1}_{A}-1_{B}$, and taking into account that $\Phi_{k}\left(u_{1}\right)=\Phi_{k}\left(u_{2}\right)(k=1,2,3)$ and that $\Phi_{3}\left(u_{2}\right) \neq 0$, we get (b).

Now, if $A \subset T_{c}, A \in F, m(A)>0$, then we construct a sequence $\left\{r_{n}\right\}$ such that $r_{0}=\mathbf{1}_{A}, \operatorname{supp} r_{n}=A, r_{n}= \pm 1$ with equal measure $=\frac{1}{2} m(A)$, and $\left\{r_{n}\right\}$ is an orthogonal sequence of $L^{2}(T, F, m)$. Substituting $r_{n}$ in place of $u$ in (12) for $n=1,2, \ldots$, we get the following system of linear equations:

$$
\begin{equation*}
D+\Phi_{3}\left(\mathbb{1}_{A}\right) \cdot \int\left\langle r_{n}, h\right\rangle^{2} v(d h)=0, \quad n=1,2, \ldots \tag{13}
\end{equation*}
$$

Adding first $n$ equations and dividing by $n$ we get

$$
D+\frac{1}{n} \Phi_{3}\left(\mathbb{1}_{A}\right) \sum_{k=0}^{n} \int\left\langle r_{k}, h\right\rangle^{2} v(d h)=0
$$

When $n \rightarrow \infty$, we get $D=0$. Substituting this value into (13) and taking into account that $\Phi_{3}\left(\mathbf{1}_{A}\right) \neq 0$ we obtain (c) for all $A \subset T_{c}, A \in F$.

Next, we show that (c) yields

$$
\begin{equation*}
\int\left\|\mathbb{1}_{A} h\right\|_{2}^{2} v(d h)=0 \tag{13}
\end{equation*}
$$

for all $A \subseteq T_{c}$. Indeed, let now $\left\{f_{k}\right\}$ denote an orthonormal basis of $L^{2}\left(T_{c}, F \cap T_{c}, m\right)$. Then

$$
\left\|\mathbf{1}_{A} h\right\|_{2}^{2}=\sum_{k=1}^{\infty}\left\langle f_{k}, \mathbf{1}_{A} h\right\rangle^{2}
$$

which, by virtue of (c), proves (13).
Observe that, in particular, (13) implies

$$
\begin{equation*}
p^{\prime}(0) \int\|h\|_{2}^{2} v(d h)=0 . \tag{14}
\end{equation*}
$$

For, putting $A=T_{c}$ in (c) and in (13) and using (12) we immediately get that $p^{\prime}(0)\left\|\mathbb{1}_{T_{d}} h\right\|_{2}^{2} v(d h)=0$, which, together with (13), gives (14).

To complete the proof, take $\left\{x_{0}\right\} \subseteq T_{d}$ and put $u=\mathbb{1}_{x_{0}}$ into (12). Applying (14) we obtain

$$
\left.\left(\Phi_{2}(u)-\Phi_{1}(u)+\Phi_{3}(u) m\left(\left\{x_{0}\right\}\right)\right)\right\}\|u h\|_{2}^{2} v(d h)=0
$$

or, equivalently,

$$
\begin{equation*}
\left(\|u\|_{M}(q-1)-\frac{p^{\prime}(0)}{p\left(1 /\|u\|_{M}\right)}\right) \Gamma\|u h\|_{2}^{2} d v(h)=0 \tag{15}
\end{equation*}
$$

Since $q \notin R_{d}$, (15) implies that (c) holds for $u=\mathbf{1}_{x_{0}}$, which clearly completes the proof.

Corollary 3. If the function $t \rightarrow p(t) / t$ is bounded or $p^{\prime}(0)=0$, then Theorem holds even if $m$ is purely atomic.

Proof. Boundedness of $t \rightarrow p(t) / t$ or condition $p^{\prime}(0)=0$ give that

$$
p^{\prime}(0) \int\left\|\mathbb{1}_{T_{d}} h\right\|_{2}^{2} v(d h)=0 .
$$

In fact, if $p^{\prime}(0) \neq 0$, then, putting in (12) $u_{n}=\mathbb{1}_{x_{n}}\left(\left\{x_{n}\right\}=T_{d}\right)$, we have

$$
\begin{equation*}
\int\|h\|_{2}^{2} v(d h)+C_{n} \int\left\|u_{n} h\right\|_{2}^{2} v(d h)=0, \tag{16}
\end{equation*}
$$

where

$$
C_{n}=\frac{\left\|u_{n}\right\|_{M} p\left(1 /\left\|u_{n}\right\|_{M}\right)}{p^{\prime}(0)}-1
$$

and $C_{n}(n=1,2, \ldots)$ is a bounded sequence. But $\sum_{n} \int\left\|u_{n} h\right\|_{2}^{2} v(d h)<\infty$, so

$$
\lim _{n \rightarrow \infty} C_{n} \int\left\|u_{n} h\right\|_{2}^{2} v(d h)=0 \quad\left(\operatorname{dim} L_{M}=\infty\right)
$$

whence $\int\|h\|_{2}^{2} v(d h)=0$, which completes the proof.
We say that $q>0$ is admissible for $L_{M}$ if $(+)$, satisfied for this particular $q$, implies that $v_{1}=v_{2}$.

Corollary 4. For $L_{M}=L^{r}$ with $r>2$ the exponent $q$ is admissible if $q \neq 1, q \neq r$. When $m$ is non-atomic or $r=2$, then $q$ is admissible if $q \neq r$.

For any $I_{M}$ with non-purely atomic measure $m$, every $q<1$ is admissible.
Remark. The following example shows that if $q=r$, then there exist two different Gaussian measures on $L^{r}$ such that $(+)$ is satisfied.

Example. Let $h_{1}, h_{2}, h$ be such functions on $T$ that $h_{1}=\mathbb{1}_{A}, h_{2}=\mathbb{1}_{A} c$, $h=h_{1}+h_{2}, A \in F, m(A)<m(T)$ and let $\theta_{1}, \theta_{2}$ be two independent and
standard Gaussian variables. If $X_{1}=h_{1} \theta_{1}+h_{2} \theta_{2}, X_{2}=h \theta_{1}$, then it is clear that $X_{1}$ and $X_{2}$ generate different measures on $L$.

We have, for any $f \in L^{\prime}$,

$$
\begin{aligned}
\mathrm{E}\left\|f+h_{1} \theta_{1}+h_{2} \theta_{2}\right\|_{r}^{r} & =\mathrm{E}\left(\underset{T}{\int}\left|f+h_{1} \theta_{1}+h_{2} \theta_{2}\right|^{r} d m\right) \\
& =\mathrm{E}\left(\int_{\boldsymbol{A}}\left|f+h_{1} \theta_{1}\right|^{r} d m\right)+\mathrm{E}\left(\int_{A^{c}}\left|f+h_{2} \theta_{2}\right|^{r} d m\right. \\
& =\mathrm{E}\left(\int_{\boldsymbol{A}}\left|f+h_{1} \theta_{1}\right|^{r} d m\right)+\mathrm{E}\left(\int_{A^{c}}\left|f+h_{2} \theta_{1}\right|^{r} d m\right) \\
& =\mathrm{E}\left\|f+h \theta_{1}\right\|_{r}^{r}
\end{aligned}
$$

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