# Eigenvalue distribution of optimal transportation 

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#### Abstract

We investigate the Brenier map $\nabla \Phi$ between the uniform measures on two convex domains in $\mathbb{R}^{n}$, or more generally, between two log-concave probability measures on $\mathbb{R}^{n}$. It is shown that the eigenvalues of the Hessian matrix $D^{2} \Phi$ exhibit remarkable concentration properties on a multiplicative scale, regardless of the choice of the two measures or the dimension $n$.


## 1 Introduction

Let $\mu$ and $\nu$ be two absolutely-continuous probability measures on $\mathbb{R}^{n}$. It was discovered by Brenier [4] and McCann [19] that there exists a convex function $\Phi$ on $\mathbb{R}^{n}$ with $(\nabla \Phi)_{*} \mu=\nu$, i.e.,

$$
\begin{equation*}
\int_{\mathbb{R}^{n}} b(\nabla \Phi(x)) d \mu(x)=\int_{\mathbb{R}^{n}} b(x) d \nu(x) \tag{1}
\end{equation*}
$$

for any $\nu$-integrable function $b: \mathbb{R}^{n} \rightarrow \mathbb{R}$. Moreover, the Brenier map $x \mapsto \nabla \Phi(x)$ is uniquely determined $\mu$-almost everywhere. In this paper we consider the case where $\mu$ and $\nu$ are log-concave probability measures. An absolutely-continuous probability measure on $\mathbb{R}^{n}$ is called log-concave if it has a density $\rho$ which satisfies

$$
\rho(\lambda x+(1-\lambda) y) \geq \rho(x)^{\lambda} \rho(y)^{1-\lambda} \quad\left(x, y \in \mathbb{R}^{n}, 0<\lambda<1\right)
$$

The uniform measure on any convex domain is log-concave, as well as the Gaussian measure. Write $\operatorname{Supp}(\mu)$ for the interior of the support of $\mu$, which is an open, convex set in $\mathbb{R}^{n}$. We make the assumption that
( $\star$ ) The function $\Phi$ is $C^{2}$-smooth in $\operatorname{Supp}(\mu)$.
It follows from the works of Caffarelli $[5,6,1]$ that $(\star)$ holds true when each of the measures $\mu$ and $\nu$ satisfies the following extra condition: Either the support of the measure is the entire $\mathbb{R}^{n}$, or else the support is a bounded, convex domain and density of the measure is bounded away from zero and from infinity in this convex domain. It is fair to say that Caffarelli's regularity theory covers most cases of interest, yet it is very plausible that $(\star)$ is in fact always correct, without any extra conditions.

[^0]As it turns out, the positive-definite Hessian matrix $D^{2} \Phi(x)$ exhibits remarkable regularity in the behavior of its eigenvalues. We write $\operatorname{Var}[X]$ for the variance of the random variable $X$.

Theorem 1.1. Let $\mu, \nu$ be absolutely-continuous, log-concave probability measures on $\mathbb{R}^{n}$. Let $\nabla \Phi$ be the Brenier map between $\mu$ and $\nu$, and assume ( $\star$ ). Write $0<\lambda_{1}(x) \leq \ldots \leq$ $\lambda_{n}(x)$ for the eigenvalues of the matrix $D^{2} \Phi(x)$, repeated according to their multiplicity. Let $X$ be a random vector in $\mathbb{R}^{n}$ that is distributed according to $\mu$. Then, for $i=1, \ldots, n$,

$$
\operatorname{Var}\left[\log \lambda_{i}(X)\right] \leq 4
$$

Thus, on a multiplicative scale, the eigenvalues of $D^{2} \Phi$ are quite stable. Note that the multiplicative scale is indeed the natural scale in the generality of Theorem 1.1: By applying appropriate linear transformations to $\mu$ and $\nu$, one may effectively multiply all eigenvalues by an arbitrary positive constant. The variance bound in Theorem 1.1 follows from a Poincaré inequality which we now formulate. For $x \in \operatorname{Supp}(\mu)$ set

$$
\Lambda(x)=\left(\log \lambda_{1}(x), \ldots, \log \lambda_{n}(x)\right) .
$$

We write $|\cdot|$ for the standard Euclidean norm in $\mathbb{R}^{n}$.
Theorem 1.2. Under the notation and assumptions of Theorem 1.1, for any locally-Lipschitz. function $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$ with $\mathbb{E}|f(\Lambda(X))|<\infty$,

$$
\operatorname{Var}[f(\Lambda(X))] \leq 4 \mathbb{E}|\nabla f|^{2}(\Lambda(X)),
$$

whenever the right-hand side is finite. At the points in which $f$ is not continuously differentiable, we define $|\nabla f|$ via (36) below.

Denote $\pi=\Lambda_{*}(\mu)$, the push-forward of the measure $\mu$ under the map $\Lambda$. Theorem 1.2 is a spectral gap estimate for the metric-measure space $\left(\mathbb{R}^{n},|\cdot|, \pi\right)$. Gromov and Milman [13] proved that a spectral gap estimate implies exponential concentration of Lipschitz functions. Therefore Theorem 1.2 admits the following immediate corollary:

Corollary 1.3. We work under the notation and assumptions of Theorem 1.1. Let $f: \mathbb{R}^{n} \rightarrow$ $\mathbb{R}$ be a 1 -Lipschitz function (i.e., $|f(x)-f(y)| \leq|x-y|)$. Denote $A=\mathbb{E} f(\Lambda(X)$ ). Then $A$ is finite and

$$
\mathbb{E} \exp (c|f(\Lambda(X))-A|) \leq 2,
$$

where $c>0$ is a universal constant.
Remark 1.4. Corollary 1.3 implies that $\mathbb{E} e^{c|\Lambda(X)|}<\infty$. Consequently, one may replace the condition $\mathbb{E}|f(\Lambda(X))|<\infty$ in Theorem 1.2 by the requirement that $e^{-c|x|}|f(x)|$ is bounded in $\mathbb{R}^{n}$, for a certain universal constant $c>0$.

Our next result is that the diagonal elements of the matrix $D^{2} \Phi(x)$ are also concentrated on a logarithmic scale, pretty much like the eigenvalues.

Theorem 1.5. We work under the notation and assumptions of Theorem 1.1. Fix $v \in \mathbb{R}^{n}$, let $H(x)=\log \left(D^{2} \Phi(x) v \cdot v\right)$ and denote $Y=H(X)$. Then,
(i) $\operatorname{Var}[Y] \leq 4$.
(ii) For any locally-Lipschitz function $f: \mathbb{R} \rightarrow \mathbb{R}$ with $\mathbb{E}|f(Y)|<\infty$,

$$
\operatorname{Var}[f(Y)] \leq 4 \mathbb{E}\left|f^{\prime}\right|^{2}(Y)
$$

(iii) For any 1-Lipschitz function $f: \mathbb{R} \rightarrow \mathbb{R}$, denoting $A=\mathbb{E} f(Y)$ we have that $A \in \mathbb{R}$ and

$$
\mathbb{E} \exp (c|f(Y)-A|) \leq 2,
$$

where $c>0$ is a universal constant.
All of the assertions made so far follow from Theorem 5.1 below, which is in fact a sound reformulation of [14, Theorem 1.4]. The results in [14] were obtained under a technical assumption dubbed "regularity at infinity", which we have to address in this paper. Our argument is based on analysis of the transportation metric: This means that we use the positive-definite Hessian $D^{2} \Phi$ in order to define a Riemannian metric in $\operatorname{Supp}(\mu)$. The weighted Riemannian manifold

$$
M_{\mu, \nu}=\left(\operatorname{Supp}(\mu), D^{2} \Phi, \mu\right)
$$

was studied in [17], where it was shown that the associated Ricci-Bakry-Émery tensor is non-negative when $\mu$ and $\nu$ are log-concave. We will also consider the map

$$
x \mapsto D^{2} \Phi(x)
$$

from $\operatorname{Supp}(\mu) \subseteq \mathbb{R}^{n}$ into the space of positive-definite matrices. The space of positivedefinite matrices is endowed with a natural Riemannian metric, which fits very nicely with computations related to the weighted Riemannian manifold $M_{\mu, \nu}$. This leads to a certain Poincaré inequality with respect to the standard Riemannian metric on the space of positivedefinite matrices, formulated in Theorem 5.1 below .

We have tried to make the exposition self-contained, apart from the regularity theory of mass-transport. The rest of this paper is organized as follows: In Section 2 we recall some well-known constructions related to positive-definite matrices. In Section 3 and Section 4 we prove the main results under regularity assumptions by employing the Bakry-Émery $\Gamma_{2^{-}}$ calculus. Section 5 is devoted to the elimination of these regularity assumptions. In Section 6 we complete the proofs of the theorems formulated above. Below we write $x \cdot y$ for the standard scalar product of $x, y \in \mathbb{R}^{n}$. We denote derivatives by $\partial_{k} f=f_{k}=\partial f / \partial x_{k}$ and $f_{i j}=\partial^{2} f /\left(\partial x_{i} \partial x_{j}\right)$. By a smooth function we mean a $C^{\infty}$-smooth one. We write $\log$ for the natural logarithm, and $\operatorname{Tr}(A)$ stands for the trace of the matrix $A$.

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## 2 Positive-definite quadratic forms

This section surveys standard material on positive-definite matrices. Denote by $M_{n}^{+}(\mathbb{R})$ the collection of all symmetric, positive-definite $n \times n$ matrices. For a function $f:(0, \infty) \rightarrow \mathbb{R}$ and $A \in M_{n}^{+}(\mathbb{R})$ we may define the symmetric matrix $f(A)$ via the spectral theorem. In other words,

$$
f\left(\sum_{i=1}^{n} \lambda_{i} v_{i} \otimes v_{i}\right)=\sum_{i=1}^{n} f\left(\lambda_{i}\right) v_{i} \otimes v_{i}
$$

for any orthonormal basis $v_{1}, \ldots, v_{n} \in \mathbb{R}^{n}$ and $\lambda_{1}, \ldots, \lambda_{n}>0$, where we write $x \otimes x=$ $\left(x_{i} x_{j}\right)_{i, j=1, \ldots, n}$ for $x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n}$.
Lemma 2.1. For any $A, B \in M_{n}^{+}(\mathbb{R})$,

$$
\begin{equation*}
\left\|\log \left(A^{1 / 2} B A^{1 / 2}\right)\right\|_{H S} \leq\|\log (A)\|_{H S}+\|\log (B)\|_{H S} \tag{2}
\end{equation*}
$$

where $\|\cdot\|_{H S}$ stands for the Hilbert-Schmidt norm.
Proof. For an $n \times n$ matrix $T$ and $k=1, \ldots, n$ we define

$$
\begin{equation*}
D_{k}(T)=\sup _{\substack{E \subseteq \mathbb{R}^{n} \\ \operatorname{dim}(E)=k}} \frac{\operatorname{Vol}_{k}\left(T\left(B^{n} \cap E\right)\right)}{\operatorname{Vol}_{k}\left(B^{n} \cap E\right)} \tag{3}
\end{equation*}
$$

where $B^{n}=\left\{x \in \mathbb{R}^{n} ;|x|<1\right\}$, and the supremum in (3) runs over all $k$-dimensional subspaces in $\mathbb{R}^{n}$. Thus, an application of the linear transformation $A$ may increase $k$ dimensional volumes by a factor of at most $D_{k}(A)$. It follows that for any $n \times n$ matrices $A$ and $B$,

$$
\begin{equation*}
D_{k}(A B) \leq D_{k}(A) D_{k}(B) \quad(k=1, \ldots, n) \tag{4}
\end{equation*}
$$

In the case in which $A \in M_{n}^{+}(\mathbb{R})$, we have $D_{k}(A)=\prod_{i=1}^{k} \lambda_{i}$, where $\lambda_{1} \geq \lambda_{2} \geq \ldots \geq$ $\lambda_{n}>0$ are the eigenvalues of $A$. Assume that $A, B \in M_{n}^{+}(\mathbb{R})$. Denote the eigenvalues of the symmetric, positive-definite matrix $A^{1 / 2} B A^{1 / 2}$ by $e^{\gamma_{1}} \geq \ldots \geq e^{\gamma_{n}}>0$. Then, for $k=1, \ldots, n$,

$$
\begin{equation*}
\prod_{i=1}^{k} e^{\gamma_{i}}=D_{k}\left(A^{1 / 2} B A^{1 / 2}\right) \leq D_{k}\left(A^{1 / 2}\right) D_{k}(B) D_{k}\left(A^{1 / 2}\right)=D_{k}(A) D_{k}(B)=\prod_{i=1}^{k}\left(e^{\alpha_{i}} e^{\beta_{i}}\right) \tag{5}
\end{equation*}
$$

where $e^{\alpha_{1}} \geq \ldots \geq e^{\alpha_{n}}>0$ are the eigenvalues of $A$ and $e^{\beta_{1}} \geq \ldots \geq e^{\beta_{n}}>0$ are the eigenvalues of $B$. We will now apply a lemma of Weyl [25], see also Polya [20]. According to the inequality of Weyl and Polya, the inequalities (5) entail that

$$
\begin{equation*}
\sum_{i=1}^{n} h\left(\gamma_{i}\right) \leq \sum_{i=1}^{n} h\left(\alpha_{i}+\beta_{i}\right) \tag{6}
\end{equation*}
$$

for any convex, non-decreasing function $h: \mathbb{R} \rightarrow \mathbb{R}$. For $t \in \mathbb{R}$ denote $t_{+}=\max \{t, 0\}$. The function $t \mapsto\left(t_{+}\right)^{2}$ is convex and non-decreasing, hence from (6),

$$
\begin{equation*}
\sum_{i=1}^{n}\left(\left(\gamma_{i}\right)_{+}\right)^{2} \leq \sum_{i=1}^{n}\left(\left(\alpha_{i}+\beta_{i}\right)_{+}\right)^{2} \tag{7}
\end{equation*}
$$

By using (4) for the inverse matrices, we conclude that for $k=1, \ldots, n$,

$$
\prod_{i=n-k+1}^{n} e^{-\gamma_{i}}=D_{k}\left(A^{-1 / 2} B^{-1} A^{-1 / 2}\right) \leq D_{k}\left(A^{-1}\right) D_{k}\left(B^{-1}\right)=\prod_{i=n-k+1}^{n}\left(e^{-\alpha_{i}} e^{-\beta_{i}}\right)
$$

The inequality of Weyl and Polya now implies that $\sum_{i=1}^{n} h\left(-\gamma_{i}\right) \leq \sum_{i=1}^{n} h\left(-\alpha_{i}-\beta_{i}\right)$ for any convex, non-decreasing function $h$. By using again $h(t)=\left(t_{+}\right)^{2}$, we get

$$
\begin{equation*}
\sum_{i=1}^{n}\left(\left(-\gamma_{i}\right)_{+}\right)^{2} \leq \sum_{i=1}^{n}\left(\left(-\alpha_{i}-\beta_{i}\right)_{+}\right)^{2} . \tag{8}
\end{equation*}
$$

Adding (7) and (8), we finally obtain

$$
\begin{equation*}
\sum_{i=1}^{n} \gamma_{i}^{2} \leq \sum_{i=1}^{n}\left(\alpha_{i}+\beta_{i}\right)^{2} \leq\left(\sqrt{\sum_{i=1}^{n} \alpha_{i}^{2}}+\sqrt{\sum_{i=1}^{n} \beta_{i}^{2}}\right)^{2} \tag{9}
\end{equation*}
$$

where we used the Cauchy-Schwartz inequality in the last passage. By taking the square root of (9) we deduce (2).

For two matrices $A, B \in M_{n}^{+}(\mathbb{R})$ set

$$
\begin{equation*}
\operatorname{dist}(A, B)=\left\|\log \left(A^{-1 / 2} B A^{-1 / 2}\right)\right\|_{H S} . \tag{10}
\end{equation*}
$$

Equivalently, $\operatorname{dist}(A, B)$ equals $\sqrt{\sum_{i} \log ^{2} \lambda_{i}}$, where $\lambda_{1}, \ldots, \lambda_{n}>0$ are the eigenvalues of the matrix $A^{-1} B$ which is conjugate to $A^{-1 / 2} B A^{-1 / 2}$. The latter equivalent definition of dist shows that for any invertible $n \times n$ matrix $T$,

$$
\begin{equation*}
\operatorname{dist}(A, B)=\operatorname{dist}\left(T^{t} A T, T^{t} B T\right) \quad\left(A, B \in M_{n}^{+}\left(\mathbb{R}^{n}\right)\right) \tag{11}
\end{equation*}
$$

where $A^{t}$ is the transpose of the matrix $A$. Observe also that $\operatorname{dist}(A, B)=\operatorname{dist}\left(A^{-1}, B^{-1}\right)$ for any $A, B \in M_{n}^{+}(\mathbb{R})$. Lemma 2.1 states that for $A, B \in M_{n}^{+}(\mathbb{R})$,

$$
\begin{equation*}
\operatorname{dist}(A, B) \leq \operatorname{dist}(A, \mathrm{Id})+\operatorname{dist}(\mathrm{Id}, B) \tag{12}
\end{equation*}
$$

where Id is the identity matrix. From (11) and (12) one realizes that dist satisfies the triangle inequality in $M_{n}^{+}(\mathbb{R})$, hence it is a metric. For $A \in M_{n}^{+}\left(\mathbb{R}^{n}\right)$ and a symmetric $n \times n$ matrix $B$ we denote

$$
\|B\|_{A}=\left\|A^{-1 / 2} B A^{-1 / 2}\right\|_{H S}=\sqrt{\operatorname{Tr}\left[\left(A^{-1} B\right)^{2}\right]} .
$$

For a smooth curve $\gamma:[a, b] \rightarrow M_{n}^{+}(\mathbb{R})$ set

$$
\begin{equation*}
\operatorname{Length}(\gamma)=\int_{a}^{b}\|\dot{\gamma}(s)\|_{\gamma(s)} d s \tag{13}
\end{equation*}
$$

where $\dot{\gamma}(s)=\frac{d \gamma(s)}{d s}$ is a symmetric $n \times n$ matrix. Then Length is invariant under conjugations. That is, the length of the curve $\gamma(s)$ equals that of the curve $T^{t} \gamma(s) T$ for any invertible $n \times n$ matrix $T$.

Lemma 2.2. (i) For any $A \in M_{n}^{+}\left(\mathbb{R}^{n}\right)$ and a symmetric $n \times n$ matrix $B$,

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0} \frac{\operatorname{dist}^{2}(A+\varepsilon B, A)}{\varepsilon^{2}}=\|B\|_{A}^{2}=\operatorname{Tr}\left[\left(A^{-1} B\right)^{2}\right] . \tag{14}
\end{equation*}
$$

(ii) Let $A, B \in M_{n}^{+}\left(\mathbb{R}^{n}\right)$ and consider the curve

$$
\gamma_{A, B}(s)=A^{1 / 2}\left(A^{-1 / 2} B A^{-1 / 2}\right)^{s} A^{1 / 2} \quad(0 \leq s \leq 1)
$$

Then $\gamma_{A, B}$ is a curve connecting $A$ and $B$ with Length $\left(\gamma_{A, B}\right)=\operatorname{dist}(A, B)$.
Proof. The invariance property (11) implies that

$$
\operatorname{dist}(A+\varepsilon B, A)=\operatorname{dist}\left(\operatorname{Id}+\varepsilon A^{-1 / 2} B A^{-1 / 2}, \mathrm{Id}\right) .
$$

It therefore suffices to prove (i) under the additional assumption that $A=\operatorname{Id}$. Let $\lambda_{1}, \ldots, \lambda_{n}>$ 0 be the eigenvalues of $B$. It follows from (10) that

$$
\lim _{\varepsilon \rightarrow 0} \frac{\operatorname{dist}^{2}(\mathrm{Id}+\varepsilon B, \mathrm{Id})}{\varepsilon^{2}}=\lim _{\varepsilon \rightarrow 0} \frac{\sum_{i=1}^{n} \log ^{2}\left(1+\varepsilon \lambda_{i}\right)}{\varepsilon}=\sum_{i=1}^{n} \lambda_{i}^{2}
$$

and (i) follows from the fact that $\|B\|_{A}^{2}=\sum_{i} \lambda_{i}^{2}$. We now turn to the proof of (ii). Again, we may reduce matters to the case where $A=\mathrm{Id}$ by noting that

$$
\gamma_{A, B}(s)=A^{1 / 2} \gamma_{\mathrm{Id}, A^{-1 / 2} B A^{-1 / 2}}(s) A^{1 / 2} \quad(0 \leq s \leq 1)
$$

Abbreviate $\gamma(s)=\gamma_{A, B}(s)=\gamma_{\mathrm{Id}, B}(s)$. Since $\gamma(s)=B^{s}$ then $\dot{\gamma}(s)=B^{s} \log (B)$ and hence, for any $0 \leq s \leq 1$,

$$
\|\dot{\gamma}(s)\|_{\gamma(s)}=\left\|B^{-s / 2}\left(B^{s} \log (B)\right) B^{-s / 2}\right\|_{H S}=\|\log (B)\|_{H S}=\operatorname{dist}(\mathrm{Id}, B) .
$$

From the definition (13) it follows that $\operatorname{Length}(\gamma)=\operatorname{dist}(\operatorname{Id}, B)$, and (ii) is proven.

The right-hand side of (14) depends quadratically on $B$, and therefore Lemma 2.2 tells us that our distance function dist on $M_{n}^{+}(\mathbb{R})$ is induced by a Riemannian metric. We refer to this Riemannian metric as the standard Riemannian metric on $M_{n}^{+}(\mathbb{R})$. The next two lemmas describe certain Lipschitz functions on $M_{n}^{+}(\mathbb{R})$.

Lemma 2.3. Fix $v \in \mathbb{R}^{n}$ and set $f(A)=\log (A v \cdot v)$ for $A \in M_{n}^{+}(\mathbb{R})$. Then $f$ is a 1 -Lipschitz function with respect to the standard Riemannian metric on $M_{n}^{+}(\mathbb{R})$.

Proof. The map $f$ is clearly smooth. Fix $A \in M_{n}^{+}(\mathbb{R})$ and let us show that the norm of the Riemannian gradient of $f$ at the point $A$ is bounded by one. For any symmetric $n \times n$ matrix $B$ we have

$$
\left.\frac{d}{d t} f(A+t B)\right|_{t=0}=\frac{B v \cdot v}{A v \cdot v}
$$

Thus, in order to prove the lemma, it suffices to show that

$$
\begin{equation*}
\frac{B v \cdot v}{A v \cdot v} \leq\|B\|_{A}=\left\|A^{-1 / 2} B A^{-1 / 2}\right\|_{H S} . \tag{15}
\end{equation*}
$$

By switching to another orthonormal basis, if necessary, we may assume that $A$ is a diagonal matrix. Denote by $\lambda_{1}, \ldots, \lambda_{n}>0$ the numbers on the diagonal of $A$. Denote $B=\left(b_{i j}\right)_{i, j=1, \ldots, n}$ and $v=\left(v_{1}, \ldots, v_{n}\right) \in \mathbb{R}^{n}$. From the Cauchy-Schwartz inequality,

$$
\sum_{i, j=1}^{n} b_{i j} v_{i} v_{j} \leq \sqrt{\sum_{i, j=1}^{n} \frac{b_{i j}^{2}}{\lambda_{i} \lambda_{j}}} \sqrt{\sum_{i, j=1}^{n} \lambda_{i} \lambda_{j} v_{i}^{2} v_{j}^{2}}=\sqrt{\sum_{i, j=1}^{n} \frac{b_{i j}^{2}}{\lambda_{i} \lambda_{j}}}\left(\sum_{i=1}^{n} \lambda_{i} v_{i}^{2}\right)
$$

which is equivalent to the desired inequality (15).

Lemma 2.4. For $A \in M_{n}^{+}(\mathbb{R})$ denote its eigenvalues by $\lambda_{1}(A) \geq \ldots \geq \lambda_{n}(A)>0$. Consider the map $\Lambda: M_{n}^{+}(\mathbb{R}) \rightarrow \mathbb{R}^{n}$ defined via

$$
\begin{equation*}
\Lambda(A)=\left(\log \left(\lambda_{1}(A)\right), \ldots, \log \left(\lambda_{n}(A)\right)\right) \tag{16}
\end{equation*}
$$

Then $\Lambda$ is a 1-Lipschitz map, with respect to the standard Riemannian metric on $M_{n}^{+}(\mathbb{R})$, and the standard Euclidean metric on $\mathbb{R}^{n}$.

Proof. Write $\mathcal{F} \subseteq M_{n}^{+}(\mathbb{R})$ the collection of all positive-definite, symmetric matrices with $n$ distinct eigenvalues. Then $\mathcal{F}$ is an open, dense set. The function $\Lambda$ is continuous, since the eigenvalues vary continuously with the matrix. It therefore suffices to prove that

$$
\left|\Lambda\left(A_{1}\right)-\Lambda\left(A_{2}\right)\right| \leq \operatorname{dist}\left(A_{1}, A_{2}\right) \quad \text { for } A_{1}, A_{2} \in \mathcal{F}
$$

Fix $A_{1}, A_{2} \in \mathcal{F}$. Consider the curve $\gamma(s)=\gamma_{A_{1}, A_{2}}\left(s / \operatorname{dist}\left(A_{1}, A_{2}\right)\right)$ where $\gamma_{A_{1}, A_{2}}(s)$ is as in Lemma 2.2. Then $\gamma$ is a length-minimizing curve between $A_{1}$ and $A_{2}$, parametrized by Riemannian arclength. We claim that $\gamma(s) \in \mathcal{F}$ for all but finitely many values of $s$. Indeed, the resultant of $\gamma(s)$ is a real-analytic function of $s$ which is not identically zero, hence its zeros are isolated. Since $\Lambda \circ \gamma$ is continuous, in order to prove the lemma it suffices to show that

$$
\begin{equation*}
\left|\frac{d \Lambda(\gamma(s))}{d s}\right| \leq 1 \tag{17}
\end{equation*}
$$

for all $s$ with $\gamma(s) \in \mathcal{F}$. Let us fix $s_{0}$ with $\gamma\left(s_{0}\right) \in \mathcal{F}$. Denote $A=\gamma\left(s_{0}\right)$ and $B=\dot{\gamma}\left(s_{0}\right)$. Since $\gamma$ is parameterized by arclength, then

$$
\begin{equation*}
\|B\|_{A}=\left\|A^{-1 / 2} B A^{-1 / 2}\right\|_{H S}=1 \tag{18}
\end{equation*}
$$

Let $v_{1}, \ldots, v_{n} \in \mathbb{R}^{n}$ be the orthonormal basis of eigenvectors that corresponds to the eigenvalues $\lambda_{1}(A), \ldots, \lambda_{n}(A)$ of the matrix $A$. Then,

$$
\begin{equation*}
\left.\frac{d \lambda_{i}(\gamma(s))}{d s}\right|_{s=s_{0}}=B v_{i} \cdot v_{i} \quad(i=1, \ldots, n) \tag{19}
\end{equation*}
$$

The relation (19) is standard, see, e.g. Reed and Simon [21, Section XII.1]. Consequently,

$$
\begin{equation*}
\left.\frac{d \Lambda(\gamma(s))}{d s}\right|_{s=s_{0}}=\left(\frac{B v_{1} \cdot v_{1}}{\lambda_{1}(A)}, \ldots, \frac{B v_{n} \cdot v_{n}}{\lambda_{n}(A)}\right) \tag{20}
\end{equation*}
$$

However, by (18),

$$
\begin{equation*}
\sum_{i=1}^{n}\left(\frac{B v_{i} \cdot v_{i}}{\lambda_{i}(A)}\right)^{2}=\sum_{i=1}^{n}\left(A^{-1 / 2} B A^{-1 / 2} v_{i} \cdot v_{i}\right)^{2} \leq\left\|A^{-1 / 2} B A^{-1 / 2}\right\|_{H S}^{2}=1 \tag{21}
\end{equation*}
$$

Now (17) follows from (20) and (21).

Corollary 2.5. Whenever $A$ and $B$ are positive-definite $n \times n$ matrices,

$$
\sum_{i=1}^{n} \log ^{2} \frac{\lambda_{i}}{\mu_{i}} \leq\left\|\log \left(A^{-1 / 2} B A^{-1 / 2}\right)\right\|_{H S}^{2}
$$

where $\lambda_{1} \geq \ldots \geq \lambda_{n}>0$ are the eigenvalues of $A$ and $\mu_{1} \geq \ldots \geq \mu_{n}>0$ are the eigenvalues of $B$.

## 3 Bakry-Émery $\Gamma_{2}$-calculus

Let $\mu$ and $\nu$ be two absolutely-continuous, log-concave probability measures on $\mathbb{R}^{n}$. Assume that $d \mu=e^{-V(x)} d x$ and $d \nu=e^{-W(x)} d x$, for certain smooth, convex functions $V, W$ : $\mathbb{R}^{n} \rightarrow \mathbb{R}$. Let $\nabla \Phi$ be the Brenier map between $\mu$ and $\nu$. Caffarelli's regularity theory states that $\Phi: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is a smooth, convex function. Therefore (1) implies that the transport equation

$$
\begin{equation*}
-V(x)=\log \operatorname{det} D^{2} \Phi(x)-W(\nabla \Phi(x)) \tag{22}
\end{equation*}
$$

holds everywhere in $\mathbb{R}^{n}$. In particular, the matrix $D^{2} \Phi(x)=\left(\Phi_{i j}(x)\right)_{i, j=1, \ldots, n}$ is invertible and hence positive-definite for any $x \in \mathbb{R}^{n}$. The inverse matrix to $D^{2} \Phi(x)$ is denoted by $\left(D^{2} \Phi(x)\right)^{-1}=\left(\Phi^{i j}(x)\right)_{i, j=1, \ldots, n}$. We use the Einstein summation convention, thus an index that appears twice in an expression, once as a subscript and once as a superscript, is being summed upon. We also use abbreviations such as $\Phi_{j k}^{i}=\Phi^{i \ell} \Phi_{j k \ell}$ and $\Phi_{k}^{i j}=$ $\Phi^{i \ell} \Phi^{j m} \Phi_{k m \ell}$. Differentiating (22), we obtain

$$
\begin{equation*}
V_{j}(x)=-\Phi_{j i}^{i}(x)+\sum_{i=1}^{n} \Phi_{i j}(x) W_{i}(\nabla \Phi(x)) \quad\left(j=1, \ldots, n, x \in \mathbb{R}^{n}\right) \tag{23}
\end{equation*}
$$

Following [17], we use the positive-definite matrices $D^{2} \Phi(x)$ in order to induce a Riemannian metric on $\mathbb{R}^{n}$, and consider the weighted Riemannian manifold

$$
M=M_{\mu, \nu}=\left(\mathbb{R}^{n}, D^{2} \Phi, \mu\right)
$$

See Grigor'yan [11] and Bakry, Gentil and Ledoux [3] for background on weighted Riemannian manifolds and the $\Gamma_{2}$-calculus. For a smooth function $u: \mathbb{R}^{n} \rightarrow \mathbb{R}$ we have $\left|\nabla_{M} u\right|_{M}^{2}=\Phi^{i j} u_{i} u_{j}$ where $\left|\nabla_{M} u\right|_{M}^{2}$ stands for the square of the Riemannian norm of the Riemannian gradient of $u$. The Dirichlet form associated with the weighted Riemannian manifold $M_{\mu, \nu}$ is defined, for smooth functions $u, v: \mathbb{R}^{n} \rightarrow \mathbb{R}$, via

$$
\Gamma(u, v)=\int_{\mathbb{R}^{n}}\left\langle\nabla_{M} u, \nabla_{M} v\right\rangle_{M} d \mu=\int_{\mathbb{R}^{n}}\left(\Phi^{i j} u_{i} v_{j}\right) d \mu
$$

whenever the integral converges. The Laplacian associated with the weighted Riemannian manifold $M_{\mu, \nu}$ is defined, for a smooth function $u: \mathbb{R}^{n} \rightarrow \mathbb{R}$, by

$$
\begin{equation*}
L u=\Phi^{i j} u_{i j}-\sum_{j=1}^{n} W_{j}(\nabla \Phi(x)) u_{j}=\Phi^{i j} u_{i j}-\left(\Phi_{i}^{i j}+\Phi^{i j} V_{i}\right) u_{j} \tag{24}
\end{equation*}
$$

where the last equality holds in view of (23). Integrating by parts, we verify that
$-\int_{\mathbb{R}^{n}}(L u) v d \mu=-\int_{\mathbb{R}^{n}}\left(\Phi^{i j} u_{i j}-\left[\Phi_{i}^{i j}+\Phi^{i j} V_{i}\right] u_{j}\right) v e^{-V}=\int_{\mathbb{R}^{n}}\left(\Phi^{i j} u_{i} v_{j}\right) d \mu=\Gamma(u, v)$
for any smooth functions $u, v: \mathbb{R}^{n} \rightarrow \mathbb{R}$ with one of them compactly-supported. The next step is to consider the Carré du Champ of $M_{\mu, \nu}$ : As in Bakry and Émery [2], for a smooth function $u: K \rightarrow \mathbb{R}$ we define

$$
\begin{equation*}
\Gamma_{2}(u)=\frac{1}{2} L\left(\left|\nabla_{M} u\right|_{M}^{2}\right)-\left\langle\nabla_{M} u, \nabla_{M}(L u)\right\rangle_{M}=\frac{1}{2} L\left(\Phi^{i j} u_{i} u_{j}\right)-\Phi^{i j}(L u)_{i} u_{j} \tag{25}
\end{equation*}
$$

Lemma 3.1. For any smooth function $u: \mathbb{R}^{n} \rightarrow \mathbb{R}$ we have the pointwise inequality

$$
\Gamma_{2}(u) \geq \frac{1}{4} \Phi_{\ell}^{i k} \Phi_{k}^{j \ell} u_{i} u_{j}
$$

Lemma 3.1 is proven in [14] by introducing a Kähler structure and interpreting the lefthand side of (26) below as the Hilbert-Schmidt norm of a certain Hessian operator restricted to a subspace. There are several additional ways to prove Lemma 3.1. The brute-force way, is to use a tedious but straightforward computation which shows that
$\Gamma_{2}(u)=\Phi^{k l} \Phi^{i j} u_{i k} u_{j \ell}-\Phi^{i j k} u_{i j} u_{k}+\frac{1}{2}\left(\Phi_{\ell}^{i k} \Phi_{k}^{j \ell}+\Phi^{i k} \Phi^{j \ell} V_{k \ell}\right) u_{i} u_{j}+\frac{1}{2} \sum_{i, j=1}^{n}\left(W_{i j} \circ \nabla \Phi\right) u_{i} u_{j}$.
This computation is more or less equivalent to reproving Bochner's formula. Then, one proves the pointwise inequality

$$
\begin{equation*}
\Phi^{k l} \Phi^{i j} u_{i k} u_{j \ell}-\Phi^{i j k} u_{i j} u_{k}+\frac{1}{4} \Phi_{\ell}^{i k} \Phi_{k}^{j \ell} u_{i} u_{j} \geq 0 \tag{26}
\end{equation*}
$$

by representing the left-hand side of (26) as the trace of the square of the matrix $B=$ $\left(b_{i}^{j}\right)_{i, j=1, \ldots, n}$ where $b_{i}^{j}=\Phi^{j k} u_{k i}-\frac{1}{2} \Phi_{i}^{j k} u_{k}$. The product $A=\left(D^{2} \Phi\right) B$ is a symmetric matrix, hence

$$
\operatorname{Tr}\left(B^{2}\right)=\operatorname{Tr}\left[\left(\left(D^{2} \Phi\right)^{-1 / 2} A\left(D^{2} \Phi\right)^{-1 / 2}\right)^{2}\right] \geq 0
$$

Lemma 3.1 follows from (26) and from the fact that $D^{2} V$ and $D^{2} W$ are positive semidefinite matrices.

Another approach to Lemma 3.1 is to use the notation of Riemannian geometry as in [17], and use the Bochner formula. We first observe that identity (23) in the case $j=1$ has the simple form

$$
\begin{equation*}
L \Phi_{1}=-V_{1} \tag{27}
\end{equation*}
$$

Differentiating (27) and using $\partial_{k}\left(\Phi^{i j}\right)=-\Phi_{k}^{i j}$, we obtain

$$
\begin{equation*}
L\left(\Phi_{11}\right)-\Phi_{1}^{j k} \Phi_{1 j k}-\sum_{j, k=1}^{n} \Phi_{j 1} \Phi_{1 k}\left(W_{j k} \circ \nabla \Phi\right)=-V_{11} \tag{28}
\end{equation*}
$$

The Bochner-Lichnerowicz-Weitzenböck formula states that for any smooth $u: \mathbb{R}^{n} \rightarrow \mathbb{R}$,

$$
\begin{equation*}
\Gamma_{2}(u)=\left\|D_{M}^{2} u\right\|_{M}^{2}+\operatorname{Ric}_{M}\left(\nabla_{M} u, \nabla_{M} u\right) \tag{29}
\end{equation*}
$$

where $\left\|D_{M}^{2} u\right\|_{M}^{2}$ is the Hilbert-Schmidt norm of the Riemannian Hessian of $u$, and $\operatorname{Ric}_{M}$ is the Bakry-Émery-Ricci tensor of the weighted Riemmannian manifold $M=M_{\mu, \nu}$. Let us analyze the term in (29) involving the Hessian of $u$. The Christofell symbols of our Riemannian metric are $\Gamma_{i j}^{k}=\frac{1}{2} \Phi_{i j}^{k}$, and therefore $\left(D_{M}^{2} u\right)_{i j}=u_{i j}-\frac{1}{2} \Phi_{i j}^{k} u_{k}$ and

$$
\left\|D_{M}^{2} u\right\|_{M}^{2}=\Phi^{i k} \Phi^{j m}\left(u_{i j}-\frac{1}{2} \Phi_{i j}^{\ell} u_{\ell}\right)\left(u_{m k}-\frac{1}{2} \Phi_{m k}^{s} u_{s}\right)
$$

In the particular case of $u=\Phi_{1}$, we obtain $\left(D_{M}^{2} \Phi_{1}\right)_{j k}=\frac{1}{2} \Phi_{1 j k}$ and hence $\left\|D_{M}^{2} \Phi_{1}\right\|_{M}^{2}=$ $\frac{1}{4} \Phi_{1 j}^{k} \Phi_{1 k}^{j}$. Furthermore, the vector field $\nabla_{M} \Phi_{1}$ satisfies $\nabla_{M} \Phi_{1}=\partial / \partial x_{1}$ and $\left|\nabla_{M} \Phi_{1}\right|_{M}^{2}=$ $\Phi_{11}$. Since $L \Phi_{1}=-V_{1}$, the Bochner formula (29) for $u=\Phi_{1}$ takes the form

$$
\begin{align*}
\frac{1}{2} L\left(\Phi_{11}\right) & =-\left\langle\nabla_{M} \Phi_{1}, \nabla_{M} V_{1}\right\rangle_{M}+\frac{1}{4} \Phi_{1 j}^{k} \Phi_{1 k}^{j}+\operatorname{Ric}_{M}\left(\nabla_{M} u, \nabla_{M} u\right) \\
& =-V_{11}+\frac{1}{4} \Phi_{1 j}^{k} \Phi_{1 k}^{j}+\left(\operatorname{Ric}_{M}\right)_{11} \tag{30}
\end{align*}
$$

From (28) and (30) we obtain a formula for the Bakry-Émery-Ricci tensor:

$$
\left(R i c_{M}\right)_{11}=\frac{1}{4} \Phi_{1 j}^{k} \Phi_{1 k}^{j}+\frac{1}{2} V_{11}+\frac{1}{2} \sum_{j, k=1}^{n} \Phi_{j 1} \Phi_{1 k}\left(W_{j k} \circ \nabla \Phi\right)
$$

It is clear that there is nothing special about the derivative $u=\Phi_{1}$, and that we could have repeated the argument with $u=\nabla \Phi \cdot \theta$ for any $\theta \in \mathbb{R}^{n}$. We thus obtain the formula

$$
\begin{equation*}
\left(R i c_{M}\right)_{i \ell}=\frac{1}{4} \Phi_{i j}^{k} \Phi_{\ell k}^{j}+\frac{1}{2} V_{i \ell}+\frac{1}{2} \sum_{j, k=1}^{n} \Phi_{j i} \Phi_{\ell k}\left(W_{j k} \circ \nabla \Phi\right) \tag{31}
\end{equation*}
$$

Since $D^{2} V$ and $D^{2} W$ are positive semi-definite, then for any smooth $u: \mathbb{R}^{n} \rightarrow \mathbb{R}$,

$$
\Gamma_{2}(u) \geq \operatorname{Ric}_{M}\left(\nabla_{M} u, \nabla_{M} u\right) \geq \frac{1}{4} \Phi_{j}^{i k} \Phi_{k}^{j \ell} u_{i} u_{\ell}
$$

and the third proof of Lemma 3.1 is complete.
Having finished with Lemma 3.1, let us introduce one of the main ideas in this paper, which was absent from [14]. The idea is to consider the map

$$
\begin{equation*}
\mathbb{R}^{n} \ni x \mapsto D^{2} \Phi(x) \in M_{n}^{+}(\mathbb{R}) \tag{32}
\end{equation*}
$$

Denote by $\left(g_{i j}(x)\right)_{i, j=1, \ldots, n}$ the pull-back of the standard Riemannian metric on $M_{n}^{+}(\mathbb{R})$ via the map (32). It follows from Lemma 2.2 that $g_{i j}$ is given by the formula

$$
\begin{equation*}
g_{i j}=\operatorname{Tr}\left[\left(D^{2} \Phi\right)^{-1} \cdot \partial_{i}\left(D^{2} \Phi\right) \cdot\left(D^{2} \Phi\right)^{-1} \cdot \partial_{j}\left(D^{2} \Phi\right)\right]=\Phi_{i k}^{\ell} \Phi_{j \ell}^{k} \tag{33}
\end{equation*}
$$

Note that the positive semi-definite matrix $\left(g_{i j}(x)\right)_{i, j=1, \ldots, n}$ is not necessarily invertible, and it could happen that distinct points of $\mathbb{R}^{n}$ have zero Riemannian distance with respect to the Riemannian metric $\left(g_{i j}\right)$. The metric $g_{i j}$ resembles an expression appearing in Lemma 3.1, a fact that will be exploited in the next section.

## 4 Dualizing the Bochner inequality

It is by now well-known that in the presence of convexity assumptions, Poincaré-type inequalities may be deduced from Bochner's formula via a dualization procedure. In this section we investigate the Poincaré inequality that is dual to Lemma 3.1. This Poincaré inequality was also obtained in [14], but in a cumbersome formulation and under an undesired assumption called "regularity at infinite", which we eliminate here.

We begin with an easy case. Throughout this section we assume, in addition to the smoothness assumptions made in the beginning of Section 3, that there exists $\varepsilon_{0}>0$ for which

$$
\begin{equation*}
D^{2} \Phi(x) \geq \varepsilon_{0} \cdot \operatorname{Id} \quad\left(x \in \mathbb{R}^{n}\right) \tag{34}
\end{equation*}
$$

in the sense of symmetric matrices. Write $C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ for the space of all compactly-supported, smooth functions on $\mathbb{R}^{n}$. The following lemma is a variant of a well-known fact (see, e.g., Strichartz [23]), that compactly-supported functions are dense in Sobolev spaces when the Riemannian manifold is complete. Our assumption (34) implies the completeness of the Riemannian manifold $M=M_{\mu, \nu}$.

Lemma 4.1. Let $f \in L^{2}(\mu)$ satisfy $\int f d \mu=0$. Then there exists a sequence $u_{k} \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ with

$$
\left\|L u_{k}-f\right\|_{L^{2}(\mu)} \xrightarrow{k \rightarrow \infty} 0 .
$$

Proof. Recall that $\int(L u) d \mu=0$ for all $u \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$. In order to show that the linear space $\left\{L u ; u \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)\right\}$ is dense, we analyze its orthogonal complement. Let $f \in L^{2}(\mu)$ be in the orthogonal complement, i.e., for any $u \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$,

$$
\begin{equation*}
\int_{\mathbb{R}^{n}} f(L u) d \mu=0 \tag{35}
\end{equation*}
$$

Our goal is to show that $f \equiv$ Const. Note that (35) means that $f$ is a weak solution of $L f \equiv 0$. Since $L$ is elliptic, then $f$ is smooth and $L f \equiv 0$ in the classical sense. Thus,

$$
L\left(f^{2}\right)=2 f L f+2\left|\nabla_{M} f\right|^{2}=2\left|\nabla_{M} f\right|^{2}
$$

Therefore, for any $\eta \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$,

$$
\begin{aligned}
\int_{\mathbb{R}^{n}}\left|\nabla_{M}(\eta f)\right|^{2} d \mu & =\int_{\mathbb{R}^{n}}\left[\eta^{2}\left|\nabla_{M} f\right|^{2}+\frac{1}{2} \nabla_{M}\left(f^{2}\right) \cdot \nabla_{M}\left(\eta^{2}\right)+f^{2}\left|\nabla_{M} \eta\right|^{2}\right] d \mu \\
& =\int_{\mathbb{R}^{n}}\left[\eta^{2}\left|\nabla_{M} f\right|^{2}-\frac{1}{2} \eta^{2} L\left(f^{2}\right)+f^{2}\left|\nabla_{M} \eta\right|^{2}\right] d \mu=\int_{\mathbb{R}^{n}}\left|\nabla_{M} \eta\right|^{2} f^{2} d \mu
\end{aligned}
$$

However, according to our assumption (34), we have $\left|\nabla_{M} \eta\right|^{2}=\Phi^{i j} \eta_{i} \eta_{j} \leq \varepsilon_{0}^{-1}|\nabla \eta|^{2}$. Let $\eta_{R}$ be a smooth cutoff function in $\mathbb{R}^{n}$, that equals one on a Euclidean ball of radius $R$ centered at the origin, that equals zero outside a Euclidean ball of radius $2 R$, and that satisfies $\left|\nabla \eta_{R}\right| \leq 2 / R$ throughout $\mathbb{R}^{n}$. Then,
$\int_{K}\left|\nabla_{M}\left(\eta_{R} f\right)\right|^{2} d \mu \leq \int_{\mathbb{R}^{n}}\left|\nabla_{M} \eta\right|^{2} f^{2} d \mu \leq \varepsilon_{0}^{-1} \int_{\mathbb{R}^{n}}\left|\nabla \eta_{R}\right|^{2} f^{2} d \mu \leq \frac{2}{R \varepsilon_{0}} \int_{\mathbb{R}^{n}} f^{2} d \mu \xrightarrow{R \rightarrow \infty} 0$,
since $f \in L^{2}(\mu)$. Therefore $\nabla f \equiv 0$ and $f$ is constant.

Suppose that $F$ is a locally-Lipschitz function on a Riemannian manifold such as $M_{n}^{+}(\mathbb{R})$. By the Rademacher theorem, the gradient $\nabla F$ is well-defined almost everywhere with respect to the Riemannian volume measure. In order to have a function $|\nabla F|$ that is defined everywhere, in this paper we set

$$
\begin{equation*}
|\nabla F|(x)=\limsup _{\substack{y \rightarrow x \\ z \rightarrow x}} \frac{|f(y)-f(z)|}{\operatorname{dist}(y, z)}=\lim _{\varepsilon \rightarrow 0^{+}} \sup _{\substack{y, z \in B(x, \varepsilon) \\ y \neq z}} \frac{|f(y)-f(z)|}{\operatorname{dist}(y, z)} \tag{36}
\end{equation*}
$$

where dist is the Riemannian distance, and $B(x, \varepsilon)=\{y ; \operatorname{dist}(x, y)<\varepsilon\}$. Since $F$ is locally-Lipschitz, then the function $|\nabla F|$ is locally-bounded and upper semi-continuous. Clearly, at any point $x$ where $F$ is continuously differentiable, $|\nabla F|(x)$ equals the Riemannian length of $\nabla F(x)$.

Proposition 4.2. Denote by $\theta$ the push-forward of the measure $\mu$ under the map (32). Then for any locally-Lipschitz function $F: M_{n}^{+}(\mathbb{R}) \rightarrow \mathbb{R}$ that belongs to $L^{2}(\theta)$ with $\int_{M_{n}^{+}(\mathbb{R})} F d \theta=0$,

$$
\int_{M_{n}^{+}(\mathbb{R})} F^{2} d \theta \leq 4 \int_{M_{n}^{+}(\mathbb{R})}|\nabla F|^{2} d \theta,
$$

whenever the right-hand side is finite.
Proof. Since $F$ is locally-Lipschitz in $L^{2}(\theta)$, then the function $f$ defined via

$$
f(x)=F\left(D^{2} \Phi(x)\right) \quad\left(x \in \mathbb{R}^{n}\right),
$$

is locally-Lipschitz in $\mathbb{R}^{n}$ and belongs to $L^{2}(\mu)$. Abbreviate $H=|\nabla F|^{2}$ and $h(x)=$ $H\left(D^{2} \Phi(x)\right)$. From the definition (36) of $|\nabla F|$, for any $x \in \mathbb{R}^{n}$ in which $f$ is differentiable,

$$
\begin{equation*}
h(x) \geq \sup \left\{\sum_{i=1}^{n} V^{i} f_{i} ; \sum_{i, j=1}^{n} g_{i j} V^{i} V^{j} \leq 1, V^{1}, \ldots, V^{n} \in \mathbb{R}\right\} \tag{37}
\end{equation*}
$$

where $f_{i}$ and $g_{i j}$ are evaluated at the point $x$. In the case where the matrix $\left(g_{i j}(x)\right)_{i, j=1, \ldots, n}$ is invertible, we may express the supremum in (37) in terms of the inverse matrix, yet it is the formula (37) which is valid in the general case. Setting $U_{i}=\Phi_{i j} V^{j}$, we reformulate (37) as

$$
\begin{equation*}
h(x) \geq \sup \left\{\Phi^{i j} U_{j} f_{i} ; g_{i j} \Phi^{k i} \Phi^{\ell j} U_{k} U_{\ell} \leq 1, U_{1}, \ldots, U_{n} \in \mathbb{R}\right\} \tag{38}
\end{equation*}
$$

The formula (38) is valid for almost any $x \in \mathbb{R}^{n}$, since $f$ is differentiable almost everywhere in $\mathbb{R}^{n}$ by the Rademacher theorem. We would like to show that for any $u \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$,

$$
\begin{equation*}
-\int_{\mathbb{R}^{n}} f(L u) d \mu \leq 2 \sqrt{\int_{\mathbb{R}^{n}} h^{2} d \mu} \cdot \sqrt{\int_{\mathbb{R}^{n}}(L u)^{2} d \mu} \tag{39}
\end{equation*}
$$

To that end, we observe that since $u$ is compactly-supported,
$\int_{\mathbb{R}^{n}} \Gamma_{2}(u) d \mu=\frac{1}{2} \int_{\mathbb{R}^{n}} L\left(\Phi^{i j} u_{i} u_{j}\right) d \mu-\int_{\mathbb{R}^{n}} \Phi^{i j}(L u)_{i} u_{j} d \mu=-\int_{\mathbb{R}^{n}} \Phi^{i j}(L u)_{i} u_{j} d \mu=\int_{\mathbb{R}^{n}}(L u)^{2} d \mu$.
Therefore Lemma 3.1 and (33) imply that for any $u \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$,

$$
\int_{\mathbb{R}^{n}}(L u)^{2} d \mu \geq \frac{1}{4} \int_{\mathbb{R}^{n}} \Phi^{i k} \Phi^{j \ell} g_{k \ell} u_{i} u_{j} d \mu
$$

Since $f$ is locally-Lipschitz, we may safely integrate by parts, and obtain that for any $u \in$ $C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$,

$$
\begin{aligned}
& -\int_{\mathbb{R}^{n}} f(L u) d \mu=\int_{\mathbb{R}^{n}} \Phi^{i j} f_{i} u_{j} d \mu \leq \int_{\mathbb{R}^{n}} h(x) \sqrt{g_{i j} \Phi^{k i} \Phi^{\ell j} u_{k} u_{\ell}} d \mu(x) \\
& \quad \leq \sqrt{\int_{\mathbb{R}^{n}} h^{2} d \mu} \sqrt{\int_{\mathbb{R}^{n}} g_{i j} \Phi^{k i} \Phi^{\ell j} u_{k} u_{\ell} d \mu} \leq \sqrt{\int_{\mathbb{R}^{n}} h^{2} d \mu} \sqrt{\int_{\mathbb{R}^{n}}(L u)^{2} d \mu}
\end{aligned}
$$

and (39) is proven. Since $\int_{M_{n}^{+}(\mathbb{R})} F d \theta=0$ then also $\int_{\mathbb{R}^{n}} f d \mu=0$. From Lemma 4.1 there exists a sequence $u_{k} \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ with $L u_{k} \rightarrow-f$ in $L^{2}(\mu)$. We plug in $u=u_{k}$ in (39), and take the limit $k \rightarrow \infty$. This yields

$$
\int_{\mathbb{R}^{n}} f^{2} d \mu \leq 2 \sqrt{\int_{\mathbb{R}^{n}} h^{2} d \mu} \cdot \sqrt{\int_{\mathbb{R}^{n}} f^{2} d \mu}
$$

Hence,

$$
\int_{\mathbb{R}^{n}} f^{2} d \mu \leq 4 \int_{K} h^{2} d \mu
$$

Since $h(x)=H\left(D^{2} \Phi\right)$ with $H=|\nabla F|^{2}$, the proposition is proven.

## 5 Regularity issues

This section explains how to eliminate assumption (34) and also the smoothness assumptions of the previous two sections.

Theorem 5.1. Assume that $\mu$ and $\nu$ are absolutely-continuous, log-concave probability measures on $\mathbb{R}^{n}$. Let $\nabla \Phi$ be the Brenier map between $\mu$ and $\nu$, and assume condition ( $\star$ ) from Section 1. Denote by $\theta$ the push-forward of the measure $\mu$ under the map $x \mapsto D^{2} \Phi(x)$.

Then for any $\theta$-integrable, locally-Lipschitz function $F: M_{n}^{+}(\mathbb{R}) \rightarrow \mathbb{R}$,

$$
\begin{equation*}
\int_{M_{n}^{+}(\mathbb{R})} F^{2} d \theta-\left(\int_{M_{n}^{+}(\mathbb{R})} F d \theta\right)^{2} \leq 4 \int_{M_{n}^{+}(\mathbb{R})}|\nabla F|^{2} d \theta \tag{40}
\end{equation*}
$$

whenever the right-hand side is finite, and $|\nabla F|$ is interpreted as in (36).
The strategy for proving Theorem 5.1 is to approximate $\Phi$ by a sequence of functions $\Phi_{N}$ that satisfy assumption (34), and to prove the pointwise (even local uniform) convergence $D^{2} \Phi_{N}(x) \xrightarrow{N \rightarrow \infty} D^{2} \Phi(x)$. We discuss below two possible justifications of this convergence, because we believe that both of them may be useful. The first proof occupies Subsection 5.1, and it is based on various results from the regularity theory of the Monge-Ampère equation. The log-concavity of the measures is not really required for the first proof, and it suffices to assume that the densities are locally Hölder.

The second proof in Subsection 5.2 is in fact an alternative approach to Caffareli's $C^{1, \alpha_{-}}$ regularity results in the log-concave case. The argument in Subsection 5.2 is more selfcontained, and it is based on integration-by-parts arguments. The log-concavity of the target measure plays an important role here, and we further assume a certain integrability condition on the logarithmic derivative of the density of $\mu$. This integrability condition is rather mild in our opinion, and it is satisfied in many cases of interest.

### 5.1 First proof of Theorem 5.1

As before, we write $e^{-V}$ and $e^{-W}$ for the densities of $\mu$ and $\nu$, respectively. By logconcavity, the functions $V$ and $W$ are locally-Lipschitz in the open sets $\operatorname{Supp}(\mu)$ and $\operatorname{Supp}(\nu)$, respectively. From condition $(\star)$ the function $\Phi$ is $C^{2}$-smooth, and the push-forward equation (1) implies that

$$
\begin{equation*}
\operatorname{det} D^{2} \Phi(x)=e^{-V(x)+W(\nabla \Phi(x))} \tag{41}
\end{equation*}
$$

for any $x \in \operatorname{Supp}(\mu)$. In particular, $D^{2} \Phi(x)$ is invertible, and hence positive-definite for all $x \in \operatorname{Supp}(\mu)$. Thus $\Phi$ is strictly-convex. The modulus of convexity of $\Phi$ at the point $x$ is defined to be

$$
\omega_{\Phi}(x ; \delta)=\inf \left\{\Phi(y)-(\Phi(x)+\nabla \Phi(x) \cdot(y-x)) ; y \in \mathbb{R}^{n},|y-x|=\delta\right\}
$$

Then $\omega_{\Phi}(x ; \delta)$ is a positive, continuous function of $x \in \operatorname{Supp}(\mu)$ and $\delta>0$, when we restrict attention to $x$ and $\delta$ for which $\overline{B(x, \delta)} \subseteq \operatorname{Supp}(\mu)$. Here, $B(x, \delta)=\left\{y \in \mathbb{R}^{n} ;|y-x|<\delta\right\}$. Next, the Legendre transform

$$
\Phi^{*}(x)=\sup _{\substack{y \in \mathbb{R}^{n} \\ \Phi(y)<\infty}}[x \cdot y-\Phi(y)]
$$

is also $C^{2}$-smooth and strictly-convex in $\operatorname{Supp}(\nu)$, with $y \mapsto \nabla \Phi^{*}(y)$ being the inverse map to $x \mapsto \nabla \Phi(x)$. Thus $\nabla \Phi$ is a $C^{1}$-diffeomorphism of $\operatorname{Supp}(\mu)$ and $\operatorname{Supp}(\nu)$. The reader is referred to Rockafellar [22] for the basic properties of the Legendre transform.

We will approximate $\mu$ and $\nu$ by sequences of probability measures $\mu_{N}$ and $\nu_{N}$ with the following properties:
(i) The probability measure $\mu_{N}$ (respectively $\nu_{N}$ ) has a density in $\mathbb{R}^{n}$ of the form $e^{-V_{N}}$ (respectively $e^{-W_{N}}$ ).
(ii) The functions $V_{N}, W_{N}: \mathbb{R}^{n} \rightarrow \mathbb{R}$ are smooth and for any $x \in \mathbb{R}^{n}$,

$$
D^{2} V_{N}(x) \geq \frac{1}{N} \cdot \operatorname{Id}, \quad D^{2} W_{N}(x) \leq N \cdot \mathrm{Id}
$$

(iii) $V_{N} \longrightarrow V$ locally uniformly in $\operatorname{Supp}(\mu)$, and similarly, $W_{N} \longrightarrow W$ locally uniformly in $\operatorname{Supp}(\nu)$.
It is quite standard to approximate $\mu$ and $\nu$ in this manner. For instance, in order to obtain $\mu_{N}\left(\right.$ or $\left.\nu_{N}\right)$, we may convolve $\mu$ (or $\nu$ ) with a Gaussian of a tiny variance, then multiply the resulting density by a Gaussian of a huge variance, and then normalize to obtain a probability density. Denote by $\nabla \Phi_{N}$ the Brenier map between $\mu_{N}$ and $\nu_{N}$. We use again Caffarelli's regularity theory, to conclude that $\Phi_{N}: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is a smooth, strictly-convex function, with

$$
\begin{equation*}
\operatorname{det} D^{2} \Phi_{N}(x)=e^{-V_{N}(x)+W_{N}\left(\nabla \Phi_{N}(x)\right)} \quad\left(x \in \mathbb{R}^{n}\right) \tag{42}
\end{equation*}
$$

The following lemma should be known to experts on the Monge-Ampère equation, yet we could not find it in the literature.

Lemma 5.2. There exists an increasing sequence $\left\{N_{j}\right\}$ such that

$$
D^{2} \Phi_{N_{j}}(x) \xrightarrow{j \rightarrow \infty} D^{2} \Phi(x)
$$

locally uniformly in $x \in \operatorname{Supp}(\mu)$.
Proof. Fix $x_{0} \in \operatorname{Supp}(\mu)$. It suffices to find $\left\{N_{j}\right\}$ such that $D^{2} \Phi_{N_{j}} \longrightarrow D^{2} \Phi$ uniformly in a neighborhood of $x_{0}$. A standard convexity argument (e.g., [15, Section 2]) based on (iii) and the fact that $\int e^{-V}=\int e^{-W}=1$ shows that there exist $A, B>0$ with

$$
\begin{equation*}
\min \left\{\inf _{N} V_{N}(x), \inf _{N} W_{N}(x), V(x), W(x)\right\} \geq A|x|-B, \quad\left(x \in \mathbb{R}^{n}\right) \tag{43}
\end{equation*}
$$

Therefore,

$$
\begin{equation*}
\sup _{N} \int_{\mathbb{R}^{n}}\left|\nabla \Phi_{N}\right|^{2} e^{-V_{N}(x)} d x=\sup _{N} \int_{\mathbb{R}^{n}}|x|^{2} e^{-W_{N}(x)} d x \leq \int_{\mathbb{R}^{n}}|x|^{2} e^{B-A|x|} d x<\infty \tag{44}
\end{equation*}
$$

Recall that $V_{N} \rightarrow V$ locally uniformly in $\operatorname{Supp}(\mu)$, according to (iii). From (44) we learn that $\sup _{N}\left\|\Phi_{N}\right\|_{\dot{H}^{1}(K)}<\infty$ for any compact $K \subset \operatorname{Supp}(\mu)$. Here,

$$
\|u\|_{\dot{H}^{1}(K)}^{2}=\int_{K}|\nabla u(x)|^{2} d x
$$

From the Rellich-Kondrachov compactness theorem (e.g., [8, Section 4.6]), we conclude that there exists a subsequence $\Phi_{N_{j}}$, numbers $C_{j} \in \mathbb{R}$ and a certain function $F: \operatorname{Supp}(\mu) \rightarrow \mathbb{R}$ such that for any compact $K \subset S u p p(\mu)$, the sequence $\Phi_{N_{j}}+C_{j}$ converges to $F$ in $L^{2}(K)$.

Passing to another subsequence, which we conveniently denote again by $\left\{\Phi_{N}\right\}$, and using [22, Theorem 10.9], we may assume that $F$ is convex and that the convergence is locallyuniform in $\operatorname{Supp}(\mu)$. Thus, from [22, Theorem 24.5],

$$
\begin{equation*}
\nabla \Phi_{N}(x) \xrightarrow{N \rightarrow \infty} \nabla F(x) \tag{45}
\end{equation*}
$$

for almost any $x \in \operatorname{Supp}(\mu)$. However, $\left(\nabla \Phi_{N}\right)_{*} \mu_{N}=\nu_{N}$. From (iii), (43) and (45) we conclude that $(\nabla F)_{*} \mu=\nu$. From the uniqueness of the Brenier map, we deduce that $\nabla F=\nabla \Phi$ almost everywhere in $\operatorname{Supp}(\mu)$. Since $\Phi$ is $C^{2}$-smooth, then we may apply [22, Theorem 25.7], and upgrade (45) to

$$
\begin{equation*}
\nabla \Phi_{N}(x) \xrightarrow{N \rightarrow \infty} \nabla \Phi(x) \tag{46}
\end{equation*}
$$

locally uniformly in $\operatorname{Supp}(\mu)$. The convexity arguments in [22, Section 25] also show that $\nabla \Phi_{N}^{*} \rightarrow \nabla \Phi^{*}$ locally uniformly in $\operatorname{Supp}(\nu)$. As for the modulus of convexity, we have

$$
\begin{equation*}
\omega_{\Phi_{N}}(x ; \delta) \xrightarrow{N \rightarrow \infty} \omega_{\Phi}(x ; \delta), \quad \text { and respectively, } \quad \omega_{\Phi_{N}^{*}}(y ; \delta) \xrightarrow{N \rightarrow \infty} \omega_{\Phi^{*}}(y ; \delta) \tag{47}
\end{equation*}
$$

locally uniformly in the set $\{(x, \delta) \in \operatorname{Supp}(\mu) \times(0, \infty) ; \overline{B(x, \delta)} \subset \operatorname{Supp}(\mu)\}$, and respectively, in the set $\{(y, \delta) \in \operatorname{Supp}(\nu) \times(0, \infty) ; \overline{B(y, \delta)} \subset \operatorname{Supp}(\nu)\}$.

We will now invoke the estimates of Gutierrez and Huang [12] and Forzani and Maldonado [9, 10], which are constructive versions of Caffarelli's $C^{1, \alpha}$-regularity theory. We are allowed to apply [12, Theorem 2.1] and [9, Theorem 15] locally near $x_{0}$, thanks to (iii), (42), (46) and (47). From [9, Theorem 15] we learn that there exist $\alpha, \delta, C>0$ such that for any $x, y \in B\left(x_{0}, \delta\right)$ and $N \geq 1$,

$$
\begin{equation*}
\left|\nabla \Phi_{N}(x)-\nabla \Phi_{N}(y)\right| \leq C|x-y|^{\alpha} \tag{48}
\end{equation*}
$$

The function $V$ is locally-Lipschitz. From (iii) and [22, Theorem 24.5], the sequence $\left\{V_{N}\right\}$ is uniformly locally-Lipschitz: This means that for any compact subset $K \subset \operatorname{Supp}(\mu)$, the Lipschitz constant of $V_{N}$ is bounded by some finite number $C_{K}$, independent of $N$. Similarly, the sequence $\left\{W_{N}\right\}$ is also uniformly locally-Lipschitz. Together with (46) and (48) we deduce that there exist $\hat{C}>0$ such that $u_{N}(x)=-V_{N}(x)+W_{N}\left(\nabla \Phi_{N}(x)\right)$ satisfies

$$
\left|u_{N}(x)-u_{N}(y)\right| \leq \hat{C}|x-y|^{\alpha} \quad\left(x, y \in B\left(x_{0}, \delta\right), N \geq 1\right)
$$

Recalling the Monge-Ampère equation (42), we learn that that there exists $\tilde{C}>0$ such that

$$
\left|\operatorname{det} D^{2} \Phi_{N}(x)-\operatorname{det} D^{2} \Phi_{N}(y)\right| \leq \tilde{C}|x-y|^{\alpha} . \quad\left(x, y \in B\left(x_{0}, \delta\right), N \geq 1\right)
$$

We are finally in good shape for applying the $C^{2, \alpha}$-estimates from Trudinger and Wang [24, Theorem 3.2]. These estimates yield the existence of $\bar{C}>0$ such that for any $x, y \in$ $B\left(x_{0}, \delta / 2\right)$ and $N \geq 1$,

$$
\begin{equation*}
\left\|D^{2} \Phi_{N}(x)-D^{2} \Phi_{N}(y)\right\|_{H S} \leq \bar{C}|x-y|^{\alpha} \tag{49}
\end{equation*}
$$

The uniform $C^{2, \alpha}$-estimate in (49) allows us to apply the Arzella-Ascoli theorem. All we need is to denote $K=B\left(x_{0}, \delta / 2\right)$ and observe that

$$
\int_{K}\left(\Delta \Phi_{N}\right) \xi=-\int_{K} \nabla \Phi_{N} \cdot \nabla \xi \xrightarrow{N \rightarrow \infty}-\int_{K} \nabla \Phi \cdot \nabla \xi=\int_{K}(\Delta \Phi) \xi,
$$

where $\xi$ is any smooth, compactly-supported function in $K$. Hence the sequence $\left\{\int_{K} \Delta \Phi_{N}\right\}_{N \geq 1}$ is bounded, and since $D^{2} \Phi_{N}$ is positive-definite, also the sequence $\left\{\int_{K}\left\|D^{2} \Phi_{N}\right\|_{H S}\right\}_{N \geq 1}$ is bounded. From (49) and the Arzella-Ascoli theorem, there exists a subsequence, denoted still by $\left\{\Phi_{N}\right\}$, such that $D^{2} \Phi_{N} \longrightarrow D^{2} \Phi$ uniformly on $K=B\left(x_{0}, \delta / 2\right)$.

Remark 5.3. Our proof of Lemma 5.2 does not make any use of the log-concavity of $\mu$ and $\nu$. By inspecting the proof above, we see that Lemma 5.2 holds true as long as $V$ and $W$ are locally Hölder, and $V_{N}, W_{N}$ are uniformly locally Hölder.

In order to simplify the notation, we denote the sequence $\left\{\Phi_{N_{j}}\right\}$ from Lemma 5.2 by $\left\{\Phi_{N}\right\}$. Properties (i), (ii) and (iii) above are still satisfied.

Corollary 5.4. Denote by $\theta_{N}$ the push-forward of the measure $\mu_{N}$ under the map $x \mapsto$ $D^{2} \Phi_{N}(x)$. Then for any bounded, continuous function $b: M_{n}^{+}(\mathbb{R}) \rightarrow \mathbb{R}$,

$$
\begin{equation*}
\int_{M_{n}^{+}(\mathbb{R})} b d \theta_{N} \xrightarrow{N \rightarrow \infty} \int_{M_{n}^{+}(\mathbb{R})} b d \theta . \tag{50}
\end{equation*}
$$

Furthermore, if $b: M_{n}^{+}(\mathbb{R}) \rightarrow \mathbb{R}$ is bounded and upper semi-continuous, then

$$
\begin{equation*}
\limsup _{N \rightarrow \infty} \int_{M_{n}^{+}(\mathbb{R})} b d \theta_{N} \leq \int_{M_{n}^{+}(\mathbb{R})} b d \theta \tag{51}
\end{equation*}
$$

Proof. In order to prove (50), we need to show that

$$
\int_{\mathbb{R}^{n}} b\left(D^{2} \Phi_{N}(x)\right) e^{-V_{N}(x)} d x \xrightarrow{N \rightarrow \infty} \int_{\mathbb{R}^{n}} b\left(D^{2} \Phi(x)\right) e^{-V(x)} d x .
$$

This follows from Lemma 5.2 and the dominated convergence theorem, since (43) provides an integrable majorant. Next, assume that $b$ is bounded and upper semi-continuous. Then for any $x \in \operatorname{Supp}(\mu)$,

$$
\limsup _{N \rightarrow \infty} b\left(D^{2} \Phi_{N}(x)\right) e^{-V_{N}(x)} \leq b\left(D^{2} \Phi(x)\right) e^{-V(x)}
$$

Now (51) follows from Fatou's lemma, since we have an integrable majorant by (43).

Proof of Theorem 5.1. Assume first that the locally-Lipschitz function $F$ is compactly supported. We observe that for any fixed $N$, assumption (34) holds true. Indeed, we may apply a refinement of Caffarelli's contraction theorem [7] which appears in [18], and obtain from (ii) that for any $x \in \mathbb{R}^{n}$,

$$
D^{2} \Phi_{N}(x) \geq \frac{1}{N^{2}} \cdot \text { Id }
$$

We may therefore apply Proposition 4.2, and conclude that for any $N \geq 1$,

$$
\int_{M_{n}^{+}(\mathbb{R})} F^{2} d \theta_{N}-\left(\int_{M_{n}^{+}(\mathbb{R})} F d \theta_{N}\right)^{2} \leq 4 \int_{M_{n}^{+}(\mathbb{R})}|\nabla F|^{2} d \theta_{N}
$$

Recall that $|\nabla F|^{2}$ is upper semi-continuous and bounded, while $F$ is continuous and bounded. By taking the limit as $N \rightarrow \infty$ and using Corollary 5.4 , we obtain that

$$
\int_{M_{n}^{+}(\mathbb{R})} F^{2} d \theta-\left(\int_{M_{n}^{+}(\mathbb{R})} F d \theta\right)^{2} \leq 4 \int_{M_{n}^{+}(\mathbb{R})}|\nabla F|^{2} d \theta
$$

and (40) is proven in the case where $F$ is a compactly-supported function.
The next step is to prove (40) under the additional assumption that $F \in L^{2}(\theta)$. To that end we pick a smooth function $\theta_{R}: M_{n}^{+}(\mathbb{R}) \rightarrow[0,1]$, such that $\theta_{R}$ equals one on $B(\operatorname{Id}, R)$ and it vanishes outside $B(\operatorname{Id}, 2 R)$, with $\left|\nabla \theta_{R}\right| \leq 2 / R$. Set $F_{R}=\theta_{R} F$. We have just proven that (40) holds true when $F$ is replaced by $F_{R}$. Clearly, $F_{R} \longrightarrow F$ in $L^{2}(\theta)$ as $R \longrightarrow \infty$. All that remains is to show that

$$
\begin{equation*}
\limsup _{R \rightarrow \infty} \int_{M_{n}^{+}(\mathbb{R})}\left|\nabla F_{R}\right|^{2} d \theta \leq \int_{M_{n}^{+}(\mathbb{R})}|\nabla F|^{2} d \theta \tag{52}
\end{equation*}
$$

The functions $\theta_{R}$ and $F$ are continuous, and therefore we may use the Leibnitz rule

$$
\left|\nabla F_{R}\right| \leq|F|\left|\nabla \theta_{R}\right|+\theta_{R}|\nabla F| \leq|\nabla F|+2|F| / R,
$$

where we interpret $|\nabla F|$ and $\left|\nabla F_{R}\right|$ in the sense of definition (36). Since $F,|\nabla F| \in L^{2}(\theta)$, then (52) follows in the case where $F \in L^{2}(\theta)$.

Finally, in order to eliminate the assumption that $F \in L^{2}(\theta)$, we replace $F$ by $F_{R}=$ $\max \{-R, \min \{F, R\}\}$, apply the inequality for $F_{R}$, and let $R$ tend to infinity. For all but countably many values of $R$, the level set $\left\{A \in M_{n}^{+}(\mathbb{R}) ; F(A)=R\right\}$ has zero $\theta$-measure. Consequently, we have the inequality $\int\left|\nabla F_{R}\right|^{2} d \theta \leq \int|\nabla F|^{2} d \theta$ for all but countably many values of $R$, and (40) follows.

### 5.2 Second proof: Log-concave target measure

In our second proof we will exploit the fact that $\nu$ is log-concave, but we will not require the log-concavity of $\mu$. Throughout this subsection we make the following additional assumption:

Assumption (A): For some $p>n$,

$$
\int_{\mathbb{R}^{n}}|\nabla V|^{p} e^{-V} d x<\infty
$$

where the derivatives $V_{i}$ are understood in the logarithmic derivative sense, i.e.

$$
\int_{\mathbb{R}^{n}} \xi V_{i} d \mu=-\int_{\mathbb{R}^{n}} \xi_{i} d \mu, \quad \xi \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right), i=1, \ldots, n
$$

By the Morrey embedding theorem (see, e.g., [8, Section 4.5]), the function $V$ is locally Hölder. We will approximate $\mu$ and $\nu$ by sequences of probability measures $\mu_{N}$ and $\nu_{N}$ having properties (i), (ii) and (iii) from Subsection 5.1. We also require a fourth property:
(iv) There exists $p>n$ such that

$$
\sup _{N} \int_{\mathbb{R}^{n}}\left|\nabla V_{N}\right|^{p} e^{-V_{N}} d x<\infty .
$$

The approach outlined in Subsection 5.1, to convolve with a tiny Gaussian and then multiply by the density of a huge Gaussian, yields also property (iv). Recall that the Brenier map $\nabla \Phi_{N}$ between $\mu_{N}$ and $\nu_{N}$ is smooth and that it satisfies (42). The central ingredient of this subsection is the following a priori estimate:

Proposition 5.5. Assume that functions $V, W$ and $\Phi$ are smooth on the entire $\mathbb{R}^{n}$ and that $\nu$ is a log-concave measure. Then for every $q \geq 2,0<\tau<1, i=1, \ldots, n$ there exists $C(q, \tau)>0$

$$
\begin{equation*}
\int_{\mathbb{R}^{n}} \Phi_{i i}^{q} d \mu \leq C(q, \tau)\left(\int_{\mathbb{R}^{n}}\left|V_{i}\right|^{\frac{2 q}{2-\tau}} d \mu+\int_{\mathbb{R}^{n}}\left|x_{i}\right|^{\frac{2 q}{\tau}} d \nu\right) . \tag{53}
\end{equation*}
$$

Proof. Assume in addition that $D^{2} W \geq \frac{1}{C} \cdot \operatorname{Id}, D^{2} V \leq C \cdot$ Id. In this case $D^{2} \Phi \leq C^{2} \cdot$ Id. Recall formula (28),

$$
L\left(\Phi_{i i}\right)-\Phi_{i}^{j k} \Phi_{i j k}-\sum_{j, k=1}^{n} \Phi_{j i} \Phi_{i k} W_{j k} \circ \nabla \Phi=-V_{i i},
$$

which is obtained by differentiating the change of variables formula (22) along $x_{i}$. Let us multiply this formula by $\Phi_{i i}^{p}, p \geq 0$ and make a formal integration by parts with respect to $\mu$. Using the convexity of $W$ we get

$$
\begin{equation*}
\int V_{i i} \Phi_{i i}^{p} d \mu \geq p \int \Phi_{i i}^{p-1}\left\langle\left(D^{2} \Phi\right)^{-1} \nabla \Phi_{i i}, \nabla \Phi_{i i}\right\rangle d \mu+\int \Phi_{i i}^{p} \Phi_{i}^{j k} \Phi_{i j k} d \mu . \tag{54}
\end{equation*}
$$

Let us justify this formula. To this end we fix a compactly supported function $\eta \geq 0$ and integrate with respect to $\eta \cdot \mu$.
$\int V_{i i} \Phi_{i i}^{p} \eta d \mu \geq \int\left\langle\left(D^{2} \Phi\right)^{-1} \nabla \eta, \nabla \Phi_{i i}\right\rangle \Phi_{i i}^{p} d \mu+p \int \Phi_{i i}^{p-1}\left\langle\left(D^{2} \Phi\right)^{-1} \nabla \Phi_{i i}, \nabla \Phi_{i i}\right\rangle \eta d \mu+\int \Phi_{i i}^{p} \Phi_{i}^{j k} \Phi_{i j k} \eta d \mu$.
Applying the Cauchy inequality we get
$-\int\left\langle\left(D^{2} \Phi\right)^{-1} \nabla \eta, \nabla \Phi_{i i}\right\rangle \Phi_{i i}^{p} d \mu \leq \frac{4}{\varepsilon} \int \frac{\left\langle\left(D^{2} \Phi\right)^{-1} \nabla \eta, \nabla \eta\right\rangle}{\eta} \Phi_{i i}^{p+1} d \mu+\varepsilon \int\left\langle\left(D^{2} \Phi\right)^{-1} \nabla \Phi_{i i}, \nabla \Phi_{i i}\right\rangle \Phi_{i i}^{p-1} \eta d \mu$.
Finally,

$$
\begin{aligned}
\int V_{i i} \Phi_{i i}^{p} \eta d \mu & +\frac{4}{\varepsilon} \int \frac{\left\langle\left(D^{2} \Phi\right)^{-1} \nabla \eta, \nabla \eta\right\rangle}{\eta} \Phi_{i i}^{p+1} d \mu \\
& \geq(p-\varepsilon) \int \Phi_{i i}^{p-1}\left\langle\left(D^{2} \Phi\right)^{-1} \nabla \Phi_{i i}, \nabla \Phi_{i i}\right\rangle \eta d \mu+\int \Phi_{i i}^{p} \Phi_{i}^{j k} \Phi_{i j k} \eta d \mu .
\end{aligned}
$$

Assume that $\eta$ has the form $\eta=\xi(\nabla \Phi)$, where $\xi$ is compactly supported. We get

$$
\int V_{i i} \Phi_{i i}^{p} \eta d \mu+\frac{4 C^{p+2}}{\varepsilon} \int \frac{|\nabla \xi|^{2}}{\xi} d \nu \geq(p-\varepsilon) \int \Phi_{i i}^{p-1}\left\langle\left(D^{2} \Phi\right)^{-1} \nabla \Phi_{i i}, \nabla \Phi_{i i}\right\rangle \eta d \mu+\int \Phi_{i i}^{p} \Phi_{i}^{j k} \Phi_{i j k} \eta d \mu
$$

It remains to construct a sequence of functions $1 \geq \xi_{N} \geq 0$ satisfying $\lim _{N} \xi_{N}(x)=1$ for $\nu$-a.e. $x$ and $\lim _{N} \int\left|\nabla \xi_{N}\right|^{2} / \xi_{N} d \nu=0$. Then applying the Fatou lemma we justify (54).

It is helpful to have in mind that $\Phi_{i}^{j k} \Phi_{i j k}=\operatorname{Tr}\left[\left(D^{2} \Phi\right)^{-\frac{1}{2}} D^{2} \Phi_{i}\left(D^{2} \Phi\right)^{-\frac{1}{2}}\right]^{2} \geq 0$. From (54),

$$
\int V_{i i} \Phi_{i i}^{p} d \mu \geq p \int \Phi_{i i}^{p-1}\left\langle\left(D^{2} \Phi\right)^{-1} \nabla \Phi_{i i}, \nabla \Phi_{i i}\right\rangle d \mu
$$

Let us integrate by parts the left-hand side $\int V_{i i} \Phi_{i i}^{p} d \mu=\int V_{i}^{2} \Phi_{i i}^{p} d \mu-p \int V_{i} \Phi_{i i}^{p-1} \Phi_{i i i} d \mu$. The justification of this integration by parts is much easier, since $D^{2} \Phi$ and $D^{2} V$ are bounded. Applying

$$
2\left|\Phi_{i i i} V_{i}\right| \leq 2\left|V_{i}\right| \sqrt{\Phi_{i i} \cdot\left\langle\left(D^{2} \Phi\right)^{-1} \nabla \Phi_{i i}, \nabla \Phi_{i i}\right\rangle} \leq V_{i}^{2} \Phi_{i i}+\left\langle\left(D^{2} \Phi\right)^{-1} \nabla \Phi_{i i}, \nabla \Phi_{i i}\right\rangle
$$

one obtains

$$
\begin{equation*}
\int V_{i}^{2} \Phi_{i i}^{p} d \mu \geq \int \Phi_{i i}^{p-1}\left\langle\left(D^{2} \Phi\right)^{-1} \nabla \Phi_{i i}, \nabla \Phi_{i i}\right\rangle d \mu \tag{55}
\end{equation*}
$$

Let us show that the right-hand side controls powers of the second derivative $\Phi_{i i}$. Indeed, for every $q \geq 2$ and $\varepsilon>0,0 \leq \tau \leq 1$ the following estimate holds

$$
\begin{aligned}
\int \Phi_{i i}^{q} d \mu & =-(q-1) \int \Phi_{i} \Phi_{i i i} \Phi_{i i}^{q-2} d \mu+\int \Phi_{i} V_{i} \Phi_{i i}^{q-1} d \mu \\
& \leq \varepsilon \int \Phi_{i}^{2} \Phi_{i i}^{q-\tau} d \mu+\frac{(q-1)^{2}}{4 \varepsilon} \int \Phi_{i i}^{q-3+\tau}\left\langle\left(D^{2} \Phi\right)^{-1} \nabla \Phi_{i i}, \nabla \Phi_{i i}\right\rangle d \mu \\
& +\frac{q-1}{q} \int \Phi_{i i}^{q} d \mu+\frac{1}{q} \int\left|\Phi_{i} V_{i}\right|^{q} d \mu
\end{aligned}
$$

Finally,

$$
\begin{aligned}
\int \Phi_{i i}^{q} d \mu & \leq \int\left|\Phi_{i} V_{i}\right|^{q} d \mu+q \varepsilon \int \Phi_{i}^{2} \Phi_{i i}^{q-\tau} d \mu+\frac{q(q-1)^{2}}{4 \varepsilon} \int \Phi_{i i}^{q-3+\tau}\left\langle\left(D^{2} \Phi\right)^{-1} \nabla \Phi_{i i}, \nabla \Phi_{i i}\right\rangle d \mu \\
& \leq \int\left|\Phi_{i} V_{i}\right|^{q} d \mu+q \varepsilon \int \Phi_{i}^{2} \Phi_{i i}^{q-\tau} d \mu+\frac{q(q-1)^{2}}{4 \varepsilon} \int \Phi_{i i}^{q-2+\tau} V_{i}^{2} d \mu
\end{aligned}
$$

Applying Hölder inequalities

$$
\begin{gathered}
\Phi_{i}^{2} \Phi_{i i}^{q-\tau} \leq \frac{q-\tau}{q} \Phi_{i i}^{q}+\frac{\tau}{q}\left|\Phi_{i}\right|^{\frac{2 q}{\tau}} \\
\Phi_{i i}^{q-2+\tau} V_{i}^{2} \leq \varepsilon \Phi_{i i}^{q}+C(\varepsilon, q, \tau)\left|V_{i}\right|^{\frac{2 q}{2-\tau}} \\
\left|\Phi_{i} V_{i}\right|^{q} \leq \frac{2-\tau}{2}\left|V_{i}\right|^{\frac{2 q}{2-\tau}}+\frac{\tau}{2}\left|\Phi_{i}\right|^{\frac{2 q}{\tau}}
\end{gathered}
$$

choosing sufficiently small $\varepsilon$, and applying the change of variables formula $\int\left|\Phi_{i}\right|^{q} d \mu=$ $\int\left|x_{i}\right|^{q} d \nu$ we easily get the claim.

Finally, let us get rid of the assumption $D^{2} W \geq \frac{1}{C} \cdot \operatorname{Id}, D^{2} V \leq C \cdot I d$. To this end we approximate $\mu$ and $\nu$ by measures with smooth potentials satisfying $D^{2} W_{N} \geq \frac{1}{C_{N}} \cdot \mathrm{Id}$, $D^{2} V_{N} \leq C_{N} \cdot \operatorname{Id}$ satisfying $\lim _{N} \int\left|\left(V_{N}\right)_{i}\right|^{2 q} d \mu_{N}=\int\left|V_{i}\right|^{2 q} d \mu$ and $\lim _{N} \int\left|x_{i}\right|^{2 q} d \nu_{N}=$
$\int\left|x_{i}\right|^{2 q} d \nu$. It remains to show that the weak $L^{q}(\mu)$-limit of $\left(\Phi_{N}\right)_{i i}$ coincides with $\Phi_{i i}$. The latter can be easily shown with the help of integration-by-parts and identifications of the poinwise limit $\lim _{N} \nabla \Phi_{N}$ with $\nabla \Phi$ (see the proof of Lemma 5.2).

Remark 5.6. The conclusion of Proposition 5.5 holds without any additional smoothness assumptions. This can be verified by smooth approximations (see again [16] for details). Finally we get that (53) holds for every log-concave measure $\nu$ and measure $\mu$ satisfying $\int\left|V_{i}\right|^{\frac{2 q}{2-\tau}} d \mu<\infty$, where $V_{i}$ is the logarithmic derivative of $\mu$ along $x_{i}$.

Second proof of Lemma 5.2: Let us show how Proposition 5.5 implies (48) above, without appealing to the works by Forzani and Maldonado [9, 10] and Gutierrez and Huang [12] related to Caffarelli's $C^{1, \alpha}$-regularity theory. We use that $\sup _{N} \int\left|\nabla V_{N}\right|^{p} e^{-V_{N}} d x<\infty$, $p>n$. Since $\nu$ is log-concave, all the moments of $\nu$ are finite. Thus Proposition 5.5 implies

$$
\sup _{N} \int\left\|D^{2} \Phi_{N}\right\|_{H S}^{p^{\prime}} e^{-V_{N}} d x<\infty
$$

for any $n<p^{\prime}<p$. Applying that $V_{N}$ are uniformly locally bounded from below, we get that $\sup _{N} \int_{B_{R}}\left\|D^{2} \Phi_{N}\right\|_{H S}^{p^{\prime}} d x<\infty$ for every $R$. Then the result follows from the Morrey embedding theorem.

## 6 Corollaries to Theorem 5.1

Proof of Theorem 1.2. For $A \in M_{n}^{+}(\mathbb{R})$ define

$$
F(A)=f\left(\log \lambda_{1}(A), \ldots, \log \lambda_{n}(A)\right)
$$

where $0<\lambda_{1}(A) \leq \ldots \leq \lambda_{n}(A)$ are the eigenvalues of $A$. According to Lemma 2.4, for any $A \in M_{n}^{+}(\mathbb{R})$,

$$
\begin{equation*}
|\nabla F|(A) \leq|\nabla f|\left(\log \lambda_{1}(A), \ldots, \log \lambda_{n}(A)\right) \tag{56}
\end{equation*}
$$

Since $f$ is locally-Lipschitz and the eigenvalues vary continuously with the matrix $A$, then (56) implies that also $F$ is locally-Lipschitz. Denote by $\theta$ the push-forward of the probability measure $\mu$ under the map $x \mapsto D^{2} \Phi(x)$. Since $\mathbb{E}|f(\Lambda(X))|<\infty$ then $F \in L^{1}(\theta)$. Since $\mathbb{E}|\nabla f|^{2}(\Lambda(X))<\infty$, then $\int|\nabla F|^{2} d \theta<\infty$. We may apply Theorem 5.1 and conclude that

$$
\int_{M_{n}^{+}(\mathbb{R})} F^{2} d \theta-\left(\int_{M_{n}^{+}(\mathbb{R})} F d \theta\right)^{2} \leq 4 \int_{M_{n}^{+}(\mathbb{R})}|\nabla F|^{2} d \theta
$$

The left-hand side equals $\operatorname{Var}[f(\Lambda(X))]$. Glancing at (56), we thus obtain

$$
\operatorname{Var}[f(\Lambda(X))] \leq 4 \mathbb{E}|\nabla f|^{2}(\Lambda(X))
$$

and the proof is complete.

Proof of Theorem 1.1. Plug in $f(x)=x_{i}$ in Theorem 1.2. Then $f$ is a 1-Lipschitz function, by Remark 1.4 we have $\mathbb{E}|f(\Lambda(X))|<\infty$. Thus the application of Theorem 1.2 is legitimate, and Theorem 1.1 follows.

Proof of Theorem 1.5. The argument is almost identical to the proof of Theorem 1.1, with Lemma 2.3 replacing the role of Lemma 2.4.

Let us end this paper with a few remarks concerning future research. If we make further assumptions regarding the log-concave measures in question, it is possible to prove concentration inequalities for the eigenvalues of $D^{2} \Phi$ themselves, and not only for their logarithms. The analysis of the weighted Riemannian manifold $M_{\mu, \nu}$ leads to such concentration inequalities. Additionally, there is a soft argument which shows that when $\nabla \Phi$ is the Brenier map between the uniform measure on $K$ and the uniform measure on $T$,

$$
\int_{K} \Delta \Phi \leq n V(K, \ldots, K, T)
$$

where $V$ stands for mixed volume. The details will be discussed elsewhere. Another possible research direction is to investigate whether phenomena similar to Theorem 1.1 occur also in a non-linear setting, when transporting measures with convexity properties supported on Riemannian manifolds.

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