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# OPTIMALITY CONDITIONS FOR MAXIMIZERS OF THE INFORMATION DIVERGENCE FROM AN EXPONENTIAL FAMILY 

František Matúš

The information divergence of a probability measure $P$ from an exponential family $\mathcal{E}$ over a finite set is defined as infimum of the divergences of $P$ from $Q$ subject to $Q \in \mathcal{E}$. All directional derivatives of the divergence from $\mathcal{E}$ are explicitly found. To this end, behaviour of the conjugate of a log-Laplace transform on the boundary of its domain is analysed. The first order conditions for $P$ to be a maximizer of the divergence from $\mathcal{E}$ are presented, including new ones when $P$ is not projectable to $\mathcal{E}$.
Keywords: Kullback-Leibler divergence, relative entropy, exponential family, information projection, log-Laplace transform, cumulant generating function, directional derivatives, first order optimality conditions, convex functions, polytopes
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## 1. INTRODUCTION

Let $\nu$ be a nonzero measure on a finite set $Z$ and $f$ a mapping from $Z$ into the $d$ dimensional Euclidean space $\mathbb{R}^{d}$. The (full) exponential family $\mathcal{E}=\mathcal{E}_{\nu, f}$ determined by $\nu$ and the directional statistic $f$, see $[5,6,7,12]$, consists of the probability measures (pm's) $Q_{\vartheta}=Q_{\nu, f, \vartheta}, \vartheta \in \mathbb{R}^{d}$, given by

$$
Q_{\vartheta}(z)=e^{\langle\vartheta, f(z)\rangle-\Lambda(\vartheta)} \nu(z), \quad z \in Z
$$

where $\langle\cdot, \cdot\rangle$ is the scalar product on $\mathbb{R}^{d}$ and

$$
\Lambda(\vartheta)=\Lambda_{\nu, f}(\vartheta)=\ln \sum_{z \in Z} e^{\langle\vartheta, f(z)\rangle} \nu(z) .
$$

The information divergence (relative entropy) of a $\mathrm{pm} P$ on $Z$ from $\nu$ is

$$
D(P \| \nu)= \begin{cases}\sum_{z \in s(P)} P(z) \ln \frac{P(z)}{\nu(z)}, & s(P) \subseteq s(\nu) \\ +\infty, & \text { otherwise }\end{cases}
$$

where $s(\nu)=\{z \in Z: \nu(z)>0\}$ is the support of $\nu$. The information divergence of $P$ from the exponential family $\mathcal{E}$ is defined by

$$
\begin{equation*}
D(P \| \mathcal{E})=\inf _{\vartheta \in \mathbb{R}^{d}} D\left(P \| Q_{\vartheta}\right) . \tag{1}
\end{equation*}
$$

This work studies the maximizers of the function $P \mapsto D(P \| \mathcal{E})$, denoted also by $D(\cdot \| \mathcal{E})$, over the pm's $P$ dominated by $\nu$, thus satisfying $s(P) \subseteq s(\nu)$.

The maximization of $D(\cdot \| \mathcal{E})$ has emerged in probabilistic models for evolution and learning in neural networks that are based on infomax principles [1, 2]. The divergence from an exponential family can be related to information theoretic measures for interdependence of stochastic units and its maximization reveals stochastic systems with high complexity w.r.t. an exponential family [3]. Dynamical versions of the problem of interactions in recurrent networks appeared in $[1,4,16]$. Two special instances of the maximization of the quantity (1) are described in [13]. Further relations to previous works on this problem are discussed in remarks of Section 5.

Let $\mu$ be the $f$-image $\nu_{f}$ of $\nu$, considered for a Borel measure on $\mathbb{R}^{d}$. Denoting by $s(\mu)$ the support $f(s(\nu))$ of $\mu$, which is the inclusion-minimal closed subset of $\mathbb{R}^{d}$ of the $\mu$-measure $\mu\left(\mathbb{R}^{d}\right)$,

$$
\Lambda(\vartheta)=\Lambda_{\mu}(\vartheta)=\ln \sum_{x \in s(\mu)} e^{\langle\vartheta, x\rangle} \mu(x)
$$

whence $\Lambda$ equals the log-Laplace transform (cumulant generating function) of $\mu$. In terms of the conjugate $\Lambda^{*}$ of $\Lambda$ [15, Section 12],

$$
\Lambda^{*}(a)=\sup _{\vartheta \in \mathbb{R}^{d}}[\langle\vartheta, a\rangle-\Lambda(\vartheta)], \quad a \in \mathbb{R}^{d}
$$

the information divergence of a pm $P$ from the exponential family $\mathcal{E}$ rewrites to
where

$$
\begin{equation*}
D(P \| \mathcal{E})=D(P \| \nu)-\Lambda^{*}\left(m\left(P_{f}\right)\right) \tag{2}
\end{equation*}
$$

is the mean of the $f$-image $P_{f}$ of $P$, or equivalently, the $P$-mean of $f$. Hence, $D(\cdot \| \mathcal{E})$ is difference of the strictly convex function $P \mapsto D(P \| \nu)$, denoted by $D(\cdot \| \nu)$, and the function $P \mapsto \Lambda^{*}\left(m\left(P_{f}\right)\right)$, which is convex because $\Lambda^{*}$ is convex and $P \mapsto m\left(P_{f}\right)$ is linear.

From now on assume that $\nu(z)$ is positive for each $z \in Z$.
This work is organized as follows. After introduction of notations and review of known facts in Section 2, directional behavior of the conjugate $\Lambda^{*}$ on a boundary of its domain is described in Theorem 3.1 of Section 3. Consequently in Section 4 it is shown, relying on (2), that the one-sided directional derivatives of the function $D(\cdot \| \mathcal{E})$ at any pm $P$ exist, possibly taking the values $\pm \infty$. Explicit formulas for the derivatives are presented in Theorems 4.1 and 4.3. The first order optimality conditions for a pm $P$ to be a maximizer of $D(\cdot \| \mathcal{E})$ emerge by requiring the derivatives not to be positive, see Theorem 5.1 in Section 5. Finally, Section 6 is devoted to a proof of Theorem 3.1.

## 2. PRELIMINARIES

This section reviews well-known facts about the log-Laplace transforms, their conjugates and exponential families, and introduces necessary notations. It is assumed throughout that $\mu$ is a nonzero positive finite measure on $\mathbb{R}^{d}$ that is concentrated on a finite set, though many of the assertions below are valid under more general assumptions $[5,6,12]$.

In accordance with [15], the affine hull of a set $B \subseteq \mathbb{R}^{d}$ is denoted by aff $(B)$, the shift of $\operatorname{aff}(B)$ containing the origin by $\operatorname{lin}(B)$ and the relative interior of $B$ by $r i(B)$, which is the interior of $B$ in the topology of aff $(B)$. The orthogonal complement of a linear subspace $E$ of $\mathbb{R}^{d}$ is denoted by $E^{\perp}$.

Since $\mu$ is concentrated on a finite set the convex support cs $(\mu)$ of $\mu$, which is the inclusion-minimal closed convex subset of $\mathbb{R}^{d}$ of the $\mu$-measure $\mu\left(\mathbb{R}^{d}\right)$, is the polytope spanned by $s(\mu)$. For $B=c s(\mu)$ above notations are abbreviated to aff $(\mu)$, lin $(\mu)$ and $r i(\mu)$.

Fact 2.1. For $a \in \operatorname{aff}(\mu)$ and $c \in \mathbb{R}^{d}$ the function $\vartheta \mapsto\langle\vartheta, a\rangle-\Lambda(\vartheta)$ is constant on $c+\operatorname{lin}(\mu)^{\perp}$.

Let $Q_{\mu, \vartheta}$ denote the pm with $\mu$-density $x \mapsto e^{\langle\vartheta, x\rangle-\Lambda(\vartheta\rangle}, \vartheta \in \mathbb{R}^{d}$, and $\mathcal{E}_{\mu}$ the family of all such pm's, thus the standard exponential family determined by $\mu$ and the identity mapping on $\mathbb{R}^{d}$, in the role of a directional statistic.

Fact 2.2. The equality $Q_{\mu, \vartheta}=Q_{\mu, \theta}$ holds if and only if $\vartheta-\theta \in \operatorname{lin}(\mu)^{\perp}$.
The mean $\sum_{x \in s(\mu)} x e^{\langle\vartheta, x\rangle-\Lambda(\vartheta)}$ of $Q_{\mu, \vartheta}$ is denoted by $m\left(Q_{\mu, \vartheta}\right)$.
Fact 2.3. $m\left(Q_{\mu, \vartheta}\right)=\nabla \Lambda(\vartheta), \vartheta \in \mathbb{R}^{d}$.
Fact 2.4. The restriction of $\nabla \Lambda$ to $\operatorname{lin}(\mu)$ is injective and onto ri $(\mu)$.
Since $\Lambda$ is smooth this restriction is a diffeomorphism. Its inverse is denoted in the sequel by $\psi=\psi_{\mu}$. With this notation the mean parametrization $a \mapsto Q_{\mu, \psi(a)}$ of $\mathcal{E}_{\mu}$ by the elements $a$ of $r i(\mu)$ is bijective.

Fact 2.5. If $m\left(Q_{\mu, \vartheta}\right)=a$ then $\Lambda^{*}(a)=\langle\vartheta, a\rangle-\Lambda(\vartheta)$.
In particular, $\Lambda^{*}(a)=\langle\psi(a), a\rangle-\Lambda(\psi(a))$ for $a \in \operatorname{ri}(\mu)$ because $m\left(Q_{\mu, \psi(a)}\right)=a$.
Fact 2.6. If $a \in r i(\mu)$ then for $b \in \operatorname{cs}(\mu)$

$$
\Lambda^{*}(b)=\Lambda^{*}(a)+\langle\psi(a), b-a\rangle+o(\|b-a\|) .
$$

The following assertion is a consequence of [8, Lemma 6].
Fact 2.7. If $a \in \operatorname{cs}(\mu) \backslash r i(\mu)$ then $+\infty>\Lambda^{*}(a)>\langle\vartheta, a\rangle-\Lambda(\vartheta)$ for all $\vartheta \in \mathbb{R}^{d}$.
Hence, the convex conjugate $\Lambda^{*}$ is finite on cs $(\mu)$, thus continuous on this polytope due to convexity. Beyond the polytope it takes the value $+\infty$, e.g. by $[8$, Proposition 1 (ii)].

If $A$ is a Borel subset of $\mathbb{R}^{d}$ then $B \mapsto \mu(B \cap A)$ is the restriction of $\mu$ by $A$. Let $\Lambda_{\mu}, \mathcal{E}_{\mu}, Q_{\mu, \vartheta}, \psi_{\mu}$, etc., with $\mu$ replaced by its nonzero restriction be denoted by $\Lambda_{A}$, $\mathcal{E}_{A}, Q_{A, \vartheta}, \psi_{A}$, etc. The following assertion follows from [8, Lemma 6].

Fact 2.8. If $a \in F$ for a face $F$ of $c s(\mu)$ then $\Lambda_{\mu}^{*}(a)=\Lambda_{F}^{*}(a)$.

Fact 2.9. If $a \in \operatorname{ri}(F)$ for a face $F$ of $c s(\mu)$ then

$$
\Lambda^{*}(a)-[\langle\vartheta, a\rangle-\Lambda(\vartheta)]=D\left(Q_{F, \psi_{F}(a)} \| Q_{\mu, \vartheta}\right), \quad \vartheta \in \mathbb{R}^{d}
$$

This identity admits a substantial generalization, see [10, Theorem 4.1] where $Q_{F, \psi_{F}(a)}$ is called a generalized maximum likelihood estimator for $\mathcal{E}_{\mu}$.

Suppose in the remaining part of this section that $\mu$ is the $f$-image $\nu_{f}$ of $\nu$ as in Introduction and recall that $s(\nu)=Z$. By (2) and the continuity of $\Lambda^{*}$ on $\operatorname{cs}(\mu)$, the function $D\left(\cdot \| \mathcal{E}_{\nu, f}\right)$ is continuous and therefore has a maximizer.

In this situation, $Q_{\mu, \vartheta}$ is the $f$-image of $Q_{\nu, f, \vartheta}, \vartheta \in \mathbb{R}^{d}$. Taking the $f$-images of pm's from $\mathcal{E}_{\nu, f}$ is a bijection onto $\mathcal{E}_{\mu}$ and $a \mapsto Q_{\nu, f, \psi(a)}$ is the mean parametrization of $\mathcal{E}_{\nu, f}$ by the elements $a$ of $r i(\mu)$.

It follows from (2) that the infimum in (1) is attained if and only if the mean $a=m\left(P_{f}\right)$ belongs to $r i(\mu)$ in which case $Q_{\nu, f, \psi(a)}$ is the unique pm $Q$ of $\mathcal{E}=$ $\mathcal{E}_{\nu, f}$ satisfying $D(P \| Q)=D(P \| \mathcal{E})$. This pm is called the reverse information (rI-) projection of $P$ on $\mathcal{E}$ and is denoted by $\Pi_{P \rightarrow \mathcal{E}}$, as in [8].

For a face $F$ of $c s(\mu)$ the $\mathrm{pm} Q_{F, \vartheta}$ is the $f$-image of the $\mathrm{pm} Q_{Y, f, \vartheta}, \vartheta \in \mathbb{R}^{d}$, where the latter denotes the pm obtained from $Q_{\nu, f, \vartheta}$ when $\nu$ is replaced by its restriction to $Y=f^{-1}(F)$. Taking the $f$-images of pm's from the exponential family $\mathcal{E}_{Y, f}=\left\{Q_{Y, f, \vartheta}: \vartheta \in \mathbb{R}^{d}\right\}$ is a bijection onto $\mathcal{E}_{F}$.

The (variation) closure $c l$ of $\mathcal{E}_{\nu, f}$, respectively $\mathcal{E}_{\mu}$, is equal to union of the families $\mathcal{E}_{f^{-1}(F), f}$, respectively $\mathcal{E}_{F}$, over the nonempty faces $F$ of $\boldsymbol{c s}(\mu)$, for a general result see [ 9 , Theorem 2]. The closures are bijectively parameterized by means of pm's, exhausting cs $(\mu)$.

Fact 2.10. For $\mathcal{E}=\mathcal{E}_{\nu, f}$ and a pm $P$ on $Z$,

$$
D(P \| \mathcal{E})=D(P \| c l(\mathcal{E}))=\min _{Q \in c l(\mathcal{E})} D(P \| Q)
$$

The minimum is attained uniquely by $Q_{f^{-1}(F), f, \psi_{F}(a)}$ where $a=m\left(P_{f}\right)$ and $F$ is that face of $c s(\mu)$ which contains $a$ in its relative interior.

This unique minimizer is denoted by $\Pi_{P \rightarrow \mathcal{E}}$, extending the above notation to the cases when the infimum in (1) is not attained, and called the generalized rI-projection of $P$ to $\mathcal{E}$.

For $a \in \operatorname{cs}(\mu)$ write $\Pi_{a \leftrightarrow \mathcal{E}}=Q_{f^{-1}(F), f, \psi_{F}(a)}$ where $F$ is the face of $c s(\mu)$ with $a \in r i(F)$. The following Pythagorean identity provides more insight into the minimization over $c l(\mathcal{E})$ in Fact 2.10, compare with [11, Proposition 4]. For a general version see [ 8 , Theorem 1].

Fact 2.11. If $\mathcal{E}=\mathcal{E}_{\nu, f}$ and $a \in \operatorname{cs}(\mu)$ then

$$
\begin{equation*}
D(P \| Q)=D\left(P \| \Pi_{a \mapsto \mathcal{E}}\right)+D\left(\Pi_{a \leftrightarrow \mathcal{E}} \| Q\right) \tag{3}
\end{equation*}
$$

holds for any pm $P$ with $m\left(P_{f}\right)=a$ and $Q \in c l(\mathcal{E})$.
Given a pm $P$ on $Z$ and a set $Y \subseteq Z$ with $P(Y)>0$ let $P^{Y}$ denote the pm, called truncation in [7], given by $P^{Y}(z)=P(z) / P(Y)$ for $z \in Y$ and $P^{Y}(z)=0$ otherwise. Note that the set $\left\{Q \in \mathcal{E}_{\nu, f}: Q^{Y}=P^{Y}\right\}$, though not given via a directional statistic, is a full exponential family provided it is nonempty. Also $\left\{Q_{\nu, f, \vartheta}: \vartheta \in \Xi\right\}$ is such a family when $\Xi$ is an affine subset of $\mathbb{R}^{d}$.

## 3. CONJUGATE OF LOG-LAPLACE TRANSFORM

In this section, $\mu$ is a positive measure on $\mathbb{R}^{d}$ concentrated on a finite set.
Each point $a$ of the polytope $\operatorname{cs}(\mu)$ belongs to the relative interior $r i(F)$ of a unique face $F$ of $c s(\mu)$. If $b \in F$ then Facts 2.6 and 2.8 combine to

$$
\begin{equation*}
\Lambda^{*}(a+\varepsilon(b-a))=\Lambda^{*}(a)+\varepsilon\left\langle\psi_{F}(a), b-a\right\rangle+o(\varepsilon), \tag{4}
\end{equation*}
$$

describing the directional behavior of the function $\varepsilon \mapsto \Lambda^{*}(a+\varepsilon(b-a))$ in a neighborhood of 0 .

Let $C$ denote the convex hull of $s(\mu) \backslash F$ and $C_{+}=C+\operatorname{lin}(F)$.
If $b \in \operatorname{cs}(\mu) \backslash F$ then it is not difficult to see that there exists a positive $t$ such that $a+t(b-a)$ belongs to $C_{+}$and $a \notin C_{+}$, see Lemmas 6.1 and 6.2. Then such a minimal $t>0$ exists. Denote this number by $t_{a b}$ and the nearest point $a+t_{a b}(b-a)$ of $C_{+}$from $a$ in the direction $b-a$ by $x_{a b}$.

Let $\Xi=\psi_{F}(a)+\operatorname{lin}(F)^{\perp}$ and

$$
\Psi_{C, \Xi}^{*}(x)=\sup _{\theta \in \Xi}\left[\langle\theta, x\rangle-\Lambda_{C}(\theta)\right], \quad x \in \mathbb{R}^{d}
$$

By Lemma 6.8 and Fact $2.5, \Psi_{C, \Xi}^{*}\left(x_{a b}\right)$ is finite.
Theorem 3.1. If $a \in \operatorname{ri}(F)$ for a face $F$ of $\operatorname{cs}(\mu), b \in \operatorname{cs}(\mu) \backslash F$ and $\varepsilon>0$ then

$$
\Lambda^{*}\left(a+\varepsilon t_{a b}(b-a)\right)=\Lambda^{*}(a)+h(\varepsilon)+\varepsilon\left[\Psi_{C, \Xi}^{*}\left(x_{a b}\right)-\Lambda^{*}(a)\right]+o(\varepsilon)
$$

where $h(\varepsilon)=\varepsilon \ln \varepsilon+(1-\varepsilon) \ln (1-\varepsilon)$.
The proof of Theorem 3.1, preceded by several lemmas, is presented in Section 6. It is independent of the material of Sections 4 and 5.

The following figure illustrates the notations presented above or used later in proofs: the support of $\mu$ consists of five points depicted as black squares, $F$ is the vertical edge of the pentagon $\operatorname{cs}(\mu), C$ is the shaded triangle and $C_{+}$a vertical infinite strip.


## 4. THE DIRECTIONAL DERIVATIVES OF $D(\cdot \| \mathcal{E})$

In this section, $\mu=\nu_{f}, \mathcal{E}=\mathcal{E}_{\nu, f}, P$ and $R$ are pm's on $Z, a=m\left(P_{f}\right)$ belongs to $r i(F)$ for a face $F$ of $c s(\mu), \vartheta=\psi_{F}(a), b=m\left(R_{f}\right)$, and $r=R(Z \backslash s(P))$.

As well-known, the one-sided directional derivative of $D(\cdot \| \mathcal{E})$ at $P$ in the direction $R-P$ is given by

$$
\lim _{\varepsilon \downarrow 0} \frac{1}{\varepsilon}[D(P+\varepsilon(R-P) \| \mathcal{E})-D(P \| \mathcal{E})]
$$

provided the limit, finite or infinite, exists. If $P$ dominates $R$ then even the two-sided limiting $\varepsilon \rightarrow 0$ makes sense.

Theorem 4.1. If $b \in F$ and $r=0$ then the two-sided derivative of $D(\cdot \| \mathcal{E})$ at $P$ in the direction $R-P$ equals

$$
\begin{equation*}
\sum_{z \in s(P)}[R(z)-P(z)] \ln \frac{P(z)}{e^{\langle\vartheta, f(z)\rangle} \nu(z)} . \tag{5}
\end{equation*}
$$

If $b \in F$ and $r>0$ then the one-sided directional derivative of $D(\cdot \| \mathcal{E})$ at $P$ in the direction $R-P$ equals $-\infty$.

If $b \notin F$ then this derivative is equal to

$$
\begin{cases}+\infty, & r t_{a b}<1 \\ -\infty, & r t_{a b}>1 \\ T-r\left[\Psi_{C, \Xi}^{*}\left(x_{a b}\right)-\Lambda^{*}(a)+\ln r\right], & r t_{a b}=1\end{cases}
$$

where

$$
T=\sum_{z \in s(R) \backslash s(P)} R(z) \ln \frac{R(z)}{\nu(z)}+\sum_{z \in s(P)}[R(z)-P(z)] \ln \frac{P(z)}{\nu(z)}
$$

A proof invokes the following simple assertion, demonstrated for convenience at the end of Section 6.

Lemma 4.2. If $\varepsilon>0$ then

$$
D(P+\varepsilon(R-P) \| \nu)=D(P \| \nu)+h(\varepsilon) r+\varepsilon T+o(\varepsilon) .
$$

If additionally $r=0$ then this holds with the $h(\varepsilon)$-term omitted also for $\varepsilon \leqslant 0$.
Proof of Theorem 4.1. If $b \in F$ and $r=0$ then $s(R) \subseteq s(P)$. On account of (2) the derivative equals the difference of coefficients at the $\varepsilon$-terms in Lemma 4.2 and (4)

$$
\sum_{z \in s(P)}[R(z)-P(z)] \ln \frac{P(z)}{\nu(z)}-\langle\vartheta, b-a\rangle
$$

which rewrites to (5).
By the same argument, if $b \in F$ and $r>0$ then the one-sided derivative is equal to $-\infty$, due to the nonzero term $h(\varepsilon) r$ in Lemma 4.2.

If $b \notin F$ then the formula of Theorem 3.1 is equivalent to

$$
\Lambda^{*}(a+\varepsilon(b-a))=\Lambda^{*}(a)+h(\varepsilon) \frac{1}{t_{a b}}+\frac{\varepsilon}{t_{a b}}\left[\Psi_{C, \Xi}^{*}\left(x_{a b}\right)-\Lambda^{*}(a)+\ln \frac{1}{t_{a b}}\right]+o(\varepsilon) .
$$

This, (2) and Lemma 4.2 imply the last assertion.
The case $b \notin F$ in Theorem 4.1 can be further simplified when assuming that there exist two different parallel hyperplanes $H_{F}$ and $H_{C}$ such that $H_{F} \supseteq F$ and $H_{C} \supseteq C$, where $C$ is the convex hull of $s(\mu) \backslash F$. This implies obviously that $F_{+}=F+\operatorname{lin}(C)$ and $C_{+}=C+\operatorname{lin}(F)$ are disjoint. The implication can be reversed. In fact, by [15, Corollary 19.3.3] if the polyhedral sets $F_{+}$and $C_{+}$are disjoint then it is possible to separate them strongly by a hyperplane $H$, and then $\operatorname{lin}(H)$ contains $\operatorname{lin}(F)$ and $\operatorname{lin}(C)$, thus different shifts of $H$ contain $F$ and $C$.

Theorem 4.3. If $b \notin F$ and $F_{+} \cap C_{+}=\emptyset$ then $r t_{a b} \geqslant 1$. The equality holds here if and only if $R\left(f^{-1}(F) \backslash s(P)\right)=0$ in which case the one-sided directional derivative of $D(\cdot \| \mathcal{E})$ at $P$ in the direction $R-P$ is equal to

$$
\begin{equation*}
r\left[D\left(R^{Y} \| \mathcal{F}\right)-D(P \| \mathcal{E})\right]+(1-r) \sum_{z \in s(P)}\left[R^{s(P)}(z)-P(z)\right] \ln \frac{P(z)}{e^{\langle\theta, f(z)\rangle} \nu(z)} \tag{6}
\end{equation*}
$$

where $Y=f^{-1}(C), \mathcal{F}$ is the exponential family consisting of $Q_{Y, f, \theta}, \theta \in \Xi$, and the truncation $R^{s(P)}$ is well-defined if $r<1$ or does not enter otherwise.

Proof. The first assumption implies $R_{f}(F)<1$ whence $s=R(Y)$ is positive. Then, $R=s R^{Y}+(1-s) Q$ for the truncation $R^{Y}$ and a pm $Q$ concentrated on $f^{-1}(F)=Z \backslash Y$. Thus, $b=m\left(R_{f}\right)$ equals $s c+(1-s) a^{\prime}$ where $c=m\left(R_{f}^{Y}\right) \in C$ and $a^{\prime}=m\left(Q_{f}\right) \in F$. Rewrite $a+\frac{1}{s}(b-a)$ to $c+\frac{1-s}{s}\left(a^{\prime}-a\right)$ to conclude that it belongs to $C_{+}$. The second assumption implies that $a \in F$ and $C_{+}$are contained in two parallel hyperplanes whence $a+t(b-a) \in C_{+}$for a unique $t$. Then, $t_{a b}=\frac{1}{s}$ and $r t_{a b} \geqslant 1$ follows from the obvious inequality $r \geqslant s$. The second assertion obtains from the equivalence of $r=s$ to $R\left(f^{-1}(F) \backslash s(P)\right)=0$.

If this holds then

$$
T=r D\left(R^{Y} \| \nu\right)+r \ln r-r D(P \| \nu)+\sum_{z \in s(P)}[R(z)-(1-r) P(z)] \ln \frac{P(z)}{\nu(z)}
$$

and the derivative equals

$$
r\left[D\left(R^{Y} \| \nu\right)-\Psi_{C, \Xi}^{*}\left(x_{a b}\right)-D(P \| \mathcal{E})\right]+(1-r) \sum_{z \in s(P)}\left[R^{s(P)}(z)-P(z)\right] \ln \frac{P(z)}{\nu(z)}
$$

where the truncation $R^{s(P)}$ is well-defined if $r<1$. Since $x_{a b}=c+\frac{1-r}{r}\left(a^{\prime}-a\right)$, $a^{\prime}-a \in \operatorname{lin}(F)$, and $a^{\prime}$ is the mean of the $f$-image of $R^{s(P)}=Q$ provided $r<1$,

$$
\begin{aligned}
r \Psi_{C, \Xi}^{*}\left(x_{a b}\right) & =r \Psi_{C, \Xi}^{*}(c)+(1-r)\left\langle\vartheta, a^{\prime}-a\right\rangle \\
& =r \Psi_{C, \Xi}^{*}(c)+(1-r) \sum_{z \in s(P)}\left[R^{s(P)}(z)-P(z)\right]\langle\vartheta, f(z)\rangle
\end{aligned}
$$

Using also the analogue of (2)

$$
D\left(R^{Y} \| \mathcal{F}\right)=\inf _{\theta \in \Xi} D\left(R^{Y} \| Q_{Y, f, \theta}\right)=D\left(R^{Y} \| \nu\right)-\Psi_{C, \Xi}^{*}(c)
$$

the above expression for the derivative rewrites to (6).

Sometimes the above simplification of Theorem 4.1 is not available but such situations are not encountered later, due to the following observation proved at the end of Section 6.

Lemma 4.4. If $F_{+} \cap C_{+} \neq \emptyset$ then there exists a pm $Q$ concentrated on $Z \backslash f^{-1}(F)$ such that the derivative of $D(\cdot \| \mathcal{E})$ at $P$ in the direction $Q-P$ is $+\infty$.

## 5. OPTIMALITY CONDITIONS

The results on directional derivatives of $D(\cdot \| \mathcal{E})$ presented in the previous section imply first order necessary conditions for a pm to be a maximizer of this function.

Theorem 5.1. If $\mathcal{E}=\mathcal{E}_{\nu, f}$ and $P$ is a maximizer of $D(\cdot \| \mathcal{E})$ then $P$ equals the truncation $\Pi^{s(P)}$ to $s(P)$ of the generalized $r I$-projection $\Pi=\Pi_{P \rightarrow \mathcal{E}}$ of $P$ to $\mathcal{E}$. If additionally $P$ is not $r I$-projectable to $\mathcal{E}$, thus $s(\Pi) \neq Z$, then $f$ maps $s(\Pi)$ and $Z \backslash s(\Pi)$ into two different parallel hyperplanes, correspondingly, and

$$
\begin{equation*}
D(P \| \mathcal{E}) \geqslant \max \left\{D\left(R \| \mathcal{E}^{\Pi}\right): R \text { is a pm on } Z \text { with } s(R) \subseteq Z \backslash s(\Pi)\right\} \tag{7}
\end{equation*}
$$

where $\mathcal{E}^{\Pi}=\left\{Q^{Z \backslash s(\Pi)}: Q \in \mathcal{E}\right.$ and $\left.Q^{s(\Pi)}=\Pi\right\}$.
Remark 5.2. The condition $P=\Pi^{s(P)}$ goes back to [1, Proposition 3.1] under the assumption that $P$ is $r I$-projectable to $\mathcal{E}$.

Remark 5.3. The maximization of $D\left(\cdot \| \mathcal{E}^{\Pi}\right)$ in (7) is a problem of the same type as that of $D(\cdot \| \mathcal{E})$, however, the family $\mathcal{E}^{\Pi}$ is of a smaller dimension than $\mathcal{E}$.

Proof. Using the notation of Section 4 and Fact 2.10, $s(\Pi)$ is equal to $f^{-1}(F)$ and $\Pi=Q_{s(\Pi), f, \vartheta}$. Since $P$ is a maximizer the two-sided derivatives of $D(\cdot \| \mathcal{E})$ at $P$ vanish, hence by Theorem 4.1

$$
\sum_{z \in s(P)}[R(z)-P(z)] \ln \frac{P(z)}{\Pi(z)}=0
$$

for all $R$ dominated by $P$. This implies $P=\Pi^{s(P)}$.
If $s(\Pi) \neq Z$ then, avoiding directional derivatives to be $+\infty$, Lemma 4.4 implies $F_{+} \cap C_{+}=\emptyset$, thus the containments in hyperplanes. By Theorem 4.3, where $Y=$ $Z \backslash s(\Pi)$, for all pm's $R$ sitting on $Y$ the expression (6) is nonnegative, and thus $D(P \| \mathcal{E}) \geqslant D(R \| \mathcal{F})$, having $r=1$. To prove the inequality (7) it suffices to show that $\mathcal{F}=\mathcal{E}^{\Pi}$. For this, observe that $Q_{Y, f, \theta}$ is the truncation of $Q_{\nu, f, \theta}$ to $Y$ while by Fact 2.2 the truncation $Q_{s(\Pi), f, \theta}$ of $Q_{\nu, f, \theta}$ equals $\Pi$ if and only if $\theta-\vartheta \in \operatorname{lin}(F)^{\perp}$, thus $\theta \in \Xi$.

Remark 5.4. It is not difficult to reverse argumentation in the previous proof and show that the conclusions of Theorem 5.1 hold for a pm $P$ if and only if no directional derivative of $D(\cdot \| \mathcal{E})$ at $P$ is positive.

Other necessary conditions for maximizers can be formulated as follows.

Proposition 5.5. If $P$ is a maximizer of $D(\cdot \| \mathcal{E})$ then $f$ restricted to $s(P)$ is injective and $f(s(P))$ is an affinely independent set.

Proof. On account of (2), the function $R \mapsto D(R \| \mathcal{E})$ is strictly convex on the polytope $\left\{R: m\left(R_{f}\right)=a\right\}$ where $a=m\left(P_{f}\right)$. Since $P$ is a maximizer it must be an extreme point of this polytope by [15, Theorem 32.1]. This implies the assertions. $\square$

Remark 5.7. As a consequence, the cardinality of $s(P)$ is at most one more than the affine dimension of $F$ where $F$ is the face of $\operatorname{cs}\left(\nu_{f}\right)$ with $m\left(P_{f}\right) \in r i(F)$. This implies that this cardinality is bounded above by $1+\operatorname{dim}(\mathcal{E})$, as observed in $[1$, Proposition 3.2] for $r I$-projectable maximizers and in [14, Corollary 2] in general.

Example 5.7. Let $Z$ consist of six points in the plane, depicted as squares in the picture below. Let $a=(0,0), b=(0,1), \nu(z)=1$ for $z \in Z \backslash\{b\}, \nu(b)=w>0$ and $f$ be the identity mapping on $\mathbb{R}^{2}$ restricted to $Z$.


Since

$$
m\left(Q_{(0, u)}\right)=\frac{(w+2) e^{u}}{3+(w+2) e^{u}} b, \quad u \in \mathbb{R},
$$

where the fraction equals a positive $\varepsilon$ if and only if $e^{u}(w+2)=\frac{3 \varepsilon}{1-\varepsilon}$

$$
\psi(\varepsilon b)=\left(0, \ln \frac{\varepsilon}{1-\varepsilon} \frac{3}{w+2}\right), \quad 0<\varepsilon<1 .
$$

By Fact 2.5 and this formula for $\psi(\varepsilon b)$,

$$
\begin{equation*}
\Lambda^{*}(\varepsilon b)=\varepsilon \ln \frac{\varepsilon}{1-\varepsilon} \frac{3}{w+2}-\ln \left[3+\frac{3 \varepsilon}{1-\varepsilon}\right]=-\ln 3+h(\varepsilon)+\varepsilon \ln \frac{3}{w+2} \tag{8}
\end{equation*}
$$

which is in accordance with Theorem 3.1 where $t_{a b}=1, x_{a b}=b, \Lambda^{*}(a)=-\ln 3$ and $\Psi_{\Xi, C}^{*}(b)=-\ln (w+2)$. Note that $\Xi$ is the vertical axis and the expression $\langle\theta, b\rangle-\Lambda_{C}(\theta)$ is constant for $\theta \in \Xi$ by Fact 2.1.

Consider the point masses $P$ and $R$ sitting at $a$ and $b$, respectively. Then, the mean of $\varepsilon P+(1-\varepsilon) R$ is $\varepsilon b$ and, by (2) and (8),

$$
D((1-\varepsilon) P+\varepsilon R \| \mathcal{E})=\ln 3+\varepsilon\left[\ln \frac{w+2}{w}-\ln 3\right], \quad 0 \leqslant \varepsilon \leqslant 1 .
$$

This is in accordance with Theorem 4.3 where $C$ is the upper horizontal edge of the rectangle $\operatorname{cs}\left(\nu_{f}\right), Y=C \cap Z$ has three elements, $r=1, \mathcal{F}$ consists of the single pm $Q_{Y, f,(0,0)}$ and $D(R \| \mathcal{F})=\ln \frac{w+2}{w}$.

Since $m\left(P_{f}\right)$ is not in the interior of $c s\left(\nu_{f}\right)$ the pm $P$ is not $r I$-projectable to $\mathcal{E}$. The generalized projection $\Pi=\Pi_{P \rightarrow \mathcal{E}}=Q_{f^{-1}(F), f,(0,0)}$ is a pm sitting and uniform on $F \cap Z$. Obviously, the first two conclusions of Theorem 5.1 hold for $P$, having the edges $F$ and $C$ contained in two parallel lines. The third conclusion (7) takes the form

$$
\ln 3 \geqslant \max \left\{r_{1} \ln \frac{r_{1}}{1 /(w+2)}+r_{2} \ln \frac{r_{2}}{w /(w+2)}+r_{3} \ln \frac{r_{3}}{1 /(w+2)}:\left(r_{1}, r_{2}, r_{3}\right) \text { is a pm }\right\} .
$$

Note that $\mathcal{E}^{\Pi}$ consists of a single pm as all $Q \in \mathcal{E}$ with $Q^{s(\Pi)}=\Pi$ have the same truncation to $Z \backslash s(\Pi)$. Thus (7) is equivalent to

$$
\ln 3 \geqslant \max \left\{\ln (w+2), \ln \frac{w+2}{w}\right\}
$$

It follows that if $w \neq 1$ then Theorem 5.1 implies that $P$ is not a maximizer while if $w=1$ then all conclusions of Theorem 5.1 take place for $P$, thus the first order necessary conditions do not exclude $P$ to be a maximizer. Actually, in the latter case it is not difficult to prove directly that $D(\cdot \| \mathcal{E}) \leqslant \ln 3$, implying that $(1-\varepsilon) P+\varepsilon R$, $0 \leqslant \varepsilon \leqslant 1$, are global maximizers.

## 6. PROOF OF THEOREM 3.1

Recall the assumptions that $\mu$ is a positive measure concentrated on a finite subset of $\mathbb{R}^{d}, a \in \operatorname{cs}(\mu)$ and $b \in c s(\mu) \backslash F$ where $F$ is the unique face of $c s(\mu)$ such that $a \in r i(F)$.

Lemma 6.1. There exists $t>0$ such that $a+t(b-a) \in C_{+}$.
Proof. Write $b$ as $\varepsilon c+(1-\varepsilon) a^{\prime}$ with $c \in C, a^{\prime} \in F$ and $0<\varepsilon \leqslant 1$, and then for $t=\frac{1}{\varepsilon}$ express $a+t(b-a)$ as $c+t(1-\varepsilon)\left(a^{\prime}-a\right)$ where the second summand belongs to $\operatorname{lin}(F)$.

Lemma 6.2. The face $F$ is contained in a hyperplane disjoint with $C_{+}$.
Proof. Since $F$ is a proper face of $c s(\mu)$ there exists a supporting hyperplane $H$ of $c s(\mu)$ such that $H \cap c s(\mu)=F$. The points of $s(\mu) \backslash F$ belong to one of the open halfspaces associated to $H$. It follows that $C_{+}$is contained in the halfspace, using $\operatorname{lin}(F) \subseteq \operatorname{lin}(H)$.

Lemma 6.3. If $G$ is a face of $C_{+}$then $G, G+\operatorname{lin}(F)$ and $(G \cap C)+\operatorname{lin}(F)$ coincide, $r i(G)=r i(G \cap C)+\operatorname{lin}(F)$ and $G \cap C$ is a face of $C$.

Proof. If $g \in G$ then $g \in C_{+}$, and thus $g=c+c^{\prime}$ for some $c \in C$ and $c^{\prime} \in \operatorname{lin}(F)$. For $c^{\prime \prime} \in \operatorname{lin}(F)$ nonzero, $g$ is inside the segment with endpoints $c+c^{\prime} \pm c^{\prime \prime}$. Since the endpoints are in $C_{+}$and $G$ is a face of $C_{+}$it contains $c+c^{\prime}+c^{\prime \prime}=g+c^{\prime \prime}$ for all $c^{\prime \prime} \in \operatorname{lin}(F)$. Therefore, $G \supseteq G+\operatorname{lin}(F)$ and $c \in G \cap C$. The latter implies $G \subseteq(G \cap C)+\operatorname{lin}(F)$, and thus the first assertion holds. The second one follows by [15, Corollary 6.6.2]. If $\varepsilon c^{\prime}+(1-\varepsilon) c^{\prime \prime} \in G \cap C$ for $c^{\prime}, c^{\prime \prime} \in C$ and $0<\varepsilon<1$ then $\varepsilon c^{\prime}+(1-\varepsilon) c^{\prime \prime} \in G$ and $c^{\prime}, c^{\prime \prime} \in C_{+}$, and using that $G$ is a face of $C_{+}$it contains $c^{\prime}, c^{\prime \prime}$. It follows that $c^{\prime}, c^{\prime \prime} \in G \cap C$ whence $G \cap C$ is a face of $C$.

By this lemma, if $G$ is the unique face of $C_{+}$that contains $x_{a b}$ in its relative interior then $G \cap C$, denoted in the sequel by $G_{a b}$, is a face of $C$.

Corollary 6.4. $\quad x_{a b} \in \operatorname{ri}\left(G_{a b}\right)+\operatorname{lin}(F)$.

Lemma 6.5. There exist two different parallel hyperplanes $H_{F}, H_{G}$ such that $H_{F} \cap c s(\mu)=F, x_{a b} \in H_{G}, H_{G} \cap C=G_{a b}$ and $H_{G}$ strongly separates $F$ from $s(\mu) \backslash\left(F \cup G_{a b}\right)$.

Proof. The segment with endpoints $a$ and $x_{a b}$ intersects $C_{+}$at its endpoint $x_{a b}$. By [15, Theorem 20.2] applied to this segment and $C_{+}$, there exists a hyperplane $H$ through $x_{a b}$ that separates $a \notin H$ from $C_{+}$. On the other hand, $x_{a b} \in r i(G)$ for a unique face $G$ of $C_{+}$, and thus there exists a supporting hyperplane $K$ of $C_{+}$that intersects this set in $G$. Then $H \cap C_{+} \supseteq G$ because $H$ contains a point from $r i(G)$.

It follows that there exist nonzero $\theta, \vartheta$ such that the hyperplanes $H$ and $K$ are defined by the equations $\left\langle\theta, x-x_{a b}\right\rangle=0$ and $\left\langle\vartheta, x-x_{a b}\right\rangle=0$, respectively. In addition, the scalar products vanish for $x \in G$, are nonnegative for $x \in C_{+}$, $\left\langle\vartheta, x-x_{a b}\right\rangle=0$ with $x \in C_{+}$implies $x \in G$ and $\left\langle\theta, a-x_{a b}\right\rangle<0$. Then the equation $\left\langle\theta+\varepsilon \vartheta, x-x_{a b}\right\rangle=0$ with $\varepsilon>0$ defines a supporting hyperplane $H_{\varepsilon}$ of $C_{+}$that intersects this set in $G$. Taking $\varepsilon$ sufficiently small, $\left\langle\theta+\varepsilon \vartheta, a-x_{a b}\right\rangle<0$, and thus $H_{\varepsilon}$ separates $a \notin H_{\varepsilon}$ and $C_{+}$.

With such a choice of $\varepsilon$, let $H_{G}=H_{\varepsilon}$ and $H_{F}$ be the shift of $H_{G}$ containing $a \notin H_{G}$. By Lemma 6.3, $G=G+\operatorname{lin}(F)$, and then $G \subseteq H_{G}$ implies that $F \subseteq H_{F}$. By the construction of $C$, the points of $s(\mu)$ are either in $F$ or in $C$, and thus $H_{F} \cap \operatorname{cs}(\mu)=F$. By the construction of $H_{\varepsilon}, x_{a b} \in H_{G}$ and $H_{G} \cap C_{+}=G$ which implies $H_{G} \cap C=G_{a b}$. Then the strict separation takes place.

Lemma 6.6. If $E$ is a linear subspace of $\mathbb{R}^{d}, \theta \in E$ and $x \in r i(\mu)+E$ then the function $\vartheta \mapsto\langle\vartheta, x\rangle-\Lambda_{\mu}(\vartheta)$ has a maximizer $\vartheta^{*}$ over the set $\theta+E^{\perp}$. The pm $Q_{\mu, \vartheta^{*}}$ does not depend on the choice of $\vartheta^{*}$ and $x-m\left(Q_{\mu, \vartheta^{*}}\right) \in E$.

Proof. Applying [10, Theorem 3.1] to $\theta+E^{\perp}$ (in the role of $\Xi$, with its barrier cone equal to $E$ ) the function has a unique maximizer over the orthogonal
projection of $\theta+E^{\perp}$ to $E_{x, \mu}=\operatorname{lin}(x-s(\mu))$. By Fact 2.1, $\langle\vartheta, x\rangle-\Lambda_{\mu}(\vartheta)$ remains unchanged when $\vartheta$ moves orthogonally to $E_{x, \mu}$, containing $\operatorname{lin}(\mu)$. It follows that the function has a maximizer $\vartheta^{*}$ over $\theta+E^{\perp}$ and the difference of two such maximizers is orthogonal to $E_{x, \mu}$. By Fact 2.2, $Q_{\mu, \vartheta^{*}}$ is independent of the choice of $\vartheta^{*}$. By [10, Theorem 3.2], $x-m\left(Q_{\mu, \vartheta^{*}}\right)$ is a normal vector of $\theta+E^{\perp}$ at $\vartheta^{*}$, thus belongs to $E . \square$

From now on $G_{a b}$ is abbreviated to $G$.
Corollary 6.7. A maximizer $\vartheta^{*}$ of the function $\vartheta \mapsto\left\langle\vartheta, x_{a b}\right\rangle-\Lambda_{G}(\vartheta)$ with $\vartheta$ in $\Xi=\psi_{F}(a)+\operatorname{lin}(F)^{\perp}$ exists, the mean $m\left(Q_{G, \vartheta^{*}}\right)$ does not depend on its choice and $x_{a b}-m\left(Q_{G, \vartheta^{*}}\right) \in \operatorname{lin}(F)$.

Proof. Lemma 6.6 applies to the restriction of $\mu$ to $G$ in the role of $\mu$, the linear space $E=\operatorname{lin}(F)$, the element $\theta=\psi_{F}(a)$ of $\operatorname{lin}(F)$ and $x=x_{a b}$, which belongs to $r i(G)+\operatorname{lin}(F)$ by Corollary 6.4.

The mean $m\left(Q_{G, \vartheta^{*}}\right)$, independent of $\vartheta^{*}$, is denoted in the sequel by $x_{a b}^{*}$.
Lemma 6.8. $\Lambda_{G}^{*}\left(x_{a b}^{*}\right)+\left\langle\psi_{F}(a), x_{a b}-x_{a b}^{*}\right\rangle=\Psi_{C, \Xi}^{*}\left(x_{a b}\right)$
Proof. By Fact 2.5, applied to $m\left(Q_{G, \vartheta^{*}}\right)=x_{a b}^{*}$, where $\vartheta^{*}$ is a maximizer from Corollary 6.7, $\Lambda_{G}^{*}\left(x_{a b}^{*}\right)=\left\langle\vartheta^{*}, x_{a b}^{*}\right\rangle-\Lambda_{G}\left(\vartheta^{*}\right)$. Since $\vartheta^{*}-\psi_{F}(a)$ is orthogonal to $\operatorname{lin}(F)$, containing $x_{a b}-x_{a b}^{*}$,

$$
\Lambda_{G}^{*}\left(x_{a b}^{*}\right)+\left\langle\psi_{F}(a), x_{a b}-x_{a b}^{*}\right\rangle=\left\langle\vartheta^{*}, x_{a b}\right\rangle-\Lambda_{G}\left(\vartheta^{*}\right) \geqslant\left\langle\vartheta, x_{a b}\right\rangle-\Lambda_{C}(\vartheta), \quad \vartheta \in \Xi
$$

using $\Lambda_{C} \geqslant \Lambda_{G}$. Maximizing over $\vartheta, \Psi_{C, \Xi}^{*}\left(x_{a b}\right)$ emerges on the right.
On the other hand, Lemma 6.5 implies that there exists a nonzero $\tau$ orthogonal to $\operatorname{lin}(F)$ such that $\left\langle\tau, x-x_{a b}\right\rangle \leqslant 0$ holds for $x \in C$, with the equality if and only if $x \in G=G_{a b}$. Hence, $\vartheta^{*}+t \tau \in \Xi, t \in \mathbb{R}$, and

$$
\Psi_{C, \Xi}^{*}\left(x_{a b}\right) \geqslant\left\langle\vartheta^{*}+t \tau, x_{a b}\right\rangle-\Lambda_{C}\left(\vartheta^{*}+t \tau\right)=-\ln \sum_{x \in s(\mu) \backslash F} e^{\left\langle\vartheta^{*}+t \tau, x-x_{a b}\right\rangle} \mu(x)
$$

where $\left\langle\vartheta^{*}, x_{a b}\right\rangle-\Lambda_{G}\left(\vartheta^{*}\right)$ emerges on the right when $t$ grows to $+\infty$.
Let $b_{\varepsilon}$ abbreviate $a+\varepsilon t_{a b}(b-a)$, equal to $a+\varepsilon\left(x_{a b}-a\right)$. The convex hull of $F \cup G$ is denoted by $A$.

Lemma 6.9. If $\varepsilon>0$ is sufficiently small then $b_{\varepsilon} \in r i(A)$.
Proof. By Corollary 6.4, $x_{a b}=c+t\left(a^{\prime}-a\right)$ with $c \in r i(G), a^{\prime} \in F$ and $t \geqslant 0$. Then

$$
b_{\varepsilon}=a+\varepsilon\left(c+t\left(a^{\prime}-a\right)-a\right)=(1-\varepsilon)\left[\frac{\varepsilon t}{1-\varepsilon} a^{\prime}+\left(1-\frac{\varepsilon t}{1-\varepsilon}\right) a\right]+\varepsilon c
$$

For small $\varepsilon>0$ the bracket is a convex combination of $a^{\prime}$ and $a \in r i(F)$ whence it belongs to $r i(F)$. Then, $b_{\varepsilon}$ is a convex combination of points from $r i(F)$ and $r i(G)$, and the assertion follows by [15, Theorem 6.9].

By Lemma 6.9 , if $\varepsilon>0$ is sufficiently small then $\vartheta_{\varepsilon}=\psi_{A}\left(b_{\varepsilon}\right)$ is well-defined. Denote the means of $Q_{F, \vartheta_{\varepsilon}}$ and $Q_{G, \vartheta_{\varepsilon}}$ by $c_{F, \varepsilon}$ and $c_{G, \varepsilon}$, respectively. By Lemma 6.5, two parallel hyperplanes contain the pairs $c_{F, \varepsilon}, a$ and $c_{G, \varepsilon}, x_{a b}$, and a geometric argument implies that $b_{\varepsilon}=(1-\varepsilon) a+\varepsilon x_{a b}$ equals $m\left(Q_{A, \vartheta_{\varepsilon}}\right)=(1-\varepsilon) c_{F, \varepsilon}+\varepsilon c_{G, \varepsilon}$. In turn,

$$
\begin{equation*}
(1-\varepsilon)\left(c_{F, \varepsilon}-a\right)=\varepsilon\left(x_{a b}-c_{G, \varepsilon}\right) . \tag{9}
\end{equation*}
$$

Lemma 6.10. If $\varepsilon$ decreases to zero then $c_{F, \varepsilon} \rightarrow a$ and $c_{G, \varepsilon} \rightarrow x_{a b}^{*}$.
Proof. The first convergence is a consequence of (9) and $c_{G, \varepsilon} \in r i(G)$, which is a bounded set. It implies that $\psi_{F}\left(c_{F, \varepsilon}\right)$, which is the projection of $\vartheta_{\varepsilon}$ to $\operatorname{lin}(F)$ by Fact 2.2, converges to $\psi_{F}(a)$. Hence, for a maximizer $\vartheta^{*}$ from Corollary 6.7

$$
\begin{aligned}
& D\left(Q_{G, \vartheta_{\varepsilon}} \| Q_{G, \vartheta^{*}}\right)+D\left(Q_{G, \vartheta^{*}} \| Q_{G, \vartheta_{\varepsilon}}\right)=\left\langle\vartheta_{\varepsilon}-\vartheta^{*}, m\left(Q_{G, \vartheta_{\varepsilon}}\right)-m\left(Q_{G, \vartheta^{*}}\right)\right\rangle \\
&=\left\langle\vartheta_{\varepsilon}-\vartheta^{*}, c_{G, \varepsilon}-x_{a b}^{*}\right\rangle=\left\langle\psi_{F}\left(c_{F, \varepsilon}\right)-\psi_{F}(a), c_{G, \varepsilon}-x_{a b}^{*}\right\rangle \rightarrow 0 .
\end{aligned}
$$

The last equality holds because $\vartheta_{\varepsilon}-\psi_{F}\left(c_{F, \varepsilon}\right)$ and $\vartheta^{*}-\psi_{F}(a)$ are orthogonal to $\operatorname{lin}(F)$ and $c_{G, \varepsilon}-x_{a b}^{*} \in \operatorname{lin}(F)$. Note that $c_{G, \varepsilon}-x_{a b}^{*}$ is sum of $c_{G, \varepsilon}-x_{a b}$, proportional to $a-c_{F, \varepsilon} \in \operatorname{lin}(F)$ by (9), and $x_{a b}-x_{a b}^{*}$, belonging to lin $(F)$ by Corollary 6.7. By Pinsker inequality, $Q_{G, \vartheta_{\varepsilon}} \rightarrow Q_{G, \vartheta^{*}}$ in variation distance which, in turn, implies the convergence of means $c_{G, \varepsilon} \rightarrow x_{a b}^{*}$.

Let $\theta_{\varepsilon}$ denote the orthogonal projection of $\vartheta_{\varepsilon}$ to $\operatorname{lin}(F)+\operatorname{lin}(G)$.

Corollary 6.11. If $\varepsilon$ decreases to 0 then $\theta_{\varepsilon}$ converges.
Proof. By Fact 2.2, $\psi_{F}\left(c_{F, \varepsilon}\right)$ is the orthogonal projection of $\vartheta_{\varepsilon}$ to $\operatorname{lin}(F)$, converging by Lemma 6.10. The arguments work also when $F$ is replaced by $G$.

Lemma 6.12. $\Lambda_{\mu}^{*}\left(b_{\varepsilon}\right)=\Lambda_{A}^{*}\left(b_{\varepsilon}\right)+o(\varepsilon)$
Proof. The assertion is trivial if $B=s(\mu) \backslash A$ is empty. Otherwise, Lemma 6.5 implies existence of a nonzero $\tau$ such that the function $x \mapsto\langle\tau, x\rangle$ equals a constant $s_{F}$ on $F$, a constant $s_{G}<s_{F}$ on $G$ and is upper bounded by $s_{B}<s_{G}$ on $B=s(\mu) \backslash A$. Scaling $\tau$ if necessary, $s_{F}-s_{G}=1$. Let

$$
r_{\varepsilon}=\Lambda_{G}\left(\theta_{\varepsilon}\right)-\Lambda_{F}\left(\theta_{\varepsilon}\right)+\ln \frac{1-\varepsilon}{\varepsilon}
$$

Since $\tau$ is orthogonal to $\operatorname{lin}(F)+\operatorname{lin}(G)$ the means of $Q_{F, \theta_{\varepsilon}+r_{\varepsilon} \tau}$ and $Q_{G, \theta_{\varepsilon}+r_{\varepsilon} \tau}$ are equal to $c_{F, \varepsilon}$ and $c_{G, \varepsilon}$, respectively. It follows from

$$
\begin{equation*}
m\left(Q_{A, \theta}\right)=e^{\Lambda_{F}(\theta)-\Lambda_{A}(\theta)} m\left(Q_{F, \theta}\right)+e^{\Lambda_{G}(\theta)-\Lambda_{A}(\theta)} m\left(Q_{G, \theta}\right), \quad \theta \in \mathbb{R}^{d} \tag{10}
\end{equation*}
$$

that the mean of $Q_{A, \theta_{\varepsilon}+r_{\varepsilon} \tau}$ equals $(1-\delta) c_{F, \varepsilon}+\delta c_{G, \varepsilon}$ where

$$
\begin{aligned}
\ln \frac{1-\delta}{\delta} & =\Lambda_{F}\left(\theta_{\varepsilon}+r_{\varepsilon} \tau\right)-\Lambda_{G}\left(\theta_{\varepsilon}+r_{\varepsilon} \tau\right) \\
& =r_{\varepsilon}\left(s_{F}-s_{G}\right)+\Lambda_{F}\left(\theta_{\varepsilon}\right)-\Lambda_{G}\left(\theta_{\varepsilon}\right)=\ln \frac{1-\varepsilon}{\varepsilon}
\end{aligned}
$$

by the choice of $r_{\varepsilon}$. Therefore, $\delta=\varepsilon$ and $m\left(Q_{A, \theta_{\varepsilon}+r_{\varepsilon} \tau}\right)$ equals the mean $b_{\varepsilon}$ of $Q_{A, \vartheta_{\varepsilon}}$. This implies

$$
\Lambda_{\mu}^{*}\left(b_{\varepsilon}\right) \geqslant\left\langle\theta_{\varepsilon}+r_{\varepsilon} \tau, b_{\varepsilon}\right\rangle-\Lambda_{\mu}\left(\theta_{\varepsilon}+r_{\varepsilon} \tau\right)=\Lambda_{A}^{*}\left(b_{\varepsilon}\right)-\Lambda_{\mu}\left(\theta_{\varepsilon}+r_{\varepsilon} \tau\right)+\Lambda_{A}\left(\theta_{\varepsilon}+r_{\varepsilon} \tau\right)
$$

using Fact 2.5. Here,

$$
\Lambda_{A}\left(\theta_{\varepsilon}+r_{\varepsilon} \tau\right)=\ln \left[e^{r_{\varepsilon} s_{F}+\Lambda_{F}\left(\theta_{\varepsilon}\right)}+e^{r_{\varepsilon} s_{G}+\Lambda_{G}\left(\theta_{\varepsilon}\right)}\right]
$$

and

$$
\Lambda_{\mu}\left(\theta_{\varepsilon}\right) \leqslant \ln \left[e^{\Lambda_{A}\left(\theta_{\varepsilon}+r_{\varepsilon} \tau\right)}+e^{r_{\varepsilon} s_{B}+\Lambda_{B}\left(\theta_{\varepsilon}\right)}\right] .
$$

Hence, $\Lambda_{\mu}^{*}\left(b_{\varepsilon}\right)-\Lambda_{A}^{*}\left(b_{\varepsilon}\right)$ is at least

$$
\begin{aligned}
-\ln \left[1+\frac{e^{r_{\varepsilon} s_{B}+\Lambda_{B}\left(\theta_{\varepsilon}\right)}}{e^{r_{\varepsilon} s_{F}+\Lambda_{F}\left(\theta_{\varepsilon}\right)}+e^{r_{\varepsilon} s_{G}+\Lambda_{G}\left(\theta_{\varepsilon}\right)}}\right] & \geqslant-\frac{e^{r_{\varepsilon}\left(s_{B}-s_{G}\right)+\Lambda_{B}\left(\theta_{\varepsilon}\right)-\Lambda_{G}\left(\theta_{\varepsilon}\right)}}{e^{r_{\varepsilon}+\Lambda_{F}\left(\theta_{\varepsilon}\right)-\Lambda_{G}\left(\theta_{\varepsilon}\right)}+1} \\
& =-\varepsilon e^{r_{\varepsilon}\left(s_{B}-s_{G}\right)+\Lambda_{B}\left(\theta_{\varepsilon}\right)-\Lambda_{G}\left(\theta_{\varepsilon}\right)}
\end{aligned}
$$

due to the choice of $r_{\varepsilon}$. By Corollary $6.11, \theta_{\varepsilon}$ converges whence $e^{-r_{\varepsilon}}$ is of the order $O(\varepsilon)$. In turn, $\varepsilon e^{r_{\varepsilon}\left(s_{B}-s_{G}\right)}$ is of the order $o(\varepsilon)$, on account of $s_{B}-s_{G}<0$. Therefore, a lower bound to $\Lambda_{\mu}^{*}\left(b_{\varepsilon}\right)-\Lambda_{A}^{*}\left(b_{\varepsilon}\right)$ is of the order $o(\varepsilon)$. The assertion follows by mentioning that $\Lambda_{\mu}^{*} \leqslant \Lambda_{A}^{*}$.

Proof of Theorem 3.1. By Lemma 6.12 and Fact 2.8, it suffices to prove that

$$
\Lambda_{A}^{*}\left(b_{\varepsilon}\right)=\Lambda_{F}^{*}(a)+h(\varepsilon)+\varepsilon\left[\Psi_{C, \Xi}^{*}\left(x_{a b}\right)-\Lambda_{F}^{*}(a)\right]+o(\varepsilon) .
$$

It follows from Fact 2.5, $b_{\varepsilon}=(1-\varepsilon) c_{F, \varepsilon}+\varepsilon c_{G, \varepsilon}$ and (10) that

$$
\begin{aligned}
\Lambda_{A}^{*}\left(b_{\varepsilon}\right)= & \left\langle\vartheta_{\varepsilon}, b_{\varepsilon}\right\rangle-(1-\varepsilon+\varepsilon) \Lambda_{A}\left(\vartheta_{\varepsilon}\right) \\
= & (1-\varepsilon)\left[\left\langle\vartheta_{\varepsilon}, c_{F, \varepsilon}\right\rangle-\Lambda_{F}\left(\vartheta_{\varepsilon}\right)+\ln (1-\varepsilon)\right] \\
& \quad+\varepsilon\left[\left\langle\vartheta_{\varepsilon}, c_{G, \varepsilon}\right\rangle-\Lambda_{G}\left(\vartheta_{\varepsilon}\right)+\ln \varepsilon\right] \\
= & h(\varepsilon)+(1-\varepsilon) \Lambda_{F}^{*}\left(c_{F, \varepsilon}\right)+\varepsilon \Lambda_{G}^{*}\left(c_{G, \varepsilon}\right) .
\end{aligned}
$$

By Lemma 6.10 and (9), the norm of $c_{F, \varepsilon}-a \in \operatorname{lin}(F)$ is of the order $o(\varepsilon)$. Then, using Fact 2.6,

$$
\Lambda_{F}^{*}\left(c_{F, \varepsilon}\right)=\Lambda_{F}^{*}(a)+\left\langle\psi_{F}(a), c_{F, \varepsilon}-a\right\rangle+o(\varepsilon)
$$

where the scalar product equals $\varepsilon\left\langle\psi_{F}(a), x_{a b}-c_{G, \varepsilon}\right\rangle+o(\varepsilon)$ by (9). Therefore,

$$
\Lambda_{A}^{*}\left(b_{\varepsilon}\right)=\Lambda_{F}^{*}(a)+h(\varepsilon)+\varepsilon\left[\Lambda_{G}^{*}\left(c_{G, \varepsilon}\right)+\left\langle\psi_{F}(a), x_{a b}-c_{G, \varepsilon}\right\rangle-\Lambda_{F}^{*}(a)\right]+o(\varepsilon)
$$

This holds also when $c_{G, \varepsilon}$ is replaced by $x_{a b}^{*}$ because $c_{G, \varepsilon} \rightarrow x_{a b}^{*}$ by Lemma 6.10 and $\Lambda_{G}^{*}$ is continuous on $G$. Using Lemma 6.8, the assertion follows.

Proof of Lemma 4.2. Let $P_{\varepsilon}=P+\varepsilon(R-P)$. Assuming first $\varepsilon>0$,

$$
D\left(P_{\varepsilon} \| \nu\right)=\sum_{z \in s(R) \backslash s(P)} \varepsilon R(z) \ln \frac{\varepsilon R(z)}{\nu(z)}+\sum_{z \in s(P)} P_{\varepsilon}(z) \ln \frac{P_{\varepsilon}(z)}{\nu(z)} .
$$

In the second sum,

$$
\ln \frac{P_{\varepsilon}(z)}{\nu(z)}=\ln \frac{P(z)}{\nu(z)}+\ln \left[1+\varepsilon \frac{R(z)-P(z)}{P(z)}\right]=\ln \frac{P(z)}{\nu(z)}+\varepsilon \frac{R(z)-P(z)}{P(z)}+o(\varepsilon) .
$$

Hence,

$$
\begin{aligned}
D\left(P_{\varepsilon} \| \nu\right)=r \varepsilon \ln \varepsilon & +\varepsilon \sum_{z \in s(R) \backslash s(P)} R(z) \ln \frac{R(z)}{\nu(z)} \\
& +D(P \| \nu)+\varepsilon \sum_{z \in s(P)}[R(z)-P(z)]\left[1+\ln \frac{P(z)}{\nu(z)}\right]+o(\varepsilon)
\end{aligned}
$$

This and

$$
\varepsilon \sum_{z \in s(P)}[R(z)-P(z)]=-r \varepsilon=r(1-\varepsilon) \ln (1-\varepsilon)+o(\varepsilon)
$$

imply the first assertion. If $r=0$ the argumentation goes through also for $\varepsilon \leqslant 0$, omitting corresponding terms.

Proof of Lemma 4.4. First, it is shown that there exists $c \in C$ such that $t_{a c}<1$. The assumption implies $a \in \operatorname{aff}(C)+\operatorname{lin}(F)$. Then $a=t c^{\prime}+(1-t) c^{\prime \prime}+b^{\prime}$ for some $c^{\prime}, c^{\prime \prime} \in C, b^{\prime} \in \operatorname{lin}(F)$ and $t \in \mathbb{R}$. By Lemma 6.2, $a \notin C_{+}$whence $t$ is not between 0 and 1 . Changing the roles of $c^{\prime}$ and $c^{\prime \prime}$ if necessary it is possible to assume that $t>1$. Let $c=c^{\prime \prime}$. It follows that $a+\frac{t-1}{t}(c-a)$ equals $c^{\prime}+\frac{1}{t} b^{\prime}$ which belongs to $C_{+}$. Hence, $t_{a c} \leqslant \frac{t-1}{t}<1$. Obviously $c=m\left(Q_{f}\right)$ for some pm $Q$ concentrated on $Z \backslash f^{-1}(F)$. Then $f^{-1}(F) \supseteq s(P)$ implies $Q(Z \backslash s(P))=1$, and the derivative in the direction $Q-P$ is $+\infty$, by Theorem 4.1.

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