# Some Properties of the Hellinger Transform and Its Application in Classification Problems 

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#### Abstract

The Hellinger transform for stochastic processes is studied. The Hellinger transform is calculated for multidimensional Gaussian autoregressive processes and for Ornstein-Uhlenbeck processes.


Keywords-Hellinger transform, Radon-Nikodym derivative, Multidimensional Gaussian autoregressive process, Ornstein-Uhlenbeck process.

## 1. INTRODUCTION AND NOTATION

In this paper, the following general classification problem is studied. Let $(\Omega, \mathcal{F}, P)$ be a probability space and $\theta: \Omega \rightarrow\{1, \ldots, N\}$ a random variable (parameter). Let

$$
\begin{equation*}
\pi_{i}=P(\theta=i)>0, \quad i=1, \ldots, N \tag{1}
\end{equation*}
$$

denote the a priori probabilities of classes.

$$
\begin{equation*}
P_{i}(\cdot)=P(\cdot \mid \theta=i), \quad i=1, \ldots, N, \tag{2}
\end{equation*}
$$

denotes the probability measure in the case of the $i^{\text {th }}$ class.
Let $\mu$ be a dominating (probability) measure for $P_{1}, \ldots, P_{N}$ and let $X_{i}$ denote the density of $P_{i}$ with respect to $\mu$. The Hellinger transform of $P_{1}, \ldots, P_{N}$ is defined as follows:

$$
\begin{equation*}
\mathbb{H}_{\underline{\alpha}}\left(P_{1}, \ldots, P_{N}\right)=E_{\mu}\left\{X_{1}^{\alpha_{1}} \cdots X_{N}^{\alpha_{N}}\right\}, \tag{3}
\end{equation*}
$$

where $\underline{\alpha}=\left(\alpha_{1}, \ldots, \alpha_{N}\right)$ is a fixed vector with $\alpha_{1} \geq 0, \ldots, \alpha_{N} \geq 0, \alpha_{1}+\cdots+\alpha_{N}=1$ and $E_{\mu}$ denotes the expectation with respect to $\mu$ (see [1]). The Hellinger transform is one of the measures of affinity among several distributions used in the literature (see, e.g., $[2,3]$ ). It is easily seen that

$$
\begin{equation*}
0 \leq \mathbb{H}_{\underline{\alpha}}\left(P_{1}, \ldots, P_{N}\right) \leq 1 \tag{4}
\end{equation*}
$$

[^0]and $\mathbb{H}_{\underline{\alpha}}=1$ holds iff measures $P_{1}, \ldots, P_{N}$ coincide. Furthermore, $\mathbb{H}_{\underline{\alpha}}$ is close to zero iff the 'distance' of $P_{1}, \ldots, P_{N}$ is large in some sense. Therefore $\mathbb{H}_{\underline{\alpha}}$ can be used to measure the possibility of discrimination among several distributions. We mention that many deep investigations are devoted to various 'distances' between two probability measures and their application in hypothesis testing (see, e.g., [4-7]), but it seems that the case of several distributions is not completely known. In [8], the limit (when the number of observations tends to infinity) of the Hellinger transform of Gaussian autoregressive processes is given in terms of spectral densities. We remark that Michálek [9] calculated the asymptotic Rényi's rate for Gaussian processes.

In this paper, upper and lower bounds are given for $\mathbb{H}_{\underline{\alpha}}$ in the case of a filtered probability space (Propositions 2.8 and 2.9). An explicit formula for $\mathbb{H}_{\underline{\alpha}}$ is calculated in the case of multidimensional Gaussian AR processes in Section 3. In Section 4, the Hellinger transform of one-dimensional Ornstein-Uhlenbeck processes is given. In Section 5, the connection between $\mathbb{H}_{\underline{\alpha}}$ and the probability of misclassification is considered.

## 2. BASIC PROPERTIES OF HELLINGER TRANSFORMS

Let us introduce the notation

$$
\begin{equation*}
\underline{X}^{\underline{\alpha}}=X_{1}^{\alpha_{1}} \cdots X_{N}^{\alpha_{N}}, \tag{5}
\end{equation*}
$$

where $\underline{X}=\left(X_{1}, \ldots, X_{N}\right)$ is a vector with nonnegative coordinates and $\underline{\alpha}=\left(\alpha_{1}, \ldots, \alpha_{N}\right)$ is an arbitrary vector.
Definition 2.1. (See [1].) Let $\underline{\alpha}=\left(\alpha_{1}, \ldots, \alpha_{N}\right)$, where $\alpha_{i} \geq 0$ for every $i$ and $\sum_{i=1}^{N} \alpha_{i}=1$. The Hellinger transform of the probability measures $P_{1}, \ldots, P_{N}$ is defined by the formula

$$
\begin{equation*}
\mathbb{H}_{\underline{\alpha}}\left(P_{1}, \ldots, P_{N}\right)=E_{\mu}\left\{X_{1}^{\alpha_{1}} \cdots X_{N}^{\alpha_{N}}\right\}=E_{\mu}\{\underline{X} \underline{\underline{\alpha}}\}, \tag{6}
\end{equation*}
$$

where $\mu$ is a dominating measure for $P_{1}, \ldots, P_{N}$ and $X_{i}=\frac{d P_{i}}{d \mu}$.
For $\underline{\alpha}_{0}=(1 / N, \ldots, 1 / N), \mathbb{H}_{\underline{\alpha}_{0}}\left(P_{1}, \ldots, P_{N}\right)$ is the Matusita affinity or Bhattacharyya coefficient (see [2,3]). For $N=2$, it is the well-known Hellinger integral. The Hellinger transform is a special case of the $f$-affinity for $f\left(s_{1}, \ldots, s_{N}\right)=\prod_{i=1}^{N} s_{i}^{\alpha_{i}}$.

Definition 2.2. Let $f(\underline{s})=f\left(s_{1}, \ldots, s_{N}\right), s_{i} \geq 0$, for $i=1, \ldots, N$, be a continuous, concave, homogeneous function. The $f$-affinity of $P_{1}, \ldots, P_{N}$ is defined by

$$
\begin{equation*}
A_{f}\left(P_{1}, \ldots, P_{N}\right)=\int f\left(X_{1}, \ldots, X_{N}\right) d \mu \tag{7}
\end{equation*}
$$

As $f$ is homogeneous, $A_{f}\left(P_{1}, \ldots, P_{N}\right)$ does not depend on the choice of the dominating measure $\mu$. The $f$-affinity is nothing else but the negative of the $f$-dissimilarity of Györfi and Nemetz [3].
Proposition 2.3. (See [3].) Let $\mathcal{F}_{s}$ be a $\sigma$-subalgebra of $\mathcal{F}$ and let $P_{1, s}, \ldots, P_{N, s}$ denote the restriction of $P_{1}, \ldots, P_{N}$ on $\mathcal{F}_{s}$. Then

$$
\begin{equation*}
A_{f}\left(P_{1}, \ldots, P_{N}\right) \leq A_{f}\left(P_{1, s}, \ldots, P_{N, s}\right) \tag{8}
\end{equation*}
$$

For the proof, let $\mu=\left(P_{1}+\cdots+P_{N}\right) / N$ and use the Jensen inequality.
Corollary 2.4.

$$
\begin{equation*}
\mathbb{H}_{\underline{\underline{\alpha}}}\left(P_{1}, \ldots, P_{N}\right) \leq \mathbb{H}_{\underline{\alpha}}\left(P_{1, s}, \ldots, P_{N, s}\right) . \tag{9}
\end{equation*}
$$

Throughout this section, $\mathcal{F}_{0} \subseteq \mathcal{F}_{1} \subseteq \cdots$ denotes a nondecreasing sequence of $\sigma$-subalgebras of $\mathcal{F}$ for which $\sigma\left(\cup_{n} \mathcal{F}_{n}\right)=\mathcal{F}$. Let $P_{i, n}$, (respectively, $\left.\mu_{n}\right)$ be the restriction of $P_{i}$, (respectively, $\mu$ ) to $\mathcal{F}_{n}$ and $X_{i, n}=\frac{d P_{i, n}}{d \mu_{n}}, \underline{X}_{n}=\left(X_{1, n}, \ldots, X_{N, n}\right)(i=1, \ldots, N, n=0,1,2, \ldots) . T$ will
denote a stopping time with respect to $\left\{\mathcal{F}_{n}, n=0,1, \ldots\right\}$. The following generalization of the approximation theorem of Liese and Vajda [6] holds for Hellinger transforms.

Proposition 2.5. With the above notation,

$$
\begin{equation*}
\lim _{n \rightarrow \infty} A_{f}\left(P_{1, n}, \ldots, P_{N, n}\right)=A_{f}\left(P_{1}, \ldots, P_{N}\right) \tag{10}
\end{equation*}
$$

For the proof, let $\mu$ be the same as in the previous proof and use the supermartingale convergence theorem for $f\left(\underline{X}_{n}\right)$.

Corollary 2.6. If $n \rightarrow \infty$, then

$$
\begin{equation*}
\mathbb{H}_{\underline{\alpha}}\left(P_{1, n}, \ldots, P_{N, n}\right) \downarrow \mathbb{H}_{\underline{\alpha}}\left(P_{1}, \ldots, P_{N}\right) . \tag{11}
\end{equation*}
$$

We shall use the following multiplicative decomposition of a nonnegative supermartingale.
Lemma 2.7. Let ( $S_{n}, \mathcal{F}_{n}, n=0,1, \ldots$ ) be a nonnegative supermartingale (with $\mathcal{F}_{0}=\{\emptyset, \Omega\}$ ). Then there exist a martingale ( $M_{n}, \mathcal{F}_{n}$ ) and a nonnegative, nonincreasing, predictable sequence $\left(G_{n}, \mathcal{F}_{n}\right)$ such that

$$
\begin{equation*}
S_{n}=M_{n} G_{n}, \quad n=0,1, \ldots \tag{12}
\end{equation*}
$$

Proof. Let $M_{0}=1, G_{0}=S_{0}$,

$$
\begin{equation*}
G_{n}=\frac{E\left\{S_{n} \mid \mathcal{F}_{n-1}\right\}}{S_{n-1}} G_{n-1}, \quad n \geq 1 \tag{13}
\end{equation*}
$$

(here we used the convention $1 / 0=0$ ). Let $M_{n}=S_{n} / G_{n}$ if $G_{n} \neq 0$ and $M_{n}=M_{n-1}$ if $G_{n}=0$. Then $M_{n}$ is a martingale. Since $G_{n}(\omega)=0$ implies that $S_{n}(\omega)=0$, then (12) is satisfied.

In the sequel, we use the notation

$$
\begin{equation*}
S_{n}(\underline{\alpha})=\underline{X} \underline{\underline{\alpha}} . \tag{14}
\end{equation*}
$$

Since $S_{n}(\underline{\alpha})$ is a nonnegative supermartingale, it admits a multiplicative decomposition

$$
\begin{equation*}
S_{n}(\underline{\alpha})=M_{n}(\underline{\alpha}) G_{n}(\underline{\alpha}), \tag{15}
\end{equation*}
$$

where $G_{n}(\underline{\alpha})$ is defined analogously to (13).
We derive inequalities between $\mathbb{H}_{\underline{\alpha}}\left(P_{1}, \ldots, P_{N}\right)$ and $G_{n}\left(\underline{\alpha}^{\prime}\right)$ for an appropriate $\underline{\alpha}^{\prime}$. Let $\underline{\alpha}=$ $\left(\alpha_{1}, \ldots, \alpha_{N}\right)$ be fixed. For any $i(1 \leq i \leq N)$ and $\gamma$, we associate with $\underline{\alpha}$ a new vector

$$
\begin{equation*}
[\underline{\alpha}, i, \gamma]=\left(\gamma \alpha_{1}, \ldots, \gamma \alpha_{i-1}, 1-\gamma\left(1-\alpha_{i}\right), \gamma \alpha_{i+1}, \ldots, \gamma \alpha_{N}\right) . \tag{16}
\end{equation*}
$$

Proposition 2.8. Let us assume that

$$
\begin{equation*}
1<\gamma_{i}<\min \left\{\frac{1}{\alpha_{i}}, \frac{1}{1-\alpha_{i}}\right\} \tag{17}
\end{equation*}
$$

is satisfied for $\gamma_{i}$ for every $i=1, \ldots, N$. Then for every stopping time $T$,

$$
\begin{equation*}
\mathbb{H}_{\underline{\alpha}}\left(P_{1, T}, \ldots, P_{N, T}\right) \leq \min _{1 \leq i \leq N}\left(E_{P_{i}}\left\{G_{T}^{1 /\left(\gamma_{i}-1\right)}\left(\left[\underline{\alpha}, i, \gamma_{i}\right]\right)\right\}\right)^{\left(\gamma_{i}-1\right) / \gamma_{i}} . \tag{18}
\end{equation*}
$$

Proof. We prove (18) for $i=1$. Let $\gamma=\gamma_{1}$ satisfy (17). By the multiplicative decomposition, we can write

$$
\begin{equation*}
\left(X_{1, T}^{\alpha_{1}} \cdots X_{N, T}^{\alpha_{N}}\right)^{\gamma}=X_{1, T}^{\gamma-1} M_{T}([\alpha, 1, \gamma]) G_{T}([\underline{\alpha}, 1, \gamma]) . \tag{19}
\end{equation*}
$$

Hence

$$
\begin{equation*}
E_{\mu}\left\{X_{1, T}^{\alpha_{1}} \cdots X_{N, T}^{\alpha_{N}}\right\}=E_{\mu}\left\{X_{1 T}^{(\gamma-1) / \gamma} G_{T}^{1 / \gamma}([\underline{\alpha}, 1, \gamma]) M_{T}^{1 / \gamma}([\underline{\alpha}, 1, \gamma])\right\} . \tag{20}
\end{equation*}
$$

Applying Hölders inequality with $p=\gamma /(\gamma-1), q=\gamma$, we get

$$
\begin{equation*}
E_{\mu}\left\{X_{1, T}^{\alpha_{1}} \cdots X_{N, T}^{\alpha_{N}}\right\} \leq\left[E_{\mu}\left\{M_{T}([\underline{\alpha}, 1, \gamma])\right\}\right]^{1 / \gamma}\left[E_{\mu}\left\{X_{1, T} G^{1 /(\gamma-1)}([\underline{\alpha}, 1, \gamma])\right\}\right]^{(\gamma-1) / \gamma} \tag{21}
\end{equation*}
$$

One can easily derive that $E_{\mu}\left\{X_{1, T} \xi\right\}=E_{P_{1}}\{\xi\}$ for every $\mathcal{F}_{T}$-measurable $\xi$. By the proof of Lemma 2.7, we have $E_{\mu}\left\{M_{n}\right\}=1$ for the martingale part of the multiplicative decomposition, so $E_{\mu}\left\{M_{T}\right\} \leq 1$. Therefore, (21) implies that

$$
\begin{equation*}
\mathbb{H}_{\underline{\alpha}}\left(P_{1, T}, \ldots, P_{N, T}\right) \leq\left[E_{P_{1}}\left\{G_{T}^{1 /(\gamma-1)}([\underline{\alpha}, 1, \gamma])\right\}\right]^{(\gamma-1) / \gamma} \tag{22}
\end{equation*}
$$

To prove the following inequality, we need the inverse Hölder inequality: if $\xi$ and $\eta$ are positive random variables, $0<p<1,1 / q+1 / p=1$, then

$$
\begin{equation*}
E\{\xi \eta\} \geq\left(E\left\{\xi^{p}\right\}\right)^{1 / p}\left(E\left\{\eta^{q}\right\}\right)^{1 / q} \tag{23}
\end{equation*}
$$

Proposition 2.9. Suppose that the probability measures $P_{1}, \ldots, P_{N}$ are equivalent and let $\mu$ be chosen such that $\mu$ and $P_{i}$ are equivalent. Then for every stopping time $T$

$$
\begin{equation*}
\mathbb{H}_{\underline{\alpha}}\left(P_{1, T}, \ldots, P_{N, T}\right) \geq \max _{1 \leq i \leq N}\left[E_{P_{i}}\left\{G_{T}^{1 /\left(\gamma_{i}-1\right)}\left(\left[\underline{\alpha}, i, \gamma_{i}\right]\right)\right\}\right]^{\left(\gamma_{i}-1\right) / \gamma_{i}}, \tag{24}
\end{equation*}
$$

where $0<\gamma_{i}<1$ for every $i=1, \ldots, N$.
Proof. First we remark that equivalence of $\mu$ and $P_{i}$ implies that $X_{i, n}>0$. Therefore $G_{n}(\cdot)>0$ and $M_{n}(\cdot)>0$. Now, let $T$ be a bounded stopping time. Let $\gamma=\gamma_{1}$ be fixed. Easy calculation and the inverse Hölder inequality (with $p=\gamma$ ) show that

$$
\begin{align*}
\mathbb{H}_{\underline{\alpha}}\left(P_{1, T}, \ldots, P_{N, T}\right) & =E_{\mu}\left\{\left(X_{1, T}\right)^{(\gamma-1) / \gamma} G_{T}^{1 / \gamma}([\underline{\alpha}, 1, \gamma]) M_{T}^{1 / \gamma}([\underline{\alpha}, 1, \gamma])\right\}  \tag{25}\\
& \geq\left[E_{\mu}\left\{X_{1, T} G_{T}^{1 /(\gamma-1)}([\underline{\alpha}, 1, \gamma])\right\}\right]^{(\gamma-1) / \gamma}\left[E_{\mu}\left\{M_{T}([\underline{\alpha}, 1, \gamma])\right\}\right]^{1 / \gamma}  \tag{26}\\
& =\left[E_{P_{1}}\left\{G_{T}^{1 /(\gamma-1)}([\underline{\alpha}, 1, \gamma])\right\}\right]^{(\gamma-1) / \gamma} . \tag{27}
\end{align*}
$$

This proves (24) for bounded stopping times. If $T$ is an arbitrary stopping time, then approximate it by $T \wedge n$ and use Corollary 2.6.
Remark 2.10. $G_{n}(\underline{\alpha})$ can easily be calculated for Markov processes. Let $y_{1}, y_{2}, \ldots$ be a homogeneous Markov process with transition density function $p_{i}(\cdot, \cdot)$ and initial density function $p_{i}(\cdot)$ in the case $\theta=i$. Then

$$
\begin{equation*}
G_{n}(\underline{\alpha})=\prod_{l=2}^{n}\left[\int_{-\infty}^{+\infty} \prod_{i=1}^{N} p_{i}^{\alpha_{i}}\left(y_{l-1}, v_{l}\right) d v_{l}\right] \cdot\left(\prod_{i=1}^{N} p_{i}^{\alpha_{i}}\left(y_{1}\right)\right) \cdot \int_{-\infty}^{+\infty} \prod_{i=1}^{N} p_{i}^{\alpha_{i}}\left(v_{1}\right) d v_{1} \tag{28}
\end{equation*}
$$

If, moreover, one of the measures $P_{i}, i=1, \ldots, N$, say $P_{1}$, is such that $y_{1}, y_{2}, \ldots$ are independent, then inequality (18) is quite simple because if we take expectation of $G_{n}$ with respect to $P_{1}$ and use (28) then the expectation sign $E_{P_{1}}$ can be interchanged with the product sign with respect to $l$.

## 3. THE HELLINGER TRANSFORM FOR GAUSSIAN AR PROCESSES

Considerations in this section have been inspired by those of Pap, Bokor and Gáspár [10]. Let $z_{(k)}=\left(z_{0}^{\top}, \ldots, z_{k}^{\top}\right)^{\top}$, where $z_{0}, \ldots, z_{k} \in \mathbb{R}^{m}$ and ${ }^{\top}$ denotes the transpose of a vector. Let $\mathcal{N}_{m}(\mu, C)$ denote the $m$-dimensional normal distribution with mean vector $\mu$ and covariance matrix $C$.

Let us suppose that the measure $P_{i}$ is generated by an $m$-dimensional Gaussian $\mathrm{AR}(1)$ process. More precisely, we assume that the observed process belongs to the model

$$
\begin{equation*}
y_{i}(l+1)=A_{i} y_{i}(l)+w_{i}(l+1), \quad l=0,1,2, \ldots, \tag{29}
\end{equation*}
$$

if $\theta=i, i=1, \ldots, N$. Here $\left\{w_{i}(l), l=1,2, \ldots\right\}$ is a Gaussian white noise with covariance $Q_{i}$; i.e., $w_{i}(l) \sim \mathcal{N}_{m}\left(0, Q_{i}\right)$ and $w_{i}(1), w_{i}(2), \ldots$ are i.i.d. Suppose that the matrices $A_{i}$ and $Q_{i}$ are known and $Q_{i}$ is invertible for every $i=1, \ldots, N$. For the $i^{\text {th }}$ model, let the initial distribution be Gaussian with mean zero and with invertible covariance matrix $D_{i}$; i.e., $y_{i}(0) \sim \mathcal{N}_{m}\left(0, D_{i}\right)$ ( $i=1, \ldots, N$ ). Suppose that the process is observed until time $k . P_{i}$ is the distribution of $y_{i}(0), \ldots, y_{i}(k)$ (in the case of the $i^{\text {th }}$ model, i.e., when $\theta=i$ ).

Then the Hellinger transform is

$$
\begin{equation*}
\mathbb{H}_{\underline{\alpha}}\left(P_{1}, \ldots, P_{N}\right)=\mathbb{H}_{k}=\int \prod_{i=1}^{N} p\left(z_{(k)} \mid H_{i}\right)^{\alpha_{i}} d z_{(k)}, \tag{30}
\end{equation*}
$$

where

$$
\begin{equation*}
p\left(z_{(k)} \mid H_{i}\right)=(2 \pi)^{-((k+1) m) / 2}\left(\operatorname{det} C_{i}\right)^{-1 / 2} \exp \left\{-\frac{1}{2} z_{(k)}^{\top} C_{i}^{-1} z_{(k)}\right\}, \tag{31}
\end{equation*}
$$

for $z_{(k)} \in \mathbb{R}^{(k+1) m}$. (To emphasize that $\mathbb{H}_{\underline{\alpha}}\left(P_{1}, \ldots, P_{N}\right)$ depends on the number of observations $y_{i}(0), \ldots, y_{i}(k)$, we shall write $\mathbb{H}_{k}$ instead of $\mathbb{H}_{\underline{\alpha}}\left(P_{1}, \ldots, P_{N}\right)$. Obviously, $P_{i}$ and $C_{i}$ depend on $k$, too.)

Using the fact that the integral of a density function is equal to 1 (for the density of $\mathcal{N}_{(k+1) m}$ $\left(0,\left(\sum_{i=1}^{N} \alpha_{i} C_{i}^{-1}\right)^{-1}\right)$ ), we obtain

$$
\begin{align*}
\mathbb{H}_{k} & =\left(\operatorname{det}\left(\sum_{i=1}^{N} \alpha_{i} C_{i}^{-1}\right)^{-1}\right)^{1 / 2} \prod_{i=1}^{N}\left(\operatorname{det} C_{i}\right)^{-\alpha_{i} / 2}  \tag{32}\\
& =\left(\operatorname{det}\left(\sum_{i=1}^{N} \alpha_{i} C_{i}^{-1}\right)\right)^{-1 / 2} \prod_{i=1}^{N}\left(\operatorname{det} C_{i}^{-1}\right)^{\alpha_{i} / 2}
\end{align*}
$$

Now, $C_{i}$ is the following matrix

$$
\left(\begin{array}{ccccc}
B_{i}^{(0)} & B_{i}^{(0)} A_{i}^{\top} & B_{i}^{(0)} A_{i}^{\top} & \cdots & B_{i}^{(0)} A_{i}^{\top k}  \tag{33}\\
A_{i} B_{i}^{(0)} & B_{i}^{(1)} & B_{i}^{(1)} A_{i}^{\top} & \cdots & B_{i}^{(1)} A_{i}^{\top(k-1)} \\
A_{i}^{2} B_{i}^{(0)} & A_{i} B_{i}^{(1)} & B_{i}^{(2)} & \cdots & B_{i}^{(2)} A_{i}^{\top(k-2)} \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
A_{i}^{k} B_{i}^{(0)} & A_{i}^{k-1} B_{i}^{(1)} & A_{i}^{k-2} B_{i}^{(2)} & \cdots & B_{i}^{(k)}
\end{array}\right),
$$

where $B_{i}^{(j)}=A_{i}^{j} D_{i} A_{i}^{\top j}+\sum_{l=0}^{j-1} A_{i}^{l} Q_{i} A_{i}^{\top l}$. The inverse of $C_{i}$ is

$$
C_{i}^{-1}=\left(\begin{array}{ccccc}
D_{i}^{-1}+A_{i}^{\top} Q_{i}^{-1} A_{i} & -A_{i}^{\top} Q_{i}^{-1} & 0 & \cdots & 0  \tag{34}\\
-Q_{i}^{-1} A_{i} & Q_{i}^{-1}+A_{i}^{\top} Q_{i}^{-1} A_{i} & \ddots & \ddots & \vdots \\
0 & \ddots & \ddots & \ddots & 0 \\
\vdots & \ddots & \ddots & Q_{i}^{-1}+A_{i}^{\top} Q_{i}^{-1} A_{i} & -A_{i}^{\top} Q_{i}^{-1} \\
0 & \cdots & 0 & -Q_{i}^{-1} A_{i} & Q_{i}^{-1}
\end{array}\right)
$$

and $C_{i}^{-1}=F_{i} F_{i}^{\top}$, where

$$
F_{i}=\left(\begin{array}{ccccc}
D_{i}^{-1 / 2} & -A_{i}^{\top} Q_{i}^{-1 / 2} & 0 & \cdots & 0  \tag{35}\\
0 & Q_{i}^{-1 / 2} & \ddots & \ddots & \vdots \\
0 & \ddots & \ddots & \ddots & 0 \\
\vdots & \ddots & \ddots & \ddots & -A_{i}^{\top} Q_{i}^{-1 / 2} \\
0 & \cdots & 0 & 0 & Q_{i}^{-1 / 2}
\end{array}\right)
$$

We remark that the above matrices depend on $k$. We obtain

$$
\begin{equation*}
\mathbb{H}_{k}=\left(\operatorname{det}\left(\sum_{i=1}^{N} \alpha_{i} C_{i}^{-1}\right)\right)^{-1 / 2} \prod_{i=1}^{N}\left(\left(\operatorname{det} Q_{i}\right)^{k} \operatorname{det} D_{i}\right)^{-\alpha_{i} / 2} \tag{36}
\end{equation*}
$$

Here $\sum_{i=1}^{N} \alpha_{i} C_{i}^{-1}$ has the form

$$
\sum_{i=1}^{N} \alpha_{i} C_{i}^{-1}=\left(\begin{array}{ccccc}
D+d_{2} & -d_{1} & 0 & \cdots & 0  \tag{37}\\
-d_{1}^{\top} & d_{0}+d_{2} & \ddots & \ddots & \vdots \\
0 & \ddots & \ddots & \ddots & 0 \\
\vdots & \ddots & \ddots & d_{0}+d_{2} & -d_{1} \\
0 & \cdots & 0 & -d_{1}^{\top} & d_{0}
\end{array}\right)
$$

where

$$
\begin{equation*}
d_{0}=\sum \alpha_{i} Q_{i}^{-1}, \quad d_{1}=\sum \alpha_{i} Q_{i}^{-1} A_{i}, \quad d_{2}=\sum \alpha_{i} A_{i}^{\top} Q_{i}^{-1} A_{i}, \quad D=\sum \alpha_{i} D_{i}^{-1} \tag{38}
\end{equation*}
$$

Eliminating the elements over the diagonal blocks, we get a matrix with diagonal blocks $M(1), \ldots$, $M(k+1)$, where

$$
\begin{align*}
M(1) & =d_{0},  \tag{39}\\
M(l+1) & =d_{0}+d_{2}-d_{1}^{\top} M(l)^{-1} d_{1}, \quad l=1, \ldots, k-1,  \tag{40}\\
M(k+1) & =D+d_{2}-d_{1}^{\top} M(k)^{-1} d_{1} . \tag{41}
\end{align*}
$$

Therefore (36) implies that

$$
\begin{equation*}
\mathbb{H}_{k}=\prod_{i=1}^{N}\left(\operatorname{det} D_{i}\right)^{-\alpha_{i} / 2}\left[\prod_{i=1}^{N}\left(\operatorname{det} Q_{i}\right)^{-\alpha_{i} / 2}\right]^{k}\left(\prod_{l=1}^{k+1} \operatorname{det} M(l)\right)^{-1 / 2} . \tag{42}
\end{equation*}
$$

Here $M(1)=d_{0}$ is positive definite. Furthermore, one can prove that $M(2) \geq M(1)$ with equality iff $A_{1}=\cdots=A_{N}$. To prove this, we need the following lemmas. (Lemmas 3.1 and 3.2 were communicated to the authors by Losonczi and Páles. For the operator version of these lemmas, see [11].)
Lemma 3.1. Let $U$ be a symmetric positive definite matrix. Then

$$
\begin{equation*}
x^{\top} U^{-1} x+z^{\top} U z \geq x^{\top} z+z^{\top} x, \tag{43}
\end{equation*}
$$

for each pair of vectors $x, z$. Equality holds in (43) if and only if $x=U z$.
Proof. (43) can be obtained from the inequality $\left\|U^{-1 / 2} x-U^{1 / 2} z\right\| \geq 0$, where $\|\cdot\|$ denotes the Euclidean norm.

Lemma 3.2. Let $V_{1}, V_{2}, \ldots, V_{N}$ be symmetric positive definite matrices and let $V=V_{1}+\cdots+V_{N}$. Then

$$
\begin{equation*}
\sum_{i=1}^{N} x_{i}^{\top} V_{i}^{-1} x_{i} \geq\left(\sum_{i=1}^{N} x_{i}\right)^{\top} V^{-1}\left(\sum_{i=1}^{N} x_{i}\right) \tag{44}
\end{equation*}
$$

for arbitrary vectors $x_{1}, \ldots, x_{N}$. Equality holds in (43) if and only if $V_{1}^{-1} x_{1}=V_{2}^{-1} x_{2}=\cdots=$ $V_{N}^{-1} x_{N}$.
Proof. For $i=1, \ldots, N$, put $U=V_{i}, x=x_{i}$, and $z=V^{-1}\left(\sum_{j=1}^{N} x_{j}\right)$ in (43) and sum the inequalities obtained. We get

$$
\begin{equation*}
\sum_{i=1}^{N} x_{i}^{\top} V_{i}^{-1} x_{i}+z^{\top}\left(\sum_{j=1}^{N} x_{j}\right) \geq \sum_{i=1}^{N} x_{i}^{\top} V^{-1} \sum_{j=1}^{N} x_{j}+\sum_{j=1}^{N} z^{\top} x_{j} \tag{45}
\end{equation*}
$$

which is equivalent to (44). The condition for equality is $x_{i}=V_{i} z$; that is, $V^{-1} x_{i}$ is independent of $i$.

Now, we can return to the proof of $M(2) \geq M(1)$. We have

$$
M(2)-M(1)=\sum_{i=1}^{N} \alpha_{i} A_{i}^{\top} Q_{i}^{-1} A_{i}-\left(\sum_{i=1}^{N} \alpha_{i} Q_{i}^{-1} A_{i}\right)^{\top}\left(\sum_{i=1}^{N} \alpha_{i} Q_{i}^{-1}\right)^{-1}\left(\sum_{i=1}^{N} \alpha_{i} Q_{i}^{-1} A_{i}\right)
$$

Therefore, with notation $V_{i}=\alpha_{i} Q_{i}^{-1}$, we have

$$
M(2)-M(1)=\sum_{i=1}^{N}\left(A_{i}^{\top} V_{i}\right) V_{i}^{-1}\left(V_{i} A_{i}\right)-\left(\sum_{i=1}^{N} V_{i} A_{i}\right)^{\top}\left(\sum_{i=1}^{N} V_{i}\right)^{-1}\left(\sum_{i=1}^{N} V_{i} A_{i}\right) \geq 0
$$

In the last step, we used (44) with $x_{i}=V_{i} A_{i} a$, where $a$ is an arbitrary vector. Equality holds if $V_{i}^{-1} x_{i}=A_{i} a$ does not depend on $i$ for every $a$. As

$$
\begin{equation*}
M(l+1)-M(l)=d_{1}^{\top}\left(M(l-1)^{-1}-M(l)^{-1}\right) d_{1}, \quad l=2, \ldots, k-1 \tag{46}
\end{equation*}
$$

it follows that

$$
\begin{equation*}
0<d_{0}=M(1) \leq M(2) \leq \cdots \leq M(k) \leq d_{0}+d_{2} \tag{47}
\end{equation*}
$$

with equality iff $A_{1}=\cdots=A_{N}$.
Substituting $M(l), l=2, \ldots, k$, by $M(1)$ in equality (42), we get the estimation

$$
\begin{equation*}
\mathbb{H}_{k} \leq \text { const } \cdot c^{k / 2} \tag{48}
\end{equation*}
$$

where

$$
\begin{equation*}
c=\prod_{i=1}^{N} \frac{\left(\operatorname{det} Q_{i}^{-1}\right)^{\alpha_{i}}}{\operatorname{det}\left(\sum_{i=1}^{N} \alpha_{i} Q_{i}^{-1}\right)} . \tag{49}
\end{equation*}
$$

$c \leq 1$ with equality iff $Q_{1}=\cdots=Q_{N}$. This estimation cannot be improved if $A_{1}=\cdots=A_{N}$. Otherwise, using $M(l), 2 \leq l \leq k$, instead of $M(1)$, one can get a better estimation than (48),(49). In the special case when $Q_{1}=\cdots=Q_{N}=Q$, one can use, e.g., $M(2)$ (for which $M(2) \geq$ $M(1) \geq Q$ with equality iff $\left.A_{1}=\cdots=A_{N}\right)$ to obtain the estimation

$$
\begin{equation*}
\mathbb{H}_{k} \leq \text { const } \cdot\left[\frac{\operatorname{det} Q}{\operatorname{det} M(2)}\right]^{k / 2} \tag{50}
\end{equation*}
$$

We mention two possible choices of the initial distribution. For the first, if $D_{i}$ is the solution of the equation $D_{i}=A_{i} D_{i} A_{i}^{\top}+Q_{i}, i=1, \ldots, N$, then we have stationary $\operatorname{AR}(1)$-processes. For the second, if all the $D_{i}$ 's are equal to the dispersion of the stationary solution of model $\theta=1$, then one can study a possible change at time 0 from the first model for one of the other models.

## 4. THE HELLINGER TRANSFORM <br> OF CONTINUOUS TIME PROCESSES

First we consider one-dimensional continuous time AR(1) or Ornstein-Uhlenbeck processes. Let $y_{i}(t), 0 \leq t \leq T$, be the solution of the stochastic differential equation

$$
\begin{equation*}
d y_{i}(t)=-\lambda_{i} y_{i}(t) d t+d w(t), \quad 0 \leq t \leq T, \tag{51}
\end{equation*}
$$

where $w(t)$ is a standard Wiener process and $\lambda_{i}>0, i=0,1, \ldots, N$. Suppose that $y_{i}(t)$ is stationary $\left(y_{i}(0) \sim \mathcal{N}\left(0,1 / 2 \lambda_{i}\right)\right)$. Then $E y_{i}(t)=0$ and $\operatorname{cov}\left(y_{i}(t), y_{i}(t+s)\right)=\exp \left[-\lambda_{i} s\right] / 2 \lambda_{i}$.
It follows from formula (2.3.41) of [12] that the Radon-Nikodym derivative with respect to $P_{\lambda_{0}}$ has the form

$$
\begin{equation*}
\frac{d P_{\lambda_{i}}}{d P_{\lambda_{0}}}(y(t))=\sqrt{\frac{\lambda_{i}}{\lambda_{0}}} \exp \left\{\frac{\lambda_{0}^{2}-\lambda_{i}^{2}}{2} \int_{0}^{T} y^{2}(t) d t+\frac{\lambda_{i}-\lambda_{0}}{2} T-\frac{\lambda_{i}-\lambda_{0}}{2}\left[y^{2}(T)+y^{2}(0)\right]\right\}, \tag{52}
\end{equation*}
$$

where $P_{\lambda_{i}}$ denotes the measure generated by the process given by equation (51). Then the Hellinger transform of measures $P_{\lambda_{1}}, \ldots, P_{\lambda_{N}}$ is

$$
\begin{align*}
\mathbb{H}_{\underline{\alpha}}\left(P_{\lambda_{1}}, \ldots, P_{\lambda_{N}}\right)= & E_{0}\left\{\prod_{i=1}^{N}\left(\frac{d P_{\lambda_{i}}}{d P_{\lambda_{0}}}\right)^{\alpha_{i}}\right\} \\
= & E_{0}\left\{\sqrt { \frac { \prod _ { i = 1 } ^ { N } \lambda _ { i } ^ { \alpha _ { i } } } { \lambda _ { 0 } } } \operatorname { e x p } \left[\frac{\lambda_{0}^{2}-\sum_{i=1}^{N} \alpha_{i} \lambda_{i}^{2}}{2} \int_{0}^{T} y^{2}(t) d t\right.\right.  \tag{53}\\
& \left.+\frac{\sum_{i=1}^{N} \alpha_{i} \lambda_{i}-\lambda_{0}}{2}\left[T-\left(y^{2}(T)+y^{2}(0)\right)\right]\right\}
\end{align*}
$$

where $E_{0}$ denotes that the expectation is taken with respect to the measure defined by the process determined by equation (51) with parameter $\lambda_{0}$. Using Novikov's method (see, e.g., [12-14]), put $\lambda_{0}^{2}=\sum_{i=1}^{N} \alpha_{i} \lambda_{i}^{2}$. Then we obtain

$$
\begin{align*}
& \mathbb{H}_{\underline{\underline{q}}}\left(P_{\lambda_{1}}, \ldots, P_{\lambda_{N}}\right) \\
& =\sqrt{\frac{\prod_{i=1}^{N} \lambda_{i}^{\alpha_{i}}}{\lambda_{0}}} \exp \left[\frac{\sum_{i=1}^{N} \alpha_{i} \lambda_{i}-\lambda_{0}}{2} T\right] E_{0}\left\{\exp \left[-\frac{\sum_{i=1}^{N} \alpha_{i} \lambda_{i}-\lambda_{0}}{2}\left(y^{2}(T)+y^{2}(0)\right)\right]\right\} . \tag{54}
\end{align*}
$$

Now we need the following fact. If the common distribution of $\xi_{1}$ and $\xi_{2}$ is normal with zero mean and $D^{2} \xi_{1}=\sigma_{1}^{2}, D^{2} \xi_{2}=\sigma_{2}^{2}, \operatorname{cov}\left(\xi_{1}, \xi_{2}\right)=\varrho$, then (the Laplace transform)

$$
\begin{equation*}
E\left\{\exp \left[-\frac{t\left(\xi_{1}^{2}+\xi_{2}^{2}\right)}{2}\right]\right\}=\left[\left(t \sigma_{1}^{2}+1\right)\left(t \sigma_{2}^{2}+1\right)-t^{2} \varrho^{2}\right]^{-1 / 2} \tag{55}
\end{equation*}
$$

Therefore equation (54) implies the following proposition.
Proposition 4.1. The Hellinger transform of Ornstein-Uhlenbeck processes defined by (51) is

$$
\begin{align*}
\mathbb{H}_{\underline{\alpha}}\left(P_{\lambda_{1}}, \ldots, P_{\lambda_{N}}\right)= & 2 \sqrt{\lambda_{0} \prod_{i=1}^{N} \lambda_{i}^{\alpha_{i}}} \exp \left[\left(\sum_{i=1}^{N} \alpha_{i} \lambda_{i}-\lambda_{0}\right) \frac{T}{2}\right] \\
& \times\left[\left(\sum_{i=1}^{N} \alpha_{i} \lambda_{i}+\lambda_{0}\right)^{2}-\left(\sum_{i=1}^{N} \alpha_{i} \lambda_{i}-\lambda_{0}\right)^{2} \exp \left(-2 \lambda_{0} T\right)\right]^{-1 / 2} \tag{56}
\end{align*}
$$

where $\lambda_{0}^{2}=\sum_{i=1}^{N} \alpha_{i} \lambda_{i}^{2}$.

Now we consider one-dimensional shifted Wiener processes. Let $y_{i}(t), 0 \leq t \leq T$, be defined by the equation

$$
\begin{equation*}
y_{i}(t)=m_{i}(t)+w(t), \quad 0 \leq t \leq T \tag{57}
\end{equation*}
$$

$i=0,1, \ldots, N$, where $w(t)$ is a standard Wiener process. Suppose that $m_{i}(0)=0$ and $m_{i}(t)$ is absolute continuous with derivative $m_{i}^{\prime}(t)$ for which $\int_{0}^{T}\left(m_{i}^{\prime}(t)\right)^{2} d t<\infty$. Then it follows from [12, formula (2.3.17)] that the Radon-Nikodym derivative with respect to $P_{m_{0}}$ has the form

$$
\begin{align*}
\frac{d P_{m_{i}}}{d P_{m_{0}}}(y(t))= & \exp \left\{-\frac{1}{2} \int_{0}^{T}\left(m_{i}^{\prime}(t)\right)^{2} d t+\int_{0}^{T} m_{i}^{\prime}(t) d y(t)\right. \\
& \left.+\frac{1}{2} \int_{0}^{T}\left(m_{0}^{\prime}(t)\right)^{2} d t-\int_{0}^{T} m_{0}^{\prime}(t) d y(t)\right\}, \tag{58}
\end{align*}
$$

where $P_{m_{i}}$ denotes the measure generated by the process given by equation (57). Then the Hellinger transform of measures $P_{m_{1}}, \ldots, P_{m_{N}}$ is

$$
\begin{align*}
\mathbb{H}_{\underline{\alpha}}\left(P_{m_{1}}, \ldots, P_{m_{N}}\right)= & E_{0}\left\{\operatorname { e x p } \left[-\frac{1}{2} \int_{0}^{T} \sum_{i=1}^{N} \alpha_{i}\left(m_{i}^{\prime}(t)\right)^{2} d t\right.\right. \\
& \left.\left.+\int_{0}^{T} \sum_{i=1}^{N} \alpha_{i} m_{i}^{\prime}(t) d y(t)+\frac{1}{2} \int_{0}^{T}\left(m_{0}^{\prime}(t)\right)^{2} d t-\int_{0}^{T} m_{0}^{\prime}(t) d y(t)\right]\right\}, \tag{59}
\end{align*}
$$

where $E_{0}$ denotes that the expectation is taken with respect to the measure $P_{m_{0}}$. Using Novikov's method, put $m_{0}^{\prime}(t)=\sum_{i=1}^{N} \alpha_{i} m_{i}^{\prime}(t)$. Then we obtain

$$
\begin{equation*}
\mathbb{H}_{\underline{\alpha}}\left(P_{m_{1}}, \ldots, P_{m_{N}}\right)=\exp \left\{-\frac{1}{2} \int_{0}^{T} \sum_{i=1}^{N} \alpha_{i}\left(m_{i}^{\prime}(t)\right)^{2} d t+\frac{1}{2} \int_{0}^{T}\left(\sum_{i=1}^{N} \alpha_{i} m_{i}^{\prime}(t)\right)^{2} d t\right\} \tag{60}
\end{equation*}
$$

Therefore we have the following proposition.
Proposition 4.2. If $m_{i}(t), i=1, \ldots, N$, satisfy conditions listed above then the Hellinger transform of the processes defined by (57) is given by (60).

## 5. APPLICATIONS TO THE PROBABILITY OF MISCLASSIFICATION

Consider the classification problem introduced in Section 1 by (1),(2). A classification rule is a measurable mapping

$$
\begin{equation*}
d: \Omega \rightarrow\{1, \ldots, N\} . \tag{61}
\end{equation*}
$$

$\theta=i$ is accepted on the event $\{d=i\}$. The Bayesian risk connected with the decision rule $d$ is given by

$$
\begin{equation*}
R(d)=P\{\theta \neq d\}=\sum_{i=1}^{N} \pi_{i} P_{i}\{d \neq i\} . \tag{62}
\end{equation*}
$$

It is well known that a decision rule $d_{0}$ with

$$
\begin{equation*}
\left\{d_{0} \neq i\right\} \supseteq\left\{\pi_{i} X_{i}<\max \left(\pi_{1} X_{1}, \ldots, \pi_{N} X_{N}\right)\right\} \tag{63}
\end{equation*}
$$

for $i=1, \ldots, N$ minimizes the risk (i.e., the probability of error). Furthermore, the resulting probability of error is

$$
\begin{equation*}
P_{e}=R\left(d_{0}\right)=1-\int \max _{1 \leq i \leq n}\left(\pi_{i} X_{i}\right) d \mu . \tag{64}
\end{equation*}
$$

One can easily give an upper bound on $P_{e}$ in terms of Hellinger integrals:

$$
\begin{equation*}
P_{e} \leq \sum_{i<j} \pi_{i}^{\alpha_{1}^{(i j)}} \pi_{j}^{\alpha_{2}^{(i j)}} \mathbb{H}_{\underline{\underline{G}}^{(i j)}}\left(P_{i}, P_{j}\right) \tag{65}
\end{equation*}
$$

(here $\underline{\alpha}^{(i j)}=\left(\alpha_{1}^{(i j)}, \alpha_{2}^{(i j)}\right)$ is a vector with nonnegative coordinates and $\left.\alpha_{1}^{(i j)}+\alpha_{2}^{(i j)}=1\right)$. The following inequality can be proved by the same way as the theorem of Bhattacharya and Toussaint [2]:

$$
\begin{equation*}
P_{e} \leq \frac{N-2}{N-1}+\frac{1}{N-1}\left(\prod_{i=1}^{N} \pi_{i}^{\alpha_{i}}\right) \mathbb{H}_{\underline{\alpha}}\left(P_{1}, \ldots, P_{N}\right) . \tag{66}
\end{equation*}
$$

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