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Working Paper

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Working Paper, No. 581

Provided in Cooperation with:

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Suggested Citation: Hu, Yingyao; Shiu, Ji-liang (2011): Nonparametric identification using instrumental variables: Sufficient conditions for completeness, Working Paper, No. 581, The Johns Hopkins University, Department of Economics, Baltimore, MD

This Version is available at: https://hdl.handle.net/10419/49869

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Nonparametric Identification Using Instrumental Variables:

Sufficient Conditions For Completeness

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June 12, 2011

Abstract

This paper provides sufficient conditions for the nonparametric identification of the regression function $m(\cdot)$ in a regression model with an endogenous regressor x and an instrumental variable z. It has been shown that the identification of the regression function from the conditional expectation of the dependent variable on the instrument relies on the completeness of the distribution of the endogenous regressor conditional on the instrument, i.e., f(x|z). We provide sufficient conditions for the completeness of f(x|z) without imposing a specific functional form, such as the exponential family. We show that if the conditional density f(x|z) coincides with an existing complete density at a limit point in the support of z, then f(x|z) itself is complete, and therefore, the regression function $m(\cdot)$ is nonparametrically identified. We use this general result provide specific sufficient conditions for completeness in three different specifications of the relationship between the endogenous regressor x and the instrumental variable z.

Keywords: nonparametric identification, instrumental variable, completeness, endogeneity.

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1. Introduction

We consider a nonparametric regression model as follows:

$$y = m(x) + u, (1)$$

where the regressor x may be correlated with a zero mean regression error u. The parameter of interest is the nonparametric regression function $m(\cdot)$. In a random sample of $\{y, x, z\}$, an instrumental variable z is conditional mean independent of the regression error u, i.e., E(u|z) = 0, which implies

$$E[y|z] = \int_{-\infty}^{+\infty} m(x)f(x|z)dx. \tag{2}$$

This paper provides sufficient conditions on the conditional density f(x|z) under which the regression function $m(\cdot)$ is nonparametrically identified from, i.e., uniquely determined by, the conditional mean E[y|z]. We show that if the conditional density f(x|z) coincides with an existing complete density at a limit point in the support of z, then f(x|z) itself is complete, and consequently, the regression function $m(\cdot)$ is nonparametrically identified. Our sufficient conditions for completeness impose no specific functional form on f(x|z), such as the exponential family.

We assume the regression function $m(\cdot)$ is in a Hilbert space \mathcal{H} of functions defined on \mathcal{X} the support of regressor x. For example, we may define $\mathcal{L}^2(\mathcal{X}) = \{h(\cdot) : \int_{\mathcal{X}} |h(x)|^2 dx < \infty, \}$ be a L^2 space with the following inner product $\langle f, g \rangle \equiv \int_{\mathcal{X}} f(x)g(x)dx$. We may define the corresponding norm as: $||f||^2 = \langle f, f \rangle$. The completion of $\mathcal{L}^2(\mathcal{X})$ under the norm $||\cdot||$ is a Hilbert space. One may show that the uniqueness of the regression function $m(\cdot)$ is implied by the completeness of the family $\{f(\cdot|z) : z \in \mathcal{O}\}$ in \mathcal{H} , where $\mathcal{O} \subseteq \mathcal{Z}$ is a subset of \mathcal{Z} the support of z. The set \mathcal{O} may be \mathcal{Z} itself or some subset of \mathcal{Z} . In particular, we consider the completeness with the set \mathcal{O} being a sequence $\{z_k : k = 1, 2, 3, ...\}$ in \mathcal{Z} . The latter case corresponds to a sequence of functions $\{f(\cdot|z_k) : k = 1, 2, ...\}$. We start with the definition of the completeness.

Definition 1. The family $\{f(\cdot|z) \in \mathcal{H} : z \in \mathcal{O}\}\$ for some set $\mathcal{O} \subseteq \mathcal{Z}$ is said to be complete in

 \mathcal{H} if for any $h(\cdot) \in \mathcal{H}$

$$\int_{\mathcal{X}} h(x)f(x|z)dx = 0 \quad \text{for all } z \in \mathcal{O}$$

implies h(x) = 0 for any $x \in \mathcal{X}$. When it is a conditional density function with support $\mathcal{X} \times \mathcal{Z}$, f(x|z) is said to be a complete density.

The completeness introduced in Definition 1 is equivalent to L^2 - completeness considered in Andrews (2011) since a L^2 space is also a Hilbert space.¹

The uniqueness (identification) of the regression function $m(\cdot)$ is implied by the completeness of the family $\{f(\cdot|z):z\in\mathcal{O}\}$ in \mathcal{H} for some set $\mathcal{O}\subseteq\mathcal{Z}$. This sufficient condition may be shown as follows. Suppose that $m(\cdot)$ is not identified so that there are two different functions $m(\cdot)$ and $\widetilde{m}(\cdot)$ in \mathcal{H} which are observationally equivalent, i.e., for any $z\in\mathcal{Z}$

$$E[y|z] = \int_{\mathcal{X}} m(x)f(x|z)dx = \int_{\mathcal{X}} \widetilde{m}(x)f(x|z)dx.$$
 (3)

We then have for some $h(x) = m(x) - \widetilde{m}(x) \neq 0$

$$\int_{\mathcal{X}} h(x)f(x|z)dx = 0 \text{ for any } z \in \mathcal{Z}$$

which implies that $\{f(\cdot|z):z\in\mathcal{O}\}$ for any $\mathcal{O}\subseteq\mathcal{Z}$ is not complete in \mathcal{H} . Therefore, if $\{f(\cdot|z):z\in\mathcal{O}\}$ for some $\mathcal{O}\subseteq\mathcal{Z}$ is complete in \mathcal{H} , then $m(\cdot)$ is uniquely determined by E[y|z] and f(x|z), and therefore, is nonparametrically identified.

The following two examples of complete f(x|z) are from Newey and Powell (2003) (See their Theorem 2.2 and 2.3 for details.):

Example 1: Suppose that the distribution of x conditional on z is $N(a + bz, \sigma^2)$ for $\sigma^2 > 0$ and the support of z contains an open set, then $E[h(\cdot)|z] = 0$ for any $z \in \mathcal{Z}$ implies h(x) = 0 for any $x \in \mathcal{X}$; equivalently, $\{f(x|z) : z \in \mathcal{Z}\}$ is complete in $L^2(\mathcal{X})$.

Another case where the $\{f(x|z):z\in\mathcal{O}\}$ is complete in \mathcal{H} is that f(x|z) belongs to an exponential family as follows:

Example 2: Let $f(x|z) = s(x)t(z) \exp [\mu(z)\tau(x)]$, where s(x) > 0, $\tau(x)$ is one-to-one in x, and support of $\mu(z)$, \mathcal{Z} , is an open set, then $E[h(\cdot)|z] = 0$ for any $z \in \mathcal{Z}$ implies h(x) = 0

Closely related definitions of L^2 -completeness can also be found in Florens, Mouchart, and Rolin (1990), Mattner (1996), and San Martin and Mouchart (2007).

for any $x \in \mathcal{X}$; equivalently, the family of conditional density functions $\{f(x|z) : z \in \mathcal{Z}\}$ is complete in $L^2(\mathcal{X})$.

These two examples show the completeness of a family $\{f(x|z): z \in \mathcal{O}\}$, where \mathcal{O} is an open set. In order to extend the completeness to general density functions, we further reduce the set \mathcal{O} from an open set to a countable set with a limit point, i.e. a converging sequence in the support \mathcal{Z} .

This paper focuses on the sufficient conditions for completeness of a conditional density. These conditions can be used to obtain global or local identification in a variety of models including the nonparametric IV regression model (see Newey and Powell (2003); Darolles, Florens, and Renault (2002); Hall and Horowitz (2005)), semiparametric IV models (see Ai and Chen (2003); Blundell, Chen, and Kristensen (2007)), nonparametric IV quantile models (see Chernozhukov and Hansen (2005); Chernozhukov, Imbens, and Newey (2007); Horowitz and Lee (2007)), measurement error models (see Hu and Schennach (2008); An and Hu (2009); Carroll, Chen, and Hu (2010); Chen and Hu (2006)), and dynamic models (see Hu and Shum (2009); Shiu and Hu (2010)), etc. We refer to D'Haultfoeuille (2011) and Andrews (2011) for more complete literature reviews.

In this paper, we provide sufficient conditions for the completeness of a general conditional density without imposing particular functional forms. We show that if the conditional density f(x|z) is close to a complete density, then f(x|z) itself is complete. We use the results in the stability of bases in Hilbert space (section 10 of chapter 1 in Young (1980)) to show that a linearly independent sequence is complete if its deviation from a complete sequence of function is finite. We then show that two sequences of density functions have a finite deviation when they have the same limit. Based on this observation, we may deviate from the existing complete density function without losing the completeness.

We apply the general results to show the completeness in three scenarios. First, we extend Example 2 to a general setting. In particular, we show the completeness of f(x|z) when x and z satisfy for some function $\mu(\cdot)$ and $\sigma(\cdot)$

$$x = \mu(z) + \sigma(z) \varepsilon$$
 with $z \perp \varepsilon$.

Second, we consider a general control function

$$x = h(z, \varepsilon)$$
 with $z \perp \varepsilon$,

and provide conditions for completeness of f(x|z) in this case. Third, our results implies that the completeness of a multidimensional conditional density, e.g.,

$$f(x_1, x_2|z_1, z_2),$$

may be implied by the completeness of two conditional densities of lower dimension, i.e., $f(x_1|z_1)$ and $f(x_2|z_2)$.

This paper is organized as follows: section 2 provides sufficient conditions for completeness; section 3 applies the main results to the three cases with different specifications of the relationship between the endogenous variable and the instrument; section 4 concludes the paper and all the proofs are in the appendix.

2. Sufficient Conditions for Completeness

In this section, we show that a sequence $\{f(\cdot|z_k)\}$ is complete if it coincides with a complete sequence $\{g(\cdot|z_k)\}$ at a limit point z_0 . We start with the introduction of two well-known complete families in Examples 1 and 2. Notice that these completeness results are established on an open set \mathcal{O} instead of a countable set with a limit point, i.e., a converging sequence. In order to extend the completeness to a new function f(x|z), we first establish the completeness on a sequence of z_k .

As we will show below, the completeness of an existing sequence $\{g(x|z_k): k=1,2,...\}$ is essential to show the completeness for a new function f(x|z). An important family of conditional distributions which admit completeness is the exponential family. Many distributions encountered in practice can be put into the form of exponential families, including Gaussian, Poisson, Binomial, and certain multivariate form of these. Another family of conditional distribution which implies completeness is in the form of a convolution density function, i.e., g(x|z) = g(x-z).

Based on the existing results, such as in Examples 1 and 2 in the introduction, we may

generate complete sequences from the exponential family or a convolution density function. We start with an introduction of a complete sequence in the exponential family. Example 2 shows the completeness of the family $\{g(\cdot|z):z\in\mathcal{O}\}$, where \mathcal{O} is an open set in \mathcal{Z} . In the next lemma, we reduce the set \mathcal{O} from an open set to a countable set with a limit point, i.e. a converging sequence in \mathcal{Z} .

Lemma 1. Let $\{z_k : k = 1, 2, ...\}$ be a sequence of distinct $z_k \in \mathcal{O}$ converging to z_0 in an open set $\mathcal{O} \subseteq \mathcal{Z}$. Define

$$g(x|z) = s(x)t(z) \exp\left[\mu(z)\tau(x)\right]$$

on $\mathcal{X} \times \mathcal{Z}$ with $s(\cdot) > 0$ and $t(\cdot) > 0$. Suppose that $g(\cdot|z) \in \mathcal{L}^2(\mathcal{X})$ for $z \in \mathcal{O}$ and

- i) $\mu(\cdot)$ is continuous with $\mu'(z_0) \neq 0$;
- ii) $\tau(\cdot)$ is monotonic over \mathcal{X} .

Then, the sequence $\{g(\cdot|z_k): k=1,2,...\}$ is complete in $\mathcal{L}^2(\mathcal{X})$.

Proof: See the appendix.

Another case where the completeness of g(x|z) is well studied is when $g(x|z) = f_{\varepsilon}(x-z)$, which is usually due to a convolution between the endogenous variable x and instrument z as follows

$$x = z + \varepsilon$$
 with $z \perp \varepsilon$.

Example 1 suggests that the completeness of the family $\{g(\cdot|z) \in \mathcal{H} : z \in \mathcal{O}\}$, where \mathcal{O} is an open set in \mathcal{Z} and ε is normal. Again, we show the completeness still holds when the set \mathcal{O} is a converging sequence. We summarize the results as follows.

Lemma 2. Let $\{z_k : k = 1, 2, ...\}$ be a sequence of distinct $z_k \in \mathcal{O}$ converging to z_0 in an open set $\mathcal{O} \subseteq \mathcal{Z}$. Define

$$g(x|z) = f_{\varepsilon}(x-z)$$

on $\mathbb{R} \times \mathcal{Z}$. Suppose that $g(\cdot|z) \in \mathcal{L}^2(\mathbb{R})$ for $z \in \mathcal{O}$ and the Fourier transform ϕ_{ε} of f_{ε} satisfies

$$0 < |\phi_{\varepsilon}(t)| < Ce^{-\delta|t|} \tag{4}$$

for all $t \in \mathbb{R}$ and some constants C, $\delta > 0$.

Then, the sequence $\{g(\cdot|z_k): k=1,2,...\}$ is complete in $\mathcal{L}^2(\mathbb{R})$.

Proof: See the appendix.

Equation (4) is equivalent to the following two conditions: i) $\phi_{\varepsilon}(t) \neq 0$ for all $t \in \mathbb{R}$; ii) $|\phi_{\varepsilon}(t)| < Ce^{-\delta|t|}$ as $|t| \to \infty$ for some constant C and $\delta > 0$. In other words, Equation (4) implies that the ch.f. does not varnish on the real line and that the ch.f. has exponentially decaying tails. Notice that $|\phi_{\varepsilon}(t)| < Ce^{-\delta|t|}$ is not binding for a finite t.

With the complete sequences explicitly specified in Lemma 1 and 2, we are ready to extend the completeness to a more general conditional density f(x|z). Our sufficient conditions for completeness are summarized as follows:

Theorem 1. For every $z \in \mathcal{Z}$, let $f(\cdot|z)$ and $g(\cdot|z)$ be in a Hilbert space \mathcal{H} of functions defined on \mathcal{X} with norm $\|\cdot\|$. Suppose that there exists a point z_0 with its open neighborhood $\mathcal{N}(z_0) \subseteq \mathcal{Z}$ such that

- i) for every sequence $\{z_k : k = 1, 2, ...\}$ of distinct $z_k \in \mathcal{N}(z_0)$ converging to z_0 , the corresponding sequence $\{g(\cdot|z_k) : k = 1, 2, ...\}$ is complete in a Hilbert space \mathcal{H} ;
- ii) the relative deviation $D(z) = \frac{\|f(\cdot|z) g(\cdot|z)\|}{\|g(\cdot|z)\|}$ is continuous in z on $\mathcal{N}(z_0)$ and $f(\cdot|z)$ coincides with $g(\cdot|z)$ at z_0 in \mathcal{H} , i.e.,

$$||f(\cdot|z_0) - g(\cdot|z_0)|| = 0;$$

iii) there exists a sequence $\{z_k : k = 1, 2, ...\}$ of distinct $z_k \in \mathcal{N}(z_0)$ converging to z_0 such that the sequence $\{f(\cdot|z_k) : k = 1, 2, ...\}$ is linearly independent, i.e.,

$$\sum_{i=1}^{I} c_i f(x|z_{k_i}) = 0 \text{ for all } x \in \mathcal{X} \text{ implies } c_i = 0 \text{ for all } i = 1, 2, ..., I.$$

Then, the family $\{f(\cdot|z): z \in \mathcal{N}(z_0)\}$ is complete in \mathcal{H} .

Proof: See the appendix.

Condition i) provides complete sequences, which may be from Lemma 1 and 2. Condition ii) requires that the relative deviation $D(z) = \frac{\|f(\cdot|z) - g(\cdot|z)\|}{\|g(\cdot|z)\|}$ is continuous in z on $\mathcal{N}(z_0)$. When the Hilbert space \mathcal{H} is the $\mathcal{L}^2(\mathcal{X})$, the relative deviation D(z) is continuous if $\|g(\cdot|z_0)\| > 0$ and the first-order derivatives $\frac{\partial}{\partial z} f(\cdot|z)$ and $\frac{\partial}{\partial z} g(\cdot|z)$ are also in $\mathcal{L}^2(\mathcal{X})$ for $z \in \mathcal{N}(z_0)$. This is because for some function $h(\cdot|z) \in \mathcal{L}^2(\mathcal{X})$ with $\frac{\partial}{\partial z} h(\cdot|z) \in \mathcal{L}^2(\mathcal{X})$ the derivative of $\|h(\cdot|z)\|^2$

w.r.t. z is finite due to the Cauchy-Schwarz inequality as follows:

$$\left| \frac{\partial}{\partial z} \left(\|h(\cdot|z)\|^2 \right) \right| \le 2 \|h(\cdot|z)\| \left\| \frac{\partial}{\partial z} h(\cdot|z) \right\|.$$

Furthermore, if \mathcal{X} is bounded, we only need $\frac{\partial}{\partial z} h(\cdot|z)$ to be bounded. In the proof of Theorem 1, we show that the continuity of D(z) implies that there exists a sequence $\{z_k\}$ converging to z_0 such that the total deviation from the sequence $\{g(\cdot|z_k)\}$ to $\{f(\cdot|z_k)\}$ is finite, i.e.,

$$\sum_{k=1}^{\infty} \frac{\|f(\cdot|z_k) - g(\cdot|z_k)\|}{\|g(\cdot|z_k)\|} < \infty.$$
 (5)

Intuitively, this condition implies that the sequence $\{f(\cdot|z_k)\}$ is close to a complete sequence $\{g(\cdot|z_k)\}$ so that the former sequence may also be complete.

The linear independence in condition iii) imposed on $\{f(\cdot|z_k)\}$ implies that there are no redundant terms in the sequence in the sense that no term can be expressed as a linear combination of some other terms. This condition imposes a mild restriction on f(x|z) because Equation $\sum_{i=1}^{I} c_i f(x|z_{k_i}) = 0$ for all $x \in \mathcal{X}$, which implies an infinite number of restrictions on a finite number of constants c_i . When the support of $f(\cdot|z_k)$ is the whole real line for all z_k , a sufficient condition for the linear independence is that

$$\lim_{x \to -\infty} \frac{f(x|z_{k+1})}{f(x|z_k)} = 0,\tag{6}$$

which implies $\lim_{x\to-\infty}\frac{f(x|z_{k+m})}{f(x|z_k)}=0$ for any $m\geq 1$. If $\sum_{i=1}^I c_i f(\cdot|z_{k_i})=0$ for all $x\in(-\infty,+\infty)$, we may have

$$-c_1 = \sum_{i=1}^{I} c_i \frac{f(x|z_{k_i})}{f(x|z_{k_1})}.$$

The limit of the RHS is zero as $x \to -\infty$ so that $c_1 = 0$. Similarly, we may show $c_2, c_3, ..., c_I = 0$ for all i sequentially. Notice that the exponential family satisfies Equation (6). When the support \mathcal{X} is bounded, for example, $\mathcal{X} = [0, 1]$, the condition (6) may become

$$\lim_{x \to 0} \frac{f(x|z_{k+1})}{f(x|z_k)} = 0. \tag{7}$$

For example, the Corollary (Müntz) on page 91 in Young (1980) implies that the family of

function $\{x^{z_1}, x^{z_2}, x^{z_3}, ...\}$ is complete in $\mathcal{L}^2([0,1])$ if $\sum_{k=1}^{\infty} \frac{1}{z_k} = \infty$. This family also satisfies the condition (7) for a strictly increasing $\{z_k\}$. For an existing function g(x|z) > 0, we may always have $f(x|z) = \frac{f(x|z)}{g(x|z)} \times g(x|z)$. If the existing sequence $\{g(\cdot|z_k)\}$ satisfies Equation (6), i.e., $\lim_{x\to -\infty} \frac{g(x|z_{k+1})}{g(x|z_k)} = 0$, then it is enough to have $0 < \left(\lim_{x\to -\infty} \frac{f(x|z_k)}{g(x|z_k)}\right) < \infty$.

Furthermore, when $f(x|z) = h(x|z) \times g(x|z)$, the condition (6) is implied by $\lim_{x \to -\infty} \frac{g(x|z_{k+1})}{g(x|z_k)} = 0$ and $\left(\lim_{x \to -\infty} \frac{h(x|z_{k+1})}{h(x|z_k)}\right) < \infty$. We may also consider

$$f(x|z) = \lambda(z) h(x|z) + [1 - \lambda(z)] g(x|z). \tag{8}$$

In this case, the conditional density $f(\cdot|z)$ is a mixture of two continuous conditional densities h, g and the weight λ in the mixture depends on z. At the limit point z_0 , $f(\cdot|z_0)$ coincides with $g(\cdot|z_0)$ if $\lim_{z_k\to z_0}\lambda(z)=0$. The linear independence condition in Equation (6) holds when h(x|z) and g(x|z) satisfy

$$\lim_{x \to -\infty} \frac{g(x|z_{k+1})}{g(x|z_k)} = 0 \text{ and } \lim_{x \to -\infty} \frac{h(x|z_k)}{g(x|z_k)} < \infty.$$
 (9)

The advantage of this condition is that there are only mild restrictions imposed on the functional form of h(x|z) and $\lambda(z)$.

Suppose the function f(x|z) is differentiable up to any finite order. We may consider the so-called Wronskian determinant as follows:

$$W(x) = \det \begin{pmatrix} f(x|z_{k_1}) & f(x|z_{k_2}) & \dots & f(x|z_{k_I}) \\ f'(x|z_{k_1}) & f'(x|z_{k_2}) & \dots & f'(x|z_{k_I}) \\ & \dots & & \dots & \dots \\ \frac{d^{(I-1)}}{dx^{(I-1)}} f(x|z_{k_1}) & \frac{d^{(I-1)}}{dx^{(I-1)}} f(x|z_{k_2}) & \dots & \frac{d^{(I-1)}}{dx^{(I-1)}} f(x|z_{k_I}) \end{pmatrix}$$

If there exists an x_0 such that the determinant $W(x_0) \neq 0$ for every $\{z_{k_i} : i = 1, 2, ..., I\}$, then $\{f(\cdot|z_k)\}$ is linear independent. For example, we may have for z > 0 and $0 \in \mathcal{X}$

$$f(x|z) = \frac{d}{dx}F_0(z \times x)$$

with

$$W(0) = \Pi_{i=1}^{I} \left(z_{k_i} \frac{d^{(i)} F_0(0)}{dx^{(i)}} \right) \times \det \begin{pmatrix} 1 & 1 & \dots & 1 \\ z_{k_1} & z_{k_2} & \dots & z_{k_I} \\ \dots & \dots & \dots & \dots \\ (z_{k_1})^{I-1} & (z_{k_2})^{I-1} & \dots & (z_{k_I})^{I-1} \end{pmatrix}.$$

According to the property of the Vandermonde matrix, the determinant W(x) is not equal to zero when $f_0(x)$ has all the nonzero derivative at x = 0 and z_k are nonzero and distinctive. In general, we may also show $\{f(\cdot|z_k)\}$ is linear independent with

$$f(x|z) = \frac{d}{dx} F_0(\mu(z)\tau(x))$$

where $\mu'(z_0) \neq 0$ and $\tau(\cdot)$ is monotonic with $\tau(0) = 0$. This is because $\sum_{i=1}^{I} c_i f(\cdot | z_{k_i}) = 0$ is implied by $\sum_{i=1}^{I} c_i F_0(\mu(z_{k_i})\tau(\cdot)) = 0$, which holds if and only if $\sum_{i=1}^{I} c_i F_0(\mu(z_{k_i})\tau) = 0$ for all $\tau \in \tau(\mathcal{X})$. We may then show the determinant W(x) of the function F_0 is nonzero at x = 0.

Another sufficient condition for the linear independence is that the so-called Gram determinant G_f is not equal to zero for every $\{z_{k_i}: i=1,2,...,I\}$, where $G_f=\det\left(\left[\left\langle f(\cdot|z_{k_i}),f(\cdot|z_{k_j})\right\rangle\right]_{i,j}\right)$. This condition does not require the function has all the derivatives.

We summarize these results on the linear independence as follows:

Lemma 3. the sequence $\{f(\cdot|z_k)\}$ corresponding to a sequence $\{z_k : k = 1, 2, ...\}$ of distinct $z_k \in \mathcal{N}(z_0)$ converging to z_0 is linearly independent if one of the following conditions hold:

- 1) $\sum_{i=1}^{I} c_i f(x|z_{k_i}) = 0$ for all $x \in \mathcal{X}$ implies $c_i = 0$ for all i = 1, 2, ..., I.
- 2) for all k, $\lim_{x\to -\infty} \frac{f(x|z_{k+1})}{f(x|z_k)} = 0$ or $\lim_{x\to x_0} \frac{f(x|z_{k+1})}{f(x|z_k)} = 0$ for some x_0 ;
- 3) $f(x|z) = \frac{d}{dx}F_0(\mu(z)\tau(x))$ with $\mu'(z_0) \neq 0$, $\tau(0) = 0$, and $\frac{d^k}{dx^k}F_0(0) \neq 0$ for k = 1, 2, ...;
- 4) for every $\{z_{k_i}: i=1,2,...,I\}$, $\det\left(\left[\left\langle f(\cdot|z_{k_i}),f(\cdot|z_{k_j})\right\rangle\right]_{i,j}\right)\neq 0$.

In order to illustrate the relationship between the complete sequence $\{g(\cdot|z_k)\}$ and the sequence $\{f(\cdot|z_k)\}$, we present numerical examples of these two functions as follows. Consider $g(x|z) = x^z$ over $\mathcal{L}^2([0,0.8])$ for $\mathcal{Z} = (\frac{2}{5},\frac{3}{5})$. We pick $z_k = \frac{1}{2} - \frac{2}{(k+1)^2}$ with $z_k \to z_0 = \frac{1}{2}$. Since $\sum_{k=1}^{\infty} \frac{1}{z_k} = \infty$, by the Corollary (Müntz) on page 91 in Young (1980) the family of function

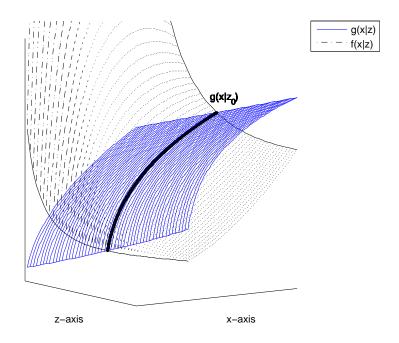


Figure 1: An example of g(x|z) and f(x|z) in Theorem 1.

 $\{x^{z_1}, x^{z_2}, x^{z_3}, ...\}$ is complete in $\mathcal{L}^2([0, 0.8])$. Let $f(x|z) = \left(1 - \frac{2(z - \frac{1}{2})}{(z - 0.62)}(x - 1)\right)x^z$. Since $\lim_{x \to 0} \frac{f(x|z_{k+1})}{f(x|z_k)} = 0$, our Theorem 1 implies that $\{f(\cdot|z_k)\}$ is also complete in $\mathcal{L}^2([0, 0.8])$ with $g(x|z_0) = f(x|z_0) = \sqrt{x}$. Figure 1 presents a 3D graph of g(x|z) and f(x|z) for (x, z) in $[0, 0.8] \times (\frac{2}{5}, \frac{3}{5})$ to illustrate the relationship between the complete sequence $\{g(\cdot|z_k)\}$ and the sequence $\{f(\cdot|z_k)\}$.

3. Applications

We consider three applications of our main results: first, we show the sufficient conditions for the completeness of f(x|z) when $x = \mu(z) + \sigma(z)\varepsilon$ with $z \perp \varepsilon$; second, we consider the completeness with a general control function $x = h(z, \varepsilon)$; finally, we show how to use our results to transform a multivariate completeness problem to a single variable one.

3.1. Extension of the convolution case

Lemma 2 provides a complete sequence when $x = z + \varepsilon$. Using Theorem 1, we may provide sufficient conditions for the completeness of f(x|z) when the endogenous variable x and the

instrument z satisfy a general heterogeneous structure as follows:

$$x = \mu(z) + \sigma(z) \varepsilon$$
 with $z \perp \varepsilon$.

We summarize the result as follows:

Lemma 4. For every $z \in \mathcal{Z}$, let $f(\cdot|z)$ be in $\mathcal{L}^2(\mathbb{R})$. Suppose that there exists a point z_0 with its open neighborhood $\mathcal{N}(z_0) \subseteq \mathcal{Z}$ such that

- i) the characteristic function $\phi_{z_0}(t)$ of $f(\cdot|z_0)$ satisfies $0 < |\phi_{z_0}(t)| < Ce^{-\delta|t|}$ for all $t \in \mathbb{R}$ and some constants C, $\delta > 0$;
 - ii) $\frac{\partial}{\partial z} f(\cdot|z)$ for $z \in \mathcal{N}(z_0)$ and $\frac{\partial}{\partial x} f(\cdot|z_0)$ are in $\mathcal{L}^2(\mathbb{R})$;
- iii) the function $f(\cdot|z)$ satisfies conditions iii) in Theorem 1, i.e., there exists a sequence $\{z_k : k = 1, 2, ...\}$ of distinct $z_k \in \mathcal{N}(z_0)$ converging to z_0 such that the sequence $\{f(\cdot|z_k) : k = 1, 2, ...\}$ is linearly independent.

Then, the family $\{f(\cdot|z):z\in\mathcal{N}(z_0)\}\$ is complete in $\mathcal{L}^2(\mathbb{R})$.

In particular, when

$$f(x|z) = \frac{1}{\sigma(z)} f_{\varepsilon} \left(\frac{x - \mu(z)}{\sigma(z)} \right)$$

on $\mathbb{R} \times \mathcal{Z}$, we assume

- i') the characteristic function $\phi_{\varepsilon}(t)$ of f_{ε} satisfies $0 < |\phi_{\varepsilon}(t)| < Ce^{-\delta|t|}$;
- ii') $\mu(\cdot)$, $\sigma(\cdot)$, and $f_{\varepsilon}(\cdot)$ are continuously differentiable with $\mu'(z_0) \neq 0$, $\sigma(z_0) \neq 0$ and $\int_{-\infty}^{+\infty} |xf'_{\varepsilon}(x)|^2 dx < \infty$;
 - iii') $\lim_{x\to-\infty} \frac{f_{\varepsilon}(x-c)}{f_{\varepsilon}(x)} = 0$ for any constant c > 0.

Then, the family $\{f(\cdot|z): z \in \mathcal{N}(z_0)\}$ is complete in $\mathcal{L}^2(\mathbb{R})$.

Proof: See the appendix.

The first part of Lemma 4 implies that one may always make a sequence coincide with a convolution sequence. Consider a sequence $\{f(\cdot|z_k): k=1,2,...\}$ with a sequence $\{z_k: k=1,2,...\}$ of distinct $z_k \in \mathcal{N}(z_0)$ converging to z_0 . We may always generate a convolution sequence $\{g(\cdot|z_k): k=1,2,...\}$ where

$$g(x|z_k) = f(x - \mu(z_k)|z_0)$$
 with $\mu(z_0) = 0$ and $\mu'(z_0) \neq 0$.

Condition i) implies that the sequence satisfies the conditions in 2 and is complete. Condition ii) guarantees that the first-order derivatives $\frac{\partial}{\partial z} f(\cdot|z)$ and $\frac{\partial}{\partial z} g(\cdot|z)$ are in $\mathcal{L}^2(\mathcal{X})$ so that the relative deviation $D(z) = \frac{\|f(\cdot|z) - g(\cdot|z)\|}{\|g(\cdot|z)\|}$ is continuous in z. Since $D(z_0) = 0$ by construction, the condition ii) in Theorem 1 holds. Thus, the completeness holds for $\{f(\cdot|z) : z \in \mathcal{N}(z_0)\}$. In the heterogeneous case where $x = \mu(z) + \sigma(z)\varepsilon$, a primitive condition for the linear independence is that $\lim_{x \to -\infty} \frac{f_{\varepsilon}(x-c)}{f_{\varepsilon}(x)} = 0$.

3.2. Completeness with a control function

We then consider a general expression of the relationship between the endogenous variable x and the instrument z. Let a control function describe the relationship between an endogenous variable x and an instrument z as follow:

$$x = h(z, \varepsilon), \text{ with } z \perp \varepsilon.$$
 (10)

We consider the case where x and ε have the support \mathbb{R} . Without loss of generality, we assume ε has a standard normal distribution with the cdf Φ . It is well known that the function h is related to the cdf $F_{x|z}$ as $h(z,\varepsilon) \equiv F_{x|z}^{-1}(\Phi(\varepsilon)|z)$ when the inverse of $F_{x|z}$ exists. Given the function h, we are interested in what restrictions on h are sufficient for the completeness of the conditional density f(x|z) implied by Equation (10).

Lemma 5. Let $\mathcal{N}(z_0) \subseteq \mathcal{Z}$ be an open neighborhood of some $z_0 \in \mathcal{Z}$. Let Equation (10) hold with $h(z_0, \varepsilon) = \varepsilon$, where ε has the support \mathbb{R} . Suppose that

- i) for $z \in \mathcal{N}(z_0)$, the function $h(z, \varepsilon)$ is strictly increasing in ε and twice differentiable in z and ε :
- ii) The density $f_{\varepsilon}(\cdot) = f_{x|z_0}(\cdot)$ and its characteristic function $\phi_{\varepsilon}(t)$ satisfy $\lim_{x \to -\infty} \frac{f_{\varepsilon}(x-c)}{f_{\varepsilon}(x)} = 0$ for any constant c > 0 and $0 < |\phi_{\varepsilon}(t)| < Ce^{-\delta|t|}$ for all $t \in \mathbb{R}$ and some constants C, $\delta > 0$
 - iii) $\frac{\partial}{\partial z} f(\cdot|z)$ for $z \in \mathcal{N}(z_0)$ and $f'_{x|z_0}(\cdot)$ are in $\mathcal{L}^2(\mathbb{R})$;
 - iv) for any $\widetilde{z} \neq \widehat{z}$ in $\mathcal{N}(z_0)$, $\lim_{\varepsilon \to -\infty} [h(\widetilde{z}, \varepsilon) h(\widehat{z}, \varepsilon)] \neq 0$.

Then, the family $\{f(\cdot|z): z \in \mathcal{N}(z_0)\}\$ is complete in $\mathcal{L}^2(\mathbb{R})$.

Proof: See the appendix.

Condition i) guarantees that the conditional density f(x|z) is continuous in both x and z around z_0 . The condition $h(z_0,\varepsilon)=\varepsilon$ is not restrictive because one may always redefine ε . Therefore, f(x|z) satisfies $f(x|z_0)=f_{\varepsilon}(x)$, which may be considered as a limit point in the convolution family such as $\{g(x|z)=f_{\varepsilon}(x-\mu(z_k)): k=1,2,...\}$ with $\mu(z_0)=0$, i.e. $f(x|z_0)=g(x|z_0)$. We may then use Theorem 1 to show f(x|z) is complete. Condition iii) implies that the deviation $D(z)=\frac{\|f(\cdot|z)-g(\cdot|z)\|}{\|g(\cdot|z)\|}$ is continuous in z. Condition iv) guarantees the linear independence of the sequence $\{f(\cdot|z_k)\}$.

Lemma 5 implies that a key sufficient assumption for the completeness of f(x|z) implied by the control function in Equation (10) is that the control function h is invertible with respect to ε around a limit point in the support of z. Our results may provide sufficient conditions for completeness with a general h. For example, we may have

$$h(z,\varepsilon) = \mu(z) + e^{z-z_0}\varepsilon + \sum_{j=0}^{J} (z-z_0)^{2j} h_j(\varepsilon),$$

where $h_{j}\left(\cdot\right)$ are increasing functions. The function h may also have a nonseparable form such as

$$h(z, \varepsilon) = \mu(z) + \ln\left[\left(z - z_0\right)^2 + \exp(\varepsilon)\right].$$

3.3. multivariate completeness

When the endogenous variable x and the instrument z are both vectors, our main results in Theorem 1 still applies. In other words, our results can be extended to the multivariate case straightforwardly. In this section, we show that one can use Theorem 1 to reduce a multivariate completeness problem to a single variate one. Without loss of generality, we consider $x = (x_1, x_2)$, $z = (z_1, z_2)$, $\mathcal{X} = \mathcal{X}_1 \times \mathcal{X}_2$, and $\mathcal{Z} = \mathcal{Z}_1 \times \mathcal{Z}_2$. One may show that the completeness of $f(x_1|z_1)$ and $f(x_2|z_2)$ implies that of $f(x_1|z_1) \times f(x_2|z_2)$. Theorem 1 then implies that if conditional density $f(x_1, x_2|z_1, z_2)$ coincides with $f(x_1|z_1) \times f(x_2|z_2)$ at a limit point in \mathcal{Z} then $f(x_1, x_2|z_1, z_2)$ is complete. We summarize the results as follows:

Lemma 6. For every $z \in \mathcal{Z} = \mathcal{Z}_1 \times \mathcal{Z}_2$, let $f(\cdot|z)$ and $g(\cdot|z)$ be in a Hilbert space \mathcal{H} of functions defined on $\mathcal{X} = \mathcal{X}_1 \times \mathcal{X}_2$ with norm $\|\cdot\|$. Suppose that there exists a point z_0 with its open neighborhood $\mathcal{N}(z_0) \subseteq \mathcal{Z}$ such that

- i) for every sequence $\{z_k : k = 1, 2, 3, ...\}$ of distinct $z_k \in \mathcal{N}(z_0)$ converging to z_0 , the corresponding sequence $\{f_{x_1|z_1}(\cdot|z_{1k}) : k = 1, 2, 3, ...\}$ and $\{f_{x_2|z_2}(\cdot|z_{2k}) : k = 1, 2, 3, ...\}$ are complete in Hilbert spaces \mathcal{H} of functions defined on \mathcal{X}_1 and \mathcal{X}_2 ;
- ii) the relative deviation $D(z_1, z_2) = \frac{\|f_{x|z}(\cdot, \cdot|z_1, z_2) f_{x_1|z_1}(\cdot|z_1)f_{x_2|z_2}(\cdot|z_2)\|}{\|f_{x_1|z_1}(\cdot|z_1)f_{x_2|z_2}(\cdot|z_2)\|}$ is continuous in $z = (z_1, z_2)$ on $\mathcal{N}(z_0)$ with $f_{x|z}(\cdot, \cdot|z_{10}, z_{20}) = f_{x_1|z_1}(\cdot|z_{10})f_{x_2|z_2}(\cdot|z_{20})$.
- iii) there exists a sequence $\{z_k : k = 1, 2, 3, ...\}$ of distinct $z_k \in \mathcal{O}$ converging to z_0 such that the sequence $\{f_{x|z}(\cdot, \cdot|z_{1k}, z_{2k}) : k = 1, 2, 3, ...\}$ is linearly independent,

Then, the sequence $\{f_{x|z}(\cdot,\cdot|z_{1k},z_{2k}): k=1,2,3,...\}$ is complete in the Hilbert space \mathcal{H} of functions defined on $\mathcal{X}_1 \times \mathcal{X}_2$.

Proof: See the appendix.

In many applications, it is difficult to show the completeness for a multivariate conditional density. The results above use Theorem 1 to extend the completeness for the one-dimensional sequences $\{f_{x_1|z_1}(\cdot|z_{1k}): k=1,2,3,...\}$ and $\{f_{x_2|z_2}(\cdot|z_{2k}): k=1,2,3,...\}$ to the multiple dimensional sequence $\{f_{x|z}(\cdot,\cdot|z_{1k},z_{2k}): k=1,2,3,...\}$. The key assumption is that the endogenous variables are conditionally independent of each other for some value of the instruments, i.e.

$$f_{x|z}(x_1, x_2|z_{10}, z_{20}) = f_{x_1|z_1}(x_1|z_{10}) f_{x_2|z_2}(x_2|z_{20}).$$

We may then use the completeness of one-dimensional conditional densities $f_{x_1|z_1}(x_1|z_{10})$ and $f_{x_2|z_2}(x_2|z_{20})$ to show the completeness of a multi-dimensional density $f_{x|z}(x_1, x_2|z_{10}, z_{20})$. Therefore, Lemma 6 may reduce the dimension as well as the difficulty of the problem.

4. Conclusion

We provide sufficient conditions for the nonparametric identification of the regression function in a regression model with an endogenous regressor x and an instrumental variable z. The identification of the regression function from the conditional expectation of the dependent variable is implied by the completeness of the distribution of the endogenous regressor conditional on the instrument, i.e., f(x|z). We provide sufficient conditions for the completeness of f(x|z) without imposing a specific functional form, such as the exponential family. We use the results in the stability of bases in Hilbert space to show that if the relative deviation

from a complete sequence of function is finite then f(x|z) itself is complete, and therefore, the regression function is nonparametrically identified.

5. Appendix: Proofs

5.1. Preliminaries

Let $\mathcal{L}^2(\mathcal{X}) = \{h(\cdot) : \int_{\mathcal{X}} |h(x)|^2 dx < \infty, \}$ be a L^2 space with the following inner product $\langle f, g \rangle \equiv \int_{\mathcal{X}} f(x)g(x)dx$. We define the corresponding norm as: $||f||^2 = \langle f, f \rangle$. The completion of $\mathcal{L}^2(\mathcal{X})$ under the norm $||\cdot||$ is a Hilbert space, which may be denoted as \mathcal{H} . The conditional density of interest f(x|z) can be considered as a sequence of functions $\{f_1, f_2, f_3, ...\}$ in \mathcal{H} with

$$f_k \equiv f(\cdot|z_k),$$

where $\{z_k : k = 1, 2, 3, ...\}$ is a sequence in \mathcal{Z} . The property of the sequence $\{f_k\}$ determines the identification of the regression function in (2).

We then introduce the definition of a basis in a Hilbert space.

Definition 2. A sequence of functions $\{f_1, f_2, f_3, ...\}$ in a Hilbert space \mathcal{H} is said to be a basis if for any $h \in \mathcal{H}$ there corresponds a unique sequence of scalars $\{c_1, c_2, c_3, ...\}$ such that

$$h = \sum_{k=1}^{\infty} c_k f_k.$$

The identification of a regression function in Equation (2) actually only requires a sequence $\{f_1, f_2, f_3, ...\}$ containing a basis, instead of a basis itself. Therefore, we consider a complete sequence of functions $\{f_1, f_2, f_3, ...\}$ which satisfies that $\langle g, f_k \rangle = 0$ for k = 1, 2, 3... implies g = 0.

In fact, one can show that a basis is complete and that a complete sequence contains a basis. Since a basis has a unique representation of every element in a Hilbert space, there is redundancy in a complete sequence. Given a complete sequence in a Hilbert space, we can extract a basis from the complete sequence. One of the important properties of a complete

sequence for a Hilbert space is that every element can be approximated arbitrarily close by finite combinations of the elements. We summarize these results as follows.

Lemma 7. (1) A basis in the Hilbert space \mathcal{H} is also a complete sequence.

- (2) Let W be a closed linear subspace of a Hilbert space. Set $W^{\perp} = \{h \in \mathcal{H} : \langle h, g \rangle = 0 \text{ for all } g \in W\}$. Then W^{\perp} is a closed linear subspace such that, $W \bigoplus W^{\perp} = \mathcal{H}$.
- (3) Given a complete sequence of functions $\{f_1, f_2, f_3, ...\}$ in a Hilbert space \mathcal{H} , there exist a subsequence $\{r_1, r_2, r_3, ...\}$ which is a basis in the Hilbert space \mathcal{H} .
- (4) Given a complete sequence of functions $\{f_1, f_2, f_3, ...\}$ in a Hilbert space \mathcal{H} , the completion of the subspace span($\{f_1, f_2, f_3, ...\}$) is \mathcal{H} . That is: for any $h \in \mathcal{H}$ there exists a sequence of scalars $\{c_1, c_2, c_3, ...\}$ such that

$$h = \sum_{k=1}^{\infty} c_k f_k;$$

in other words, $||h - \sum_{k=1}^{n} c_k f_k|| \to 0$ as $n \to \infty$.

Proof of Lemma 7(1): Given a basis $\{f_1, f_2, f_3, ...\}$ in a Hilbert space \mathcal{H} , applying Gram-Schmidt process to the basis yields an orthonormal sequence $\{g_1, g_2, g_3, ...\}$ and $\mathrm{span}(\{f_1, f_2, f_3, ...\}) = \mathrm{span}(\{g_1, g_2, g_3, ...\})$. This implies that $\{g_1, g_2, g_3, ...\}$ is also a basis of the Hilbert space \mathcal{H} and $f = \sum_{k=1}^{\infty} \langle f, g_k \rangle g_k$ for any $f \in \mathcal{H}$. Suppose that $\int f_k(x)h(x)dx = 0$ for all k. It follows that $\langle h, g_k \rangle = 0$ for all k. Thus, $h = \sum_{k=1}^{\infty} \langle h, g_k \rangle g_k = 0$. $\{f_1, f_2, f_3, ...\}$ is a complete sequence. QED.

Proof of Lemma 7(3): We will choose r_k using Gram-Schmidt procedure. First, let $r_1 = f_1$ and $g_1 = \frac{r_1}{\|r_1\|}$. Then $r_2 = f_{s_2}$ where s_2 is the smallest index among $\{2, 3, 4, ...\}$ such that $\tilde{g}_2 \equiv f_{s_2} - \langle f_{s_2}, g_1 \rangle g_1 \neq 0$. Denote $g_2 = \frac{\tilde{g}_2}{\|\tilde{g}_2\|}$. Keep the selection process going, in the k-th step, we have $r_k = f_{s_k}$ where s_k is the smallest index among $\{s_{k-1} + 1, s_{k-1} + 2, s_{k-1} + 3, ...\}$ such that $\tilde{g}_k \equiv f_{s_k} - \sum_{i=1}^{k-1} \langle f_{s_i}, g_i \rangle g_i \neq 0$ and $g_k = \frac{\tilde{g}_k}{\|\tilde{g}_k\|}$. This selection procedure produces three sequences with the same span space, i.e., $\operatorname{span}(\{f_1, f_2, f_3, ...\}) = \operatorname{span}(\{r_1, r_2, r_3, ...\}) = \operatorname{span}(\{g_1, g_2, g_3, ...\})$. In addition, $\{g_1, g_2, g_3, ...\}$ is an orthonormal sequence. To prove $\{r_1, r_2, r_3, ...\}$ is a basis, it is sufficient to show (i) the completion of $\operatorname{span}(\{r_1, r_2, r_3, ...\}) = \mathcal{H}$, and (ii) every finite linear combinations of elements in $\{r_1, r_2, r_3, ...\}$ has a unique representation. Let W be the completion of the subspace $\operatorname{span}(\{r_1, r_2, r_3, ...\})$

under the norm $\|\cdot\|$. Let $W^{\perp}=\{h\in\mathcal{H}:\langle h,g\rangle=0 \text{ for all }g\in W\}$. By Lemma 7, $W\bigoplus W^{\perp}=\mathcal{H}$. Since the sequence $\{f_1,f_2,f_3,...\}$ is complete and $\mathrm{span}(\{f_1,f_2,f_3,...\})=\mathrm{span}(\{r_1,r_2,r_3,...\})$ then $W^{\perp}=\{0\}$ and $W=\mathcal{H}$. On the other hand, suppose that $\sum\limits_{i=1}^n c_k r_k=0$ for some scalars $c_1,...,c_n$. From the selection of r_k , we have $r_k=\sum\limits_{i=1}^{k-1}\langle f_{s_i},g_i\rangle g_i+\|\tilde{g}_k\|g_k\equiv\sum\limits_{i=1}^k a_{ik}g_i$ where $a_{kk}=\|\tilde{g}_k\|\neq 0$. Plugging the expression into the previous equation, $\sum\limits_{i=1}^n c_k r_k=\sum\limits_{i=1}^n c_k \left(\sum\limits_{i=1}^k a_{ik}g_i\right)=0$. Consider the inner products of this term with $g_k, k=1,...,n$. We obtain the system of linear equations

$$\sum_{k=1}^{n} c_k a_{1k} = 0,$$

$$\sum_{k=2}^{n} c_k a_{2k} = 0,$$
...
$$\sum_{k=1}^{n} c_k a_{nk} = 0.$$

The matrix expression of the above system is

$$\begin{bmatrix} a_{11} & a_{12} & \dots & a_{1n} \\ 0 & a_{22} & \dots & a_{2n} \\ 0 & 0 & \dots & \vdots \\ 0 & 0 & 0 & a_{nn} \end{bmatrix} \begin{bmatrix} c_1 \\ c_2 \\ \vdots \\ c_n \end{bmatrix} = 0.$$

Since $a_{kk} = \|\tilde{g}_k\| \neq 0$ for all k, using backward induction results in $c_k = 0$ for all k. This proves the condition (ii). Therefore, the sequence $\{r_1, r_2, r_3, ...\}$ is a basis. QED.

Proof of Lemma 7(4): Let W be the completion of the subspace $span(\{f_1, f_2, f_3, ...\})$ under the norm $\|\cdot\|$. Set $W^{\perp} = \{h \in \mathcal{H} : \langle h, g \rangle = 0 \text{ for all } g \in W\}$. By Lemma 7, $W \bigoplus W^{\perp} = \mathcal{H}$. Since the sequence $\{f_1, f_2, f_3, ...\}$ is complete then $W^{\perp} = \{0\}$ and $W = \mathcal{H}$. QED.

Notice that the sequence of scalars corresponding to a complete sequence in Lemma 7 may not be unique. However, any function f in a Hilbert space can be expressed as a linear

combination of the basis function with a unique sequence of scalars $\{c_1, c_2, c_3, ...\}$. Therefore, we can consider c_n as a function of f. In fact, $c_n(\cdot)$ is the so-called coefficient functional.

Definition 3. If $\{f_1, f_2, f_3, ...\}$ is a basis in a Hilbert space \mathcal{H} , then every function f in \mathcal{H} has a unique series $\{c_1, c_2, c_3, ...\}$ such that

$$f = \sum_{n=1}^{\infty} c_n(f) f_k.$$

Each c_n is a function of f. The functionals c_n (n = 1, 2, 3, ...) are called the coefficient functionals associated with the basis $\{f_1, f_2, f_3, ...\}$.

It is clear that c_n is a linear function of f. Although these functionals are defined in a Hilbert space, they can also be defined in Banach space and are useful tools in Banach space theory. The following results regarding the coefficient functionals are from Theorem 3 in section 6 in Young (1980).

Lemma 8. If $\{f_1, f_2, f_3, ...\}$ is a basis in a Hilbert space \mathcal{H} . Define c_n as coefficient functionals associated with the basis. Then, there exists a constant M such that

$$1 \le ||f_n|| \cdot ||c_n|| \le M,\tag{11}$$

for all n.

In our proofs, we limit our attention to linearly independent sequences when providing sufficient conditions for completeness. We introduce the linear independence of an infinite sequence as follows. We first introduce the finite linear independence

Definition 4. A sequence of functions $\{f_n(\cdot)\}$ of a Hilbert space \mathcal{H} is said to be linearly independent if the equality for any finite K

$$\sum_{i=1}^{K} c_i f_{n_i}(x) = 0 \text{ for all } x \in \mathcal{X}$$

is possible only for $c_i = 0$, (i = 1, 2, ..., K).

The linear independence of an infinite sequence is considered as follows.

Definition 5. A sequence of functions $\{f_n(\cdot)\}\$ of a Hilbert space \mathcal{H} is said to be ω -independent if the equality

$$\sum_{n=1}^{\infty} c_n f_n(x) = 0 \text{ for all } x \in \mathcal{X}$$

is possible only for $c_n = 0$, (n = 1, 2, 3, ...).

It is obvious that the ω -independence implies that linear independence. But the converse argument does not hold. A complete sequence may not be ω -independent, but it contains a basis, and therefore, contains an ω -independent subsequence.

Our proofs also need a uniqueness theorem of complex differential functions. Let w = a+ib, where a, b are real number and $i = \sqrt{-1}$. Define $\mathbb{C} = \{w = a + ib : a, b \in \mathbb{R}\}$ and it is called a complex plane. The complex differential function is defined as follows.

Definition 6. Suppose f is a complex function defined in Ω . If $z_0 \in \Omega$ and

$$\lim_{z \to z_0} \frac{f(z) - f(z_0)}{z - z_0}$$

exists, we denote this limit by $f'(z_0)$ and call it the derivative of f at z_0 . If $f'(z_0)$ exists for every $z_0 \in \Omega$, f is called a complex differential function in Ω .

A complex differential function has a large number of interesting properties which are different from a real differential function. One of them is the following uniqueness theorem, as stated in a corollary on page 209 in Rudin (1987).

Lemma 9. If g and f are complex differential functions in an open connect set Ω and if f(z) = g(z) for all z in some set which has a limit point in Ω , then f(z) = g(z) for all $z \in \Omega$.

5.2. Proofs of completeness of existing sequences

Proof of Lemma 1: In order to use the above uniqueness result of complex differential functions, we consider a converging sequence $\{z_k : k = 1, 2, ...\}$ in \mathcal{Z} as the set with a limit point. Since $\mu(\cdot)$ is continuous with $\mu'(z_0) \neq 0$ for some limit point $z_0 \in \mathcal{Z}$, there exists $\delta > 0$ and a sequence $\{z_k : k = 1, 2, ...\}$ converging to z_0 such that $\{\mu(z_k) : k = 1, 2, ...\} \in (\mu(z_0) - 1)$

 $\delta, \mu(z_0) + \delta \subset \mu(\mathcal{N}(z_0))$ be a sequence of distinct numbers converging to an interior point $\mu(z_0) \in \mu(\mathcal{N}(z_0))$. Suppose that $\int_{-\infty}^{\infty} s(x)t(z)e^{\mu(z_k)\tau(x)}h_0(x)dx = 0$ for some $h_0 \in \mathcal{L}^2(\mathcal{X})$. Notice that t(z) > 0 is irrelevant to completeness, hence we may set t(z) = 1 for simplicity. Consider a complex function with the following form

$$f(w) = \int_{\mathcal{X}} s(x)e^{w\tau(x)}h_0(x)dx,\tag{12}$$

where $w \in \mathbb{C}$ the set of complex numbers. Let w = a + ib, where a, b are real numbers. Applying Cauchy-Schwarz inequality along with the assumption (i) and $h_0 \in \mathcal{L}^2(\mathcal{X})$, we have for $a \in (\mu(z_0) - \delta, \mu(z_0) + \delta)$

$$|f(w)|^{2} \leq \left| \int_{\mathcal{X}} s(x)e^{w\tau(x)}h_{0}(x)dx \right|^{2}$$

$$\leq \left(\int_{\mathcal{X}} s(x)e^{a\tau(x)}|h_{0}(x)|dx \right)^{2}$$

$$\leq \left(\int_{\mathcal{X}} s(x)^{2}e^{2a\tau(x)}dx \right) \left(\int_{\mathcal{X}} h_{0}(x)^{2}dx \right) < \infty.$$
(13)

This suggests that f(w) defined in Equation (12) exists and is finite on the vertical strip $\{w: \mu(z_0) - \delta < \operatorname{Re}(w) < \mu(z_0) + \delta\}$. Since the integration in Equation (12) is with respect to x instead of w, f(w) is a complex differential function on $\{w: \mu(z_0) - \delta < \operatorname{Re}w < \mu(z_0) + \delta\}$ according to the definition introduced above. The condition $\int_{\mathcal{X}} s(x)e^{\mu(z_k)\tau(x)}h_0(x)dx = 0$ is equivalent to $f(\mu(z_k)) = 0$ by Equation (12). This implies that the complex differential function f is equal to zeros in the sequence $\{\mu(z_1), \mu(z_2), \mu(z_3), ...\}$ which has a limit point $\mu(z_0)$. Applying the uniqueness theorem (Lemma 9) quoted above to f results in f(w) = 0 on $\{w: \mu(z_0) - \delta < \operatorname{Re}(w) < \mu(z_0) + \delta\}$. If \mathcal{X} is a finite domain, we extend h_0 to a function in $\mathcal{L}^2(\mathbb{R})$ by

$$\tilde{h}_0(x) = \begin{cases} h_0 & \text{if } x \in \mathcal{X}, \\ 0 & \text{otherwise.} \end{cases}$$

In particular, choose $w = \mu(z_0) + it$ for any real t, we have

$$f(w) = \int_{\mathcal{X}} s(x)e^{\mu(z_0)\tau(x)}e^{it\tau(x)}h_0(x)dx = 0$$

$$= \int_{-\infty}^{\infty} s(\tau^{-1}(x))e^{\mu(z_0)x}e^{itx}\tilde{h}_0(\tau^{-1}(x))\frac{1}{\tau'(x)}dx$$

$$\equiv \int_{-\infty}^{\infty} e^{itx}\hat{h}_0(x)dx.$$

The second step is due to the monotonicity of $\tau(\cdot)$. The last step implies that the Fourier transform of $\hat{h}_0(x)$ is zero on the whole real line. And Eq. (13) implies $\hat{h}_0 \in \mathcal{L}^1(\mathbb{R})$. By the uniqueness Theorem 9.12 in Rudin (1987) for $\hat{h}_0 \in \mathcal{L}^1(\mathbb{R})$, we have $\hat{h}_0 = 0$ and therefore the function $h_0 = 0$. This shows that the sequence $\{g(\cdot|z_k) = s(\cdot)t(z_k)e^{\mu(z_k)\tau(\cdot)} : k = 1, 2, ...\}$ is complete in $\mathcal{L}^2(\mathcal{X})$. QED.

Proof of Lemma 2: Choose a sequence of distinct numbers $\{z_k\}$ in the support \mathcal{Z} converging to $z_0 \in \mathcal{Z}$. Suppose that $\int_{-\infty}^{\infty} f_{\epsilon}(x-z_k) h_0(x) dx = 0$ for some $h_0 \in \mathcal{L}^2(\mathbb{R})$. Consider

$$g(z) \equiv \int_{\mathcal{X}} h_0(x) f_{\epsilon}(x-z) dx,$$

which is a convolution. Let ϕ_g stands for the Fourier transformation of g as follows:

$$\phi_g(t) = \int_{-\infty}^{\infty} e^{itz} g(z) dz$$

$$= \int_{-\infty}^{\infty} e^{it(x-(x-z))} \int_{\mathcal{X}} h_0(x) f_{\epsilon}(x-z) dx dz$$

$$= \int_{\mathcal{X}} e^{itx} h_0(x) \left(\int_{-\infty}^{\infty} e^{it(z-x)} f_{\epsilon}(-(z-x)) dz \right) dx$$

$$= \phi_{h_0}(t) \phi_{-\epsilon}(t)$$

$$= \phi_{h_0}(t) \phi_{\epsilon}(-t).$$

We have $g(z) = \frac{1}{2\pi} \int_{-\infty}^{\infty} e^{-itz} \phi_{h_0}(t) \phi_{\epsilon}(-t) dt$. We define

$$f(w) = \int_{-\infty}^{\infty} e^{-itw} \phi_{h_0}(t) \phi_{\epsilon}(-t) dt,$$

for

w=z+ib for $z\in\mathbb{R}$ and b around zero.

Consider $|b| < \delta$ for δ in Equation (4), we have

$$\begin{aligned} \left| f(w) \right|^2 &= \left| \int_{-\infty}^{\infty} e^{-itw} \phi_{h_0}(t) \phi_{\epsilon}(-t) dt \right|^2 \\ &\leq \left(\int_{-\infty}^{\infty} \left| \phi_{\epsilon}(-t) \right| e^{bt} \left| \phi_{h_0}(t) \right| dt \right)^2 \\ &\leq \left(\int_{-\infty}^{\infty} \left| \phi_{\epsilon}(-t) \right|^2 e^{2bt} dt \right) \left(\int_{-\infty}^{\infty} \left| \phi_{h_0}(t) \right|^2 dt \right) \\ &\leq C^2 \left(\int_{-\infty}^{\infty} e^{-2(\delta - b)|t|} dt \right) \left(\int_{-\infty}^{\infty} \left| \phi_{h_0}(t) \right|^2 dt \right) < \infty, \end{aligned}$$

by $\phi_{h_0}(t) \in \mathcal{L}^2(\mathbb{R})$ since $h_0 \in \mathcal{L}^2(\mathbb{R})$. Since RHS is finite, then f(w) is analytic (complex differentiable) in $\Omega = \{z + ib : |b| < \delta\}$. Consequently, the fact that f(w) equals zero for a sequence $\{z_1, z_2, z_3, ...\}$ converging to z_0 implies that f(w) is equal to zero in Ω by the uniqueness theorem cited in the proof of Lemma 1. This suggests that f(w) is equal to zero for all w = z on the real line, i.e., $\int_{-\infty}^{\infty} e^{-itz} \phi_{h_0}(t) \phi_{\epsilon}(-t) dt = 0$ for all $z \in \mathbb{R}$. Since $\int_{-\infty}^{\infty} |\phi_{h_0}(t)\phi_{\epsilon}(-t)| dt \leq (\int_{-\infty}^{\infty} |\phi_{h_0}(t)|^2 dt)^{1/2} (\int_{-\infty}^{\infty} |\phi_{\epsilon}(-t)|^2 dt)^{1/2} < \infty$, $\phi_{h_0}(t)\phi_{\epsilon}(-t) \in \mathcal{L}^1(\mathbb{R})$. Thus, the ch.f. $\phi_h(t)\phi_{\epsilon}(-t) = 0$ for all t. By Eq. (4), i.e., $\phi_{\epsilon}(t) \neq 0$, we have $\phi_h(t) = 0$ for all $t \in \mathbb{R}$ so h = 0. The family $\{g(x|z) = f_{\epsilon}(x - z_k) : k = 1, 2, ...\}$ is complete in $\mathcal{L}^2(\mathcal{X})$. QED.

5.3. Proof of Theorem 1

We prove Theorem 1 in four steps:

- 1. We prove that if the total deviation from a basis to an ω independent sequence is finite, then the latter sequence is also a basis. This result is summarized in Lemma 10 as the cornerstone of the proof of Theorem 1.
- 2. Condition ii) implies that the total deviation from a complete sequence $\{g(\cdot|z_k)\}$ to the corresponding sequence $\{f(\cdot|z_k)\}$ is finite in the sense that

$$\sum_{j=1}^{\infty} \frac{\|f(\cdot|z_k) - g(\cdot|z_k)\|}{\|g(\cdot|z_k)\|} < \infty, \tag{14}$$

²See Theorem 9.13 on page 186 in Rudin (1987).

³See Theorem 9.12 on page 185 in Rudin (1987).

3. A linearly independent sequence $\{f(\cdot|z_k)\}$ in a normed space contains an ω - independent subsequence $\{f(\cdot|z_{k_l})\}$. Finally, for a complete sequence $\{g(\cdot|z_{k_l})\}$ and the ω -independent sequence $\{f(\cdot|z_{k_l})\}$, Equation (14) and Lemma 10 imply that the sequence $\{f(\cdot|z_{k_l})\}$ is complete, and therefore, $\{f(\cdot|z): z \in \mathcal{N}(z_0)\}$ is complete.

Step 1: We prove that if the total deviation from a basis to an ω - independent sequence is finite, then the latter sequence is also a basis. This result is summarized in the following lemma as the cornerstone of the proof of Theorem 1.

Lemma 10. Suppose that

- i) the sequence $\{e_n(\cdot): n=1,2,...\}$ is a basis in a Hilbert space \mathcal{H} ;
- ii) the sequence $\{f_{n}\left(\cdot\right):n=1,2,...\}$ in \mathcal{H} is ω -independent;
- iii) $\sum_{n=1}^{\infty} \frac{\|f_n(\cdot) e_n(\cdot)\|}{\|e_n(\cdot)\|} < \infty.$

Then, the sequence $\{f_n(\cdot): n=1,2,...\}$ is a basis in \mathcal{H} .

Proof of Lemma 10: As in the proof of Theorem 15 on page 45 of Young (1980), we consider for any function $f \in \mathcal{H}$

$$f = \sum_{n=1}^{\infty} c_n(f)e_n,$$

where $c_n(f)$ is the so-called coefficient functional corresponding to the basis $\{e_n\}$. It is clear that $c_n(f)$ is a linear function of f. Define an operator $T: \mathcal{H} \to \mathcal{H}$ as

$$Tf = \sum_{n=1}^{\infty} c_n(f) (e_n - f_n).$$

It is clear that T is linear. Since $c_n(e_n) = 1$ and $c_k(e_n) = 0$ for $k \neq n$, we have

$$Te_n = \sum_{n=1}^{\infty} c_n(e_n) (e_n - f_n) = e_n - f_n.$$

By using the triangle inequality and the definition of functional, we have

$$||Tf|| = \left\| \sum_{n=1}^{\infty} c_n(f) (e_n - f_n) \right\|$$

$$\leq \sum_{n=1}^{\infty} ||c_n(f) (e_n - f_n)||$$

$$\leq \left(\sum_{n=1}^{\infty} ||e_n - f_n|| \, ||c_n|| \right) ||f||.$$

Lemma 8 suggests that

$$1 \le ||e_n|| \, ||c_n|| \le M.$$

Therefore, we have

$$||Tf|| \leq \left(\sum_{n=1}^{\infty} \frac{||e_n - f_n||}{||e_n||} ||e_n|| ||c_n||\right) ||f||$$

$$\leq M \left(\sum_{n=1}^{\infty} \frac{||e_n - f_n||}{||e_n||}\right) ||f||$$

The relationship above implies that the linear operator T is bounded if

$$\sum_{n=1}^{\infty} \frac{\|e_n - f_n\|}{\|e_n\|} < \infty,$$

which will be shown in the next step to be implied by condition (ii). We then show that T is a compact operator. Set

$$T_N f = \sum_{n=1}^{N} c_n(f) (e_n - f_n).$$

Since each T_N has finite dimensional range and $||T - T_N|| \to 0$ as $N \to \infty$, T is an compact operator.⁴

 $^{^4}$ If an bounded linear operator T is the limit of operators of finite rank, then T is compact. See Exercise 13 on page 112 in Rudin (1991).

Next, we show that $Ker(I-T) = \{0\}$, i.e., (I-T) is invertible. Consider

$$0 = (I - T) f$$

$$= f - \sum_{n=1}^{\infty} c_n(f) (e_n - f_n)$$

$$= \sum_{n=1}^{\infty} c_n(f) e_n - \sum_{n=1}^{\infty} c_n(f) e_n + \sum_{n=1}^{\infty} c_n(f) f_n$$

$$= \sum_{n=1}^{\infty} c_n(f) f_n$$

Since $\{f_n(\cdot)\}$ is an ω -independent sequence, we have $c_n(f) = 0$ for all n, and therefore, 0 = (I - T)f implies f = 0.

Therefore, T is a compact operator defined in a Hilbert space \mathcal{H} with $Ker(I-T)=\{0\}$. By the Fredholm alternative, this shows that (I-T) is a bounded invertible operator.⁵

Since T is bounded, (I-T) is also bounded. Therefore, we have shown that (I-T) is a bounded invertible operator. Clearly, we have $(I-T)e_n=f_n$. Consider any $h\in\mathcal{H}$. Then, $(I-T)^{-1}h$ has an unique series expression $(I-T)^{-1}h=\sum\limits_{n=1}^{\infty}c_ne_n$ since $\{e_n(\cdot)\}$ is a basis. Since (I-T) is bounded, applying (I-T) to the expression above results in $h=\sum\limits_{n=1}^{\infty}c_nf_n$. The argument above shows that every element $h\in\mathcal{H}$ has a unique series expansion in term of f_n . Thus, $\{f_n(\cdot)\}$ is also a basis for \mathcal{H} . QED

Step 2: We show condition ii) implies Equation (14), i.e.,

$$\sum_{k=1}^{\infty} \frac{\|f(\cdot|z_k) - g(\cdot|z_k)\|}{\|g(\cdot|z_k)\|} < \infty.$$

$$(15)$$

We choose a sequence $\{z_k : k = 1, 2, ...\} \subset \mathcal{N}(z_0)$ converging to $z_0 \in \mathcal{N}(z_0)$. In other words, z_0 is a limit point in $\mathcal{N}(z_0)$. Define

$$D(z) \equiv \frac{\|f(\cdot|z) - g(\cdot|z)\|}{\|g(\cdot|z)\|}$$

⁵See the Fredholm alternative in Rudin (1991), Exercise 13 on page 112.

for z close to z_0 . Condition ii) imply that D(z) is continuous at z_0 with $D(z_0) = 0$. Then, we have for some constant C and z close to z_0

$$\frac{\|f(\cdot|z) - g(\cdot|z)\|}{\|g(\cdot|z)\|} = D(z) - D(z_0) \le C |z - z_0|.$$

Therefore, we may choose $|z_k - z_0| = O(k^{-p})$ for p > 1 so that Equation (14) holds with $\sum_{k=1}^{\infty} D(z_k) = O(\sum_{k=1}^{\infty} k^{-p}) < \infty$. Thus, there exists a sequence $\{z_k : k = 1, 2, ...\}$ converging to z_0 such that Equation (14) holds.

Step 3: Condition iii) implies that there exists a linearly independent sequence in $\{f(\cdot|z_k)\}$. According to the second Theorem in Erdos and Straus (1953), any linearly independent sequence in a normed space contains an ω - independent subsequence. We obtain an ω - independent subsequence $\{f(\cdot|z_{k_l})\}$ in $\{f(\cdot|z_k)\}$.

We then show that the ω - independent subsequence $\{f(\cdot|z_{k_l})\}$ is complete in the Hilbert space \mathcal{H} . Since the sequence $\{z_{k_l}\}$ corresponding to $\{f(\cdot|z_{k_l})\}$ is a subsequence of $\{z_k\}$ and also converges to z_0 , condition i) implies that the corresponding sequence $\{g(\cdot|z_{k_l})\}$ is complete in the Hilbert space defined on \mathcal{X} . The two sequences also satisfies Equation (14) i.e.,

$$\sum_{l=1}^{\infty} \frac{\|f(\cdot|z_{k_l}) - g(\cdot|z_{k_l})\|}{\|g(\cdot|z_{k_l})\|} < \infty.$$
 (16)

Let $\{e_n\}$ denote a basis contained in the complete sequence $\{g(\cdot|z_{k_l})\}$ and $\{f_n\}$ be the corresponding subsequence in $\{f(\cdot|z_{k_l})\}$, which is also ω — independent. Then $\{e_n\}$ and $\{f_n\}$ also satisfies $\sum_{n=1}^{\infty} \frac{\|f_n(\cdot) - e_n(\cdot)\|}{\|e_n(\cdot)\|} < \infty$. Lemma 10 implies that $\{f_n\}$ is a basis and therefore $\{f(\cdot|z_{k_l})\}$ is complete in the Hilbert space \mathcal{H} . Since the sequence $\{z_k\}$ is in $\mathcal{N}(z_0)$, the family $\{f(\cdot|z): z \in \mathcal{N}(z_0)\}$ is complete in the Hilbert space \mathcal{H} . QED.

5.4. Proof of completeness in applications

Proof of Lemma 4: Let $\mathcal{N}(z_0)$ be an open neighborhood of z_0 . Since the characteristic function $\phi_{z_0}(t)$ of $f(\cdot|z_0)$ satisfies Equation (4) in Lemma 2, we may generate a complete sequence $\{g(x|z_k) = f(x - \mu(z_k)|z_0) : k = 1, 2, ...\}$ satisfying condition i) in Lemma 2 with $\mu(z_0) = 0$ and $\mu'(z_0) \neq 0$. We have $f(\cdot|z_0) = g(\cdot|z_0)$ and $\|g(\cdot|z_0)\| > 0$ due to $|\phi_{z_0}(t)| > 0$.

As discussed below Theorem 1, when the Hilbert space \mathcal{H} is the $\mathcal{L}^2(\mathcal{X})$, the relative deviation D(z) is continuous if $||g(\cdot|z_0)|| > 0$ and the first-order derivatives $\frac{\partial}{\partial z} f(\cdot|z)$ and $\frac{\partial}{\partial z} g(\cdot|z)$ are also in $\mathcal{L}^2(\mathcal{X})$ for $z \in \mathcal{N}(z_0)$. This is because the derivative of $||f(\cdot|z)||^2$ w.r.t. z is bounded by the Cauchy-Schwarz inequality as follows:

$$\begin{split} \left| \frac{\partial}{\partial z} \left(\| f(\cdot | z) \|^2 \right) \right| &= \left| \frac{\partial}{\partial z} \int f(x | z)^2 dx \right| \\ &= \left| \int 2 f(x | z) \frac{\partial}{\partial z} f(x | z) dx \right| \\ &\leq 2 \left\| f(\cdot | z) \right\| \left\| \frac{\partial}{\partial z} f(\cdot | z) \right\|. \end{split}$$

For $g(x|z) = f(x - \mu(z)|z_0)$, we have

$$\frac{\partial}{\partial z}g(x|z) = \frac{\partial}{\partial z}f(x - \mu(z)|z_0)$$
$$= \frac{\partial}{\partial x}f(x - \mu(z)|z_0)(-\mu'(z)).$$

The condition ii) of Lemma 4 implies that $\frac{\partial}{\partial z}g(\cdot|z)$ is in $\mathcal{L}^2(\mathcal{X})$ so that the relative deviation $D(z) = \frac{\|f(\cdot|z) - g(\cdot|z)\|}{\|g(\cdot|z)\|}$ is continuous in z. Since $D(z_0) = 0$ by definition, the condition ii) in Theorem 1 holds. Thus, the completeness holds for $\{f(\cdot|z) : z \in \mathcal{N}(z_0)\}$.

We then consider the special case $f(x|z) = \frac{1}{\sigma(z)} f_{\varepsilon} \left(\frac{x - \mu(z)}{\sigma(z)} \right)$. WLOG, we set $\sigma(z_0) = 1$ because we may always redefine $\frac{1}{\sigma(z_0)} f_{\varepsilon} \left(\frac{x}{\sigma(z_0)} \right)$ as $f_{\varepsilon}(x)$. Since $\mu(\cdot)$ is continuous with $\mu'(z_0) \neq 0$, the sequence $\{\mu(z_k) : k = 1, 2, 3, ...\} \subset \mu(\mathcal{N}(z_0))$ may be a distinct sequence converging to $\mu(z_0) \in \mu(\mathcal{N}(z_0))$. Applying the results in Lemma 2 with the sequence $\{\mu(z_k) : k = 1, 2, 3, ...\}$, we may show that $\{g(x|z_k) = f_{\varepsilon}(x - \mu(z_k)) : k = 1, 2, ...\}$ is complete. We then extend the completeness of $\{g(x|z_k) = f_{\varepsilon}(x - \mu(z_k)) : k = 1, 2, ...\}$ to $\{f(x|z_k) = \frac{1}{\sigma(z_k)} f_{\varepsilon} \left(\frac{x - \mu(z_k)}{\sigma(z_k)} \right) : k = 1, 2, ...\}$. Since $\sigma(z_0) = 1$, we have $f(x|z_0) = g(x|z_0)$.

We then check $f(\cdot|z) \in \mathcal{L}^2(\mathbb{R})$ for any $z \in \mathcal{N}(z_0)$. We have for some constant C

$$||f(\cdot|z)|| = \int_{\mathcal{X}} \left| \frac{1}{\sigma(z)} f_{\varepsilon} \left(\frac{x - \mu(z)}{\sigma(z)} \right) \right|^{2} dx$$

$$= \int_{\mathcal{X}} \left| \frac{1}{\sigma(z)} f_{\varepsilon}(\varepsilon) \right|^{2} \sigma(z) d\varepsilon$$

$$\leq \frac{C}{\sigma(z_{0})} \int_{\mathbb{R}} |f_{\varepsilon}(\varepsilon)|^{2} d\varepsilon.$$

The last step is due to the continuity of $\sigma(\cdot)$ and $\sigma(z_0) > 0$. Since $|\phi_{\varepsilon}(t)| < Ce^{-\delta|t|}$, we have $\int_{\mathbb{R}} |\phi_{\varepsilon}(t)|^2 dt < \infty$, the last expression is finite, and therefore, $f(\cdot|z)$ is in $\mathcal{L}^2(\mathbb{R})$ for $z \in \mathcal{N}(z_0)$. In order to show the continuity of D(z), we show $\frac{\partial}{\partial z} f(\cdot|z)$ is also in $\mathcal{L}^2(\mathcal{X})$. We have

$$\frac{\partial}{\partial z} f(x|z) = \frac{-\sigma'(z)}{\sigma^2(z)} f_{\varepsilon} \left(\frac{x - \mu(z)}{\sigma(z)} \right) + f'_{\varepsilon} \left(\frac{x - \mu(z)}{\sigma(z)} \right) \left(\frac{-\mu'(z)}{\sigma^2(z)} \right) + \frac{x - \mu(z)}{\sigma(z)} f'_{\varepsilon} \left(\frac{x - \mu(z)}{\sigma(z)} \right) \left(\frac{-\sigma'(z)}{\sigma^2(z)} \right).$$

The function $\frac{\partial}{\partial z} f(\cdot|z)$ for $z \in \mathcal{N}(z_0)$ is in $\mathcal{L}^2(\mathbb{R})$ because of condition ii'). Therefore, the total deviation

$$D(z) = \frac{\left\| \frac{1}{\sigma(z)} f_{\varepsilon} \left(\frac{x - \mu(z)}{\sigma(z)} \right) - f_{\varepsilon} \left(x - \mu(z) \right) \right\|}{\left\| f_{\varepsilon} \left(x - \mu(z) \right) \right\|}$$

is continuous in z.

We show the linear independence of $\{f(\cdot|z_k)\}$ as follows:

$$\lim_{x \to -\infty} \frac{f(x|z_{k+1})}{f(x|z_k)} = \lim_{x \to -\infty} \frac{\left| \frac{1}{\sigma(z_{k+1})} \right| f_{\varepsilon} \left(\frac{x - \mu(z_{k+1})}{\sigma(z_{k+1})} \right)}{\left| \frac{1}{\sigma(z_k)} \right| f_{\varepsilon} \left(\frac{x - \mu(z_k)}{\sigma(z_k)} \right)},$$

where

$$\begin{split} \frac{f_{\varepsilon}\left(\frac{x-\mu(z_{k+1})}{\sigma(z_{k+1})}\right)}{f_{\varepsilon}\left(\frac{x-\mu(z_{k})}{\sigma(z_{k})}\right)} &= \frac{f_{\varepsilon}\left(\frac{x-\mu(z_{k})}{\sigma(z_{k})} - \left(\frac{x-\mu(z_{k})}{\sigma(z_{k})} - \frac{x-\mu(z_{k+1})}{\sigma(z_{k})}\right)\right)}{f_{\varepsilon}\left(\frac{x-\mu(z_{k})}{\sigma(z_{k})}\right)} \\ &= \frac{f_{\varepsilon}\left(\frac{x-\mu(z_{k})}{\sigma(z_{k})} - \left(\frac{[\sigma(z_{k+1})-\sigma(z_{k})]x-\sigma(z_{k+1})\mu(z_{k})+\sigma(z_{k})\mu(z_{k+1})}{\sigma(z_{k})\sigma(z_{k+1})}\right)\right)}{f_{\varepsilon}\left(\frac{x-\mu(z_{k})}{\sigma(z_{k})}\right)} \\ &< \frac{f_{\varepsilon}\left(\frac{x-\mu(z_{k})}{\sigma(z_{k})} - c\right)}{f_{\varepsilon}\left(\frac{x-\mu(z_{k})}{\sigma(z_{k})}\right)}. \end{split}$$

If $\sigma'(z) = 0$, i.e., $\sigma(z_{k+1}) = \sigma(z_k)$, we may pick z_k such that $\mu(z_{k+1}) > \mu(z_k)$ so that the last inequality holds because $f_{\varepsilon}(x)$ decreases as $x \to -\infty$. If $\sigma'(z) \neq 0$, we may pick z_k such that $\sigma(z_{k+1}) < \sigma(z_k)$ and therefore $\left(\frac{[\sigma(z_{k+1}) - \sigma(z_k)]x - \sigma(z_{k+1})\mu(z_k) + \sigma(z_k)\mu(z_{k+1})}{\sigma(z_k)\sigma(z_{k+1})}\right) > c > 0$ for some constant c as $x \to -\infty$. Therefore, condition iii') implies that condition (2) in Lemma 3 holds.

Finally, Theorem 1 implies that the family $\{f(\cdot|z):z\in\mathcal{N}(z_0)\}$ is complete in $\mathcal{L}^2(\mathbb{R})$.

In fact, the proof of Theorem 1 suggests that the sequence $\{f(\cdot|z_k): k=1,2,...\}$ is also complete. QED.

Proof of Lemma 5: We choose distinct $z_k \uparrow z_0$ such that $|z_k - z_0| < \frac{1}{k^p}$ for some p > 2. We use the complete sequence $\{g(x|z_k) = f_{\varepsilon}(x - \mu(z_k)) : k = 1, 2, ...\}$ with $\mu(z_0) = 0$ and $\mu'(z_0) \neq 0$ from Lemma 4. Condition iii) implies that $g(x|z_0) = f_{\varepsilon}(x) = f(x|z_0)$. We may check that the family $\{f(x|z_k) = \left|\frac{\partial}{\partial x}h^{-1}(z_k,x)\right| f_{\varepsilon}\left(h^{-1}(z_k,x)\right) : k = 1, 2, ...\}$ is in $\mathcal{L}^2(\mathbb{R})$. Consider for some constant c_1 and $z \in \mathcal{N}(z_0)$

$$\int_{\mathbb{R}} |f(x|z)|^2 dx = \int_{\mathbb{R}} \left| \frac{\partial h^{-1}(z,x)}{\partial x} f_{\varepsilon} \left(h^{-1}(z,x) \right) \right|^2 dx$$

$$= \int_{\mathbb{R}} \left| \left(\frac{\partial h(z,\varepsilon)}{\partial \varepsilon} \right)^{-1} f_{\varepsilon} (\varepsilon) \right|^2 \frac{\partial h(z,\varepsilon)}{\partial \varepsilon} d\varepsilon$$

$$= \int_{\mathbb{R}} \left| \frac{\partial h(z,\varepsilon)}{\partial \varepsilon} \right|^{-1} |f_{\varepsilon} (\varepsilon)|^2 d\varepsilon$$

$$\leq c_1 \int_{\mathbb{R}} \left| \frac{\partial h(z_0,\varepsilon)}{\partial \varepsilon} \right|^{-1} |f_{\varepsilon} (\varepsilon)|^2 d\varepsilon$$

$$= \frac{c_1}{C} \int_{\mathbb{R}} |f_{\varepsilon} (\varepsilon)|^2 d\varepsilon < \infty$$

The last step is because conditions i) and ii) implies $\left|\frac{\partial h(z_0,\varepsilon)}{\partial \varepsilon}\right| > C > 0$ and $\int_{\mathbb{R}} |f_{\varepsilon}(\varepsilon)|^2 d\varepsilon < \infty$. That means $f(x|z) \in \mathcal{L}^2(\mathbb{R})$ for $z \in \mathcal{N}(z_0)$.

The condition iii) of Lemma 5 implies that $\frac{\partial}{\partial z} f(\cdot|z)$ and $\frac{\partial}{\partial z} g(\cdot|z)$ are in $\mathcal{L}^2(\mathbb{R})$ so that the relative deviation $D(z) = \frac{\|f(\cdot|z) - g(\cdot|z)\|}{\|g(\cdot|z)\|}$ is continuous in z. Since $D(z_0) = 0$ by definition, the condition ii) in Theorem 1 holds.

We show the linear independence of $\{f(\cdot|z_k)\}$ and the corresponding CDF sequence $\{F(\cdot|z_k)\}$ as follows. We consider

$$\lim_{x \to -\infty} \frac{F_{\varepsilon}(x|z_{k+1})}{F_{\varepsilon}(x|z_{k})} = \lim_{x \to -\infty} \frac{F_{\varepsilon}(h^{-1}(z_{k+1},x))}{F_{\varepsilon}(h^{-1}(z_{k},x))}$$

$$= \lim_{x \to -\infty} \frac{F_{\varepsilon}(h^{-1}(z_{k},x) - (h^{-1}(z_{k},x) - h^{-1}(z_{k+1},x)))}{F_{\varepsilon}(h^{-1}(z_{k},x))}.$$

Since the function $h(z,\varepsilon)$ is strictly increasing in ε for $z \in \mathcal{N}(z_0)$, condition iv) implies that

$$h^{-1}(z_{k}, x) - h^{-1}(z_{k+1}, x)$$

$$\equiv \varepsilon_{k} - h^{-1}(z_{k+1}, h(z_{k}, \varepsilon_{k}))$$

$$= \varepsilon_{k} - h^{-1}(z_{k+1}, h(z_{k+1}, \varepsilon_{k}) + [h(z_{k}, \varepsilon_{k}) - h(z_{k+1}, \varepsilon_{k})]).$$

WLOG, we let $[h(z_k, \varepsilon_k) - h(z_{k+1}, \varepsilon_k)] = c' \neq 0$ for $\varepsilon_k \to \infty$. We have

$$h^{-1}(z_k, x) - h^{-1}(z_{k+1}, x)$$

$$= \varepsilon_k - h^{-1}(z_{k+1}, h(z_{k+1}, \varepsilon_k) + c')$$

$$= \varepsilon_k - h^{-1}(z_{k+1}, h(z_{k+1}, \varepsilon_k)) + c$$

$$= \varepsilon_k - \varepsilon_k + c \neq 0$$

for some constant $c \neq 0$ as $\varepsilon_k \to -\infty$. Given F_{ε} is increasing, we may pick z_k such that c > 0 to have

$$\lim_{x \to -\infty} \frac{F_{\varepsilon}(x|z_{k+1})}{F_{\varepsilon}(x|z_k)} < \lim_{x \to -\infty} \frac{F_{\varepsilon}(h^{-1}(z_k, x) - c)}{F_{\varepsilon}(h^{-1}(z_k, x))} = 0.$$

The last step is because F_{ε} satisfies $\lim_{x\to-\infty}\frac{F_{\varepsilon}(x-c)}{F_{\varepsilon}(x)}=\lim_{x\to-\infty}\frac{f_{\varepsilon}(x-c)}{f_{\varepsilon}(x)}=0$. Therefore, condition (2) in Lemma 3 holds. Theorem 1 then implies that completeness of $\{f(\cdot|z):z\in\mathcal{N}(z_0)\}$ in $\mathcal{L}^2(\mathbb{R})$. QED.

Proof of Lemma 6: Without loss of generality, we consider $x=(x_1,x_2), z=(z_1,z_2),$ $\mathcal{X}=\mathcal{X}_1\times\mathcal{X}_2$, and $\mathcal{Z}=\mathcal{Z}_1\times\mathcal{Z}_2$. Condition i) implies that $\{f_{x_1|z_1}(\cdot|z_{1k}): k=1,2,3,...\}$ and $\{f_{x_2|z_2}(\cdot|z_{2k}): k=1,2,3,...\}$ are complete in their corresponding Hilbert spaces.

We then show the sequence $\{f_{x_1|z_1}(\cdot|z_{1k})f_{x_2|z_2}(\cdot|z_{2k}): k=1,2,3,...\}$ is complete because $\{f_{x_1|z_1}(\cdot|z_{1k}): k=1,2,3,...\}$ and $\{f_{x_2|z_2}(\cdot|z_{2k}): k=1,2,3,...\}$ are complete in corresponding Hilbert spaces. Consider

$$\int \int h(x_1, x_2) f(x_1|z_1) f(x_2|z_2) dx_1 dx_2 = \int \left(\int h(x_1, x_2) f(x_1|z_1) dx_1 \right) f(x_2|z_2) dx_2
\equiv \int h'(x_2, z_1) f(x_2|z_2) dx_2.$$

If the LHS is equal to zero for any $(z_1, z_2) \in \mathcal{Z}_1 \times \mathcal{Z}_2$, then for any given $z_1 \int h'(x_2, z_1) f(x_2|z_2) dx_2$

equals to zero for any z_2 . Since $f(x_2|z_2)$ is complete, we have $h'(x_2,z_1)=0$ for any $x_2\in\mathcal{X}_2$ and any given $z_1\in\mathcal{Z}_1$. Furthermore, for any given $x_2\in\mathcal{X}_2$, $h'(x_2,z_1)=0$ for any $z_1\in\mathcal{Z}_1$ implies $h(x_1,x_2)=0$ for any $x_1\in\mathcal{X}_1$. Therefore, the sequence $\{f_{x_1|z_1}(\cdot|z_{1k})f_{x_2|z_2}(\cdot|z_{2k}):k=1,2,3,...\}$ is complete. We then apply Theorem 1 to show that the sequence $\{f_{x_1,x_2|z_1,z_2}(\cdot,\cdot|z_{1k},z_{2k}):k=1,2,3,...\}$ is complete because it is close to a complete sequence $\{f_{x_1|z_1}(\cdot|z_{1k})f_{x_2|z_2}(\cdot|z_{2k}):k=1,2,3,...\}$. QED.

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