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# MAPS THAT TAKE GAUSSIAN MEASURES TO GAUSSIAN MEASURES

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For Don Burkholder, with apologies for what may be abstract nonsense

ABSTRACT. Given a pair of separable, real Banach spaces E and F and a centered Gaussian measure  $\mu$  on E, one can ask what sort of Borel measurable maps  $\Phi: E \longrightarrow F$  map  $\mu$  to a centered Gaussian measure on F. Obviously, a sufficient condition is that  $\Phi$  be linear. On the other hand, linearity is far more than is really needed. Indeed, it suffices to know that  $\Phi$  has the property that

$$\Phi\left(\frac{x_1+x_2}{\sqrt{2}}\right) = \frac{\Phi(x_1) + \Phi(x_2)}{\sqrt{2}}$$

for  $\mathcal{W}^2$ -almost every  $(x_1, x_2) \in E^2$ . In this article, I will first prove a structure theorem which shows that any map  $\Phi$  which satisfies this property arises from a linear map on the Cameron–Martin space associated with  $\mu$  on E. I will then investigate which linear maps on the Cameron–Martin space determine a  $\Phi$ , and finally I will discuss some of the properties of  $\Phi$  which reflect properties of the linear map from which it is determined.

### 1. Abstract Wiener spaces

In this section, I will summarize a few facts about Gaussian measures on a Banach space. My treatment derives from L. Gross's theory of abstract Wiener space. For more details, I refer the reader to [2], [5], or Chapter 8 in [6]. In particular, it is important to know that any non-degenerate, centered Gaussian measure on a separable, real Banach space can be realized as the measure in an abstract Wiener space.

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I will call the pair (E, H) a *potential abstract Wiener space* if H is a real, separable Hilbert space, E is a real separable Banach space, and H is continuously embedded in E as a dense subspace. The following lemma summarizes some elementary facts about potential abstract Wiener spaces.

LEMMA 1.1. Let (H, E) be a potential abstract Wiener space. For each  $x^* \in E^*$  there exists a unique  $h_{x^*} \in H$  with the property that  $\langle h, x^* \rangle = (h, h_{x^*})_H$  for all  $h \in H$ . Moreover,  $x^* \in E^* \longmapsto h_{x^*} \in H$  is a continuous, one-to-one, linear map whose range is dense, and, as a map from  $E^*$  into  $E, x^* \rightsquigarrow h_{x^*}$  is continuous. Finally,  $\{h_{x^*} : x^* \in E^*\}$  contains an orthonormal basis for H.

Given a potential abstract Wiener space (H, E), the triple (H, E, W) is called an *abstract Wiener space* if W is a Borel probability measure on Ewith the property that, for each  $x^* \in E^*$ , the random variable  $x \rightsquigarrow \langle x, x^* \rangle$ under W is a centered Gaussian with variance  $||h_{x^*}||_H^2$ . Equivalently, the Fourier transform  $\widehat{W}$  of W is given by

$$\widehat{\mathcal{W}}(x^*) = \mathbb{E}^{\mathcal{W}}\left[e^{\sqrt{-1}\langle x, x^* \rangle}\right] = e^{-\frac{\|h_{x^*}\|_{H}^2}{2}} \quad \text{for } x^* \in E^*.$$

The Hilbert space H in (H, E, W) is called the *Cameron–Martin space* associated with W on E.

Although the uniqueness of  $\mathcal{W}$  is obvious, its existence is a highly nontrivial matter. Nonetheless, for each H there always exists an E on which there is a  $\mathcal{W}$  for which it is the Cameron–Martin space (i.e.,  $(H, E, \mathcal{W})$  is an abstract Wiener space). When  $N = \dim(H) < \infty$  and one thinks of H as  $\mathbb{R}^{\mathbb{N}}$ with some Hilbert norm, all choices of E can also be identified with  $\mathbb{R}^N$ , and  $\mathcal{W}$  is the distribution of

$$(\xi_1, \dots, \xi_N) \in \mathbb{R}^{\mathbb{N}} \longrightarrow \sum_{k=1}^N \xi_k h_k \quad \text{under } \gamma_{0,1}^N,$$

where  $\gamma_{0,1}$  is the standard Gauss measure on  $\mathbb{R}$  and  $\{h_k : 1 \leq k \leq N\}$  is an orthonormal basis for H. When  $\dim(H) = \infty$ , one has the following criterion for the existence of  $\mathcal{W}$ .

LEMMA 1.2. Let (H, E) be an infinite dimensional (i.e.,  $\dim(H) = \infty$ ) potential abstract Wiener space. Then there exists a W on E for which (H, E, W) is an abstract Wiener space if there exists an orthonormal basis  $\{h_k : k \ge 1\}$  in H for which the series

(1.1) 
$$\sum_{k=1}^{\infty} \xi_k h_k \text{ converges in } E \text{ for } \gamma_{0,1}^{\mathbb{Z}^+} \text{-almost every } (\xi_1, \dots, \xi_k, \dots) \in \mathbb{R}^{\mathbb{Z}^+}$$

Conversely, if (H, E, W) is an abstract Wiener space, then (1.1) holds for every choice of orthonormal basis, the convergence is in  $L^p(\gamma_{0,1}^{\mathbb{Z}^+}; E)$  for every  $p \in [1, \infty)$ , and W is the distribution of the series under  $\gamma_{0,1}^{\mathbb{Z}^+}$ . Given an abstract Wiener space (H, E, W), there is a unique, linear isometric map  $h \in H \longrightarrow \mathcal{I}(h) \in L^2(W; \mathbb{R})$ , known as the *Paley–Wiener map*, such that  $\mathcal{I}(h_{x^*}) = \langle \cdot, x^* \rangle$  *W*-almost surely for each  $x^* \in E^*$ . Indeed, the existence and uniqueness of  $\mathcal{I}$  follow immediately from the facts that  $||h_{x^*}||_H$  is the  $L^2(W; \mathbb{R})$ -norm of  $\langle \cdot, x^* \rangle$  and that  $\{h_{x^*} : x^* \in E^*\}$  is dense in H. Moreover, since  $\langle \cdot, x^* \rangle$  is a centered Gaussian random variable under  $\mathcal{W}$  for each  $x^* \in E^*$ , it follows that  $\mathcal{I}(h)$  under  $\mathcal{W}$  is a centered Gaussian random variable with variance  $||h||_H^2$  for each  $h \in H$ . Hence,  $\{\mathcal{I}(h) : h \in H\}$  is a closed, centered Gaussian family in  $L^2(\mathcal{W}; \mathbb{R})$ .

Recall the (unnormalized) Hermite polynomials  $H_n, n \ge 0$ , given by

$$H_n(\xi) = (-1)^n e^{\frac{\xi^2}{2}} \frac{d^n}{d\xi^n} e^{-\frac{\xi^2}{2}}, \quad \xi \in \mathbb{R}.$$

Familiar facts about these polynomials are that

(1.2) 
$$(H_m, H_n)_{L^2(\gamma_{0,1};\mathbb{R})} = m! \delta_{m,n} \text{ and } \frac{dH_m}{d\xi} = mH_{(m-1)^+}$$

and the span of  $\{H_n : n \ge 0\}$  is dense in  $L^2(\gamma_{0,1}; \mathbb{R})$ . Now suppose that  $(H, E, \mathcal{W})$  is an abstract Wiener space with  $\dim(H) = \infty$ , choose an orthonormal basis  $\{h_k : k \ge 1\}$  in H, and, for  $\alpha = (\alpha_1, \ldots, \alpha_k, \ldots) \in \mathbb{N}^{\mathbb{Z}^+}$  with  $\|\alpha\| = \sum_{k=1}^{\infty} \alpha_k < \infty$ , define

$$\mathcal{H}_{\alpha} = \prod_{k=1}^{\infty} H_{\alpha_k} \big( \mathcal{I}(h_k) \big).$$

Then

$$(\mathcal{H}_{\alpha}, \mathcal{H}_{\beta})_{L^{2}(\mathcal{W};\mathbb{R})} = \alpha! \delta_{\alpha,\beta},$$

where  $\alpha! = \prod_{k=1}^{\infty} \alpha_k!$ . Moreover, because  $\{H_n : n \ge 0\}$  is an orthogonal basis for  $L^2(\gamma_{0,1}; \mathbb{R})$ , one can use the results in Lemma 1.2 to check that  $\{\mathcal{H}_{\alpha} : \|\alpha\| < \infty\}$  is an orthogonal basis in  $L^2(\mathcal{W}; \mathbb{R})$ . In particular, if

$$Z^{(n)}(\mathcal{W}) = \overline{\operatorname{span}(\{\mathcal{H}_{\alpha} : \|\alpha\| = n\})}^{L^{2}(\mathcal{W};\mathbb{R})}$$

then  $Z^{(m)}(\mathcal{W}) \perp Z^{(n)}(\mathcal{W})$  for  $m \neq n$  and  $L^2(\mathcal{W}; \mathbb{R}) = \bigoplus_{n=0}^{\infty} Z^{(n)}(\mathcal{W})$ . This is Wiener's decomposition of  $L^2(\mathcal{W}; \mathbb{R})$  into spaces of homogeneous chaos. It is important to recognize that  $Z^{(n)}(\mathcal{W})$  does not depend on the particular choice of orthonormal basis  $\{h_k : k \geq 1\}$  in terms of which it is defined. In particular,  $Z^{(0)}(\mathcal{W})$  consists of the  $\mathcal{W}$ -almost surely constant elements of  $L^2(\mathcal{W}; \mathbb{R})$  and  $Z^{(1)}(\mathcal{W}) = \{\mathcal{I}(h) : h \in H\}.$ 

Now choose  $\{x_k^*: k \ge 1\} \subseteq E^*$  so that  $\{h_k: k \ge 1\}$  is an orthonormal basis for H when  $h_k = h_{x_k^*}$ , and use this basis and the choice of  $\langle \cdot, x_k^* \rangle$  to represent  $\mathcal{I}(h_k)$  to define the  $\mathcal{H}_{\alpha}$ 's. Clearly, each  $\varphi$  from the span of the  $\mathcal{H}_{\alpha}$ 's is a polynomial in variable  $\{\langle \cdot, x_k^* \rangle : \alpha_k \neq 0\}$ , and therefore  $\partial_k \varphi \equiv \frac{d}{dt} \varphi(x + th_k)|_{t=0}$  exists and is a polynomial function of the same variables as  $\varphi$ . Moreover, from (1.2) one sees that if  $\varphi \in Z^{(n)}(\mathcal{W})$  for some  $n \ge 1$ , then  $\partial_k \varphi \in Z^{(n-1)}(\mathcal{W})$  and

(1.3) 
$$\sum_{k=1}^{\infty} \|\partial_k \varphi\|_{L^2(\mathcal{W};\mathbb{R})}^2 = n \|\varphi\|_{L^2(\mathcal{W};\mathbb{R})}^2.$$

Hence, for each  $n \geq 1$ ,  $\partial_k$  admits a unique continuous extension as a linear map from  $Z^{(n)}(\mathcal{W})$  into  $Z^{(n-1)}(\mathcal{W})$  for which (1.3) continues to hold. Similarly, for  $h \in H$ , there is a linear map  $\partial_h$  on the span of the  $\mathcal{H}_{\alpha}$ 's with the properties that  $\partial_h \varphi = \sum_{k=1}^{\infty} (h, h_k)_H \partial_k \varphi$  when  $h \in \text{span}(\{h_k : k \geq 1\})$  and, for each  $n \geq 1$ ,  $\partial_h$ takes  $Z^{(n)}(\mathcal{W})$  into  $Z^{(n-1)}(\mathcal{W})$  with

(1.4) 
$$\|\partial_h \varphi\|^2_{L^2(\mathcal{W};\mathbb{R})} \le n \|h\|^2_H \|\varphi\|^2_{L^2(\mathcal{W};\mathbb{R})} \quad \text{for } \varphi \in Z^{(n)}(\mathcal{W}).$$

Hence, for each  $h \in H$ ,  $\partial_h$  has a unique extension to  $\operatorname{span}(\bigcup_{n=0}^{\infty} Z^{(n)}(\mathcal{W}))$ as a linear operator with the property that (1.4) holds. Moreover,  $\partial_h$  maps  $Z^{(0)}(\mathcal{W})$  to 0 and, when  $n \geq 1$ ,  $Z^{(n)}(\mathcal{W})$  to  $Z^{(n-1)}(\mathcal{W})$ .

## 2. Wiener maps

Given an abstract Wiener space  $(H, E, \mathcal{W})$ , it is clear that  $(E^2, H^2, \mathcal{W}^2)$  is also an abstract Wiener space. Further, if  $S: E^2 \longrightarrow E$  is given by

$$S(x_1, x_2) = \frac{x_1 + x_2}{\sqrt{2}},$$

then  $S_*\mathcal{W}^2 = \mathcal{W}$ .

Now suppose that (H, E, W) is an abstract Wiener space and that F is a second real, separable Banach space. A map  $\Phi : E \longrightarrow F$  is a *Wiener map* if  $\Phi$  is Borel measurable and

(2.1) 
$$\Phi \circ S = \frac{\Phi \circ \pi_1 + \Phi \circ \pi_2}{\sqrt{2}} \quad \mathcal{W}^2\text{-almost surely},$$

where  $\pi_i: E^2 \longrightarrow E$  is the projection map  $\pi_i(y_1, y_2) = y_i$  for  $i \in \{1, 2\}$ . Notice that if  $\Phi: E \longrightarrow F$  is a Wiener map, then  $\Phi_* \mathcal{W}$  is a centered Gaussian measure on F. Indeed, given any  $y^* \in E^*$ , the distribution  $\mu$  of  $\langle \Phi, y^* \rangle$  will satisfy the convolution equation  $\mu = \mu_{2^{-\frac{1}{2}}} \star \mu_{2^{-\frac{1}{2}}}$ , where, for  $\alpha \in \mathbb{R}$ ,  $\mu_{\alpha}$  is the distribution of  $x \rightsquigarrow \alpha x$  under  $\mu$ , and (cf. Exercise 2.3.21 in Chapter 2 of [6]) the only solutions to this equation are centered Gaussians. Hence, by Fernique's theorem,<sup>1</sup>

(2.2) 
$$\mathbb{E}^{\mathcal{W}}\left[e^{\lambda\|\Phi\|_{F}^{2}}\right] < \infty \quad \text{for some } \lambda \in (0,\infty).$$

In addition, if  $\Psi : E \longrightarrow F$  is a second Borel measurable map which is  $\mathcal{W}$ -almost surely equal to  $\Phi$ , then  $\Psi$  is also a Wiener map, and, more generally, a

<sup>&</sup>lt;sup>1</sup> The statement of Fernique's theorem to which I am referring asserts that there is a  $C < \infty$  such that for each R > 0 there is a  $\lambda > 0$  for which  $\mathbb{E}^{\mu}[e^{\lambda \|x\|_{E}^{2}}] \leq C$  whenever  $\mu$  is a centered Gaussian measure on a separable, real Banach space E with  $\mu(\|x\|_{E} \geq R) \leq \frac{1}{4}$ .

Borel measurable  $\Phi$  is a Wiener map if it is the W-almost sure limit of Wiener maps.

In this section, I will investigate the structure of Wiener maps, and, since there is nothing more to say when  $\dim(H) < \infty$ , I will assume throughout that  $\dim(H) = \infty$ .

LEMMA 2.1. If  $\varphi \in Z^{(n)}(\mathcal{W})$  for some  $n \ge 0$ , then  $\varphi \circ S \in Z^{(n)}(\mathcal{W}^2)$  and  $\partial_{(h,-h)}\varphi \circ S = 0$   $\mathcal{W}$ -almost surely for each  $h \in H$ .

*Proof.* Obviously there is nothing to do when n = 0, and so I will assume that  $n \ge 1$ . In addition, since the set of  $\varphi \in Z^{(n)}(\mathcal{W})$  for which these properties hold is a closed subspace of  $L^2(\mathcal{W}; \mathbb{R})$ , it suffices to prove them when  $\varphi = \mathcal{H}_{\alpha}$  for some  $\alpha$  with  $\|\alpha\| = n \ge 1$ . Further, I will assume that the  $\mathcal{H}_{\alpha}$ 's are defined in terms of an orthonormal basis  $\{h_k : k \ge 1\}$  where, for each  $k \ge 1$ ,  $h_k = h_{x_k^*}$  for some  $x_k^* \in E^*$ . Thus,  $\mathcal{H}_{\alpha}$  can be taken to be a polynomial in the variables  $\{\langle \cdot, x_k^* \rangle : \alpha_k \ne 0\}$ .

That  $\partial_{(h,-h)}\mathcal{H}_{\alpha} \circ S = 0$  is essentially trivial. Indeed, for any  $k \geq 1$ ,  $\sqrt{2} \times \partial_{(h_k,-h_k)}\mathcal{H}_{\alpha} \circ S(x_1,x_2)$  equals

$$H_{\alpha_{k}}^{\prime}\left(\frac{\langle x_{1}+x_{2},x_{k}^{*}\rangle}{\sqrt{2}}\right)\prod_{j\neq k}H_{\alpha_{j}}\left(\frac{\langle x_{1}+x_{2},x_{j}^{*}\rangle}{\sqrt{2}}\right)$$
$$-H_{\alpha_{k}}^{\prime}\left(\frac{\langle x_{1}+x_{2},x_{k}^{*}\rangle}{\sqrt{2}}\right)\prod_{j\neq k}H_{\alpha_{j}}\left(\frac{\langle x_{1}+x_{2},x_{j}^{*}\rangle}{\sqrt{2}}\right)=0.$$

To prove that  $\mathcal{H}_{\alpha} \circ S \in Z^{(n)}(\mathcal{W}^2)$  when  $\|\alpha\| = n$ , use the generating function

$$e^{\lambda\xi - \frac{\lambda^2}{2}} = \sum_{n=0}^{\infty} \frac{\lambda^n}{n!} H_n(\xi)$$

to see that

$$H_n\left(\frac{\xi_1+\xi_2}{\sqrt{2}}\right) = 2^{-\frac{n}{2}} \sum_{m=0}^n \binom{n}{m} H_m(\xi_1) H_{n-m}(\xi_2),$$

and from this conclude that

$$\mathcal{H}_{\alpha} \circ S = 2^{-\frac{n}{2}} \sum_{\beta \leq \alpha} {\alpha \choose \beta} \mathcal{H}_{\beta} \circ \pi_1 \mathcal{H}_{\alpha-\beta} \circ \pi_2 \in Z^{(n)}(\mathcal{W}^2),$$

where  $\beta \leq \alpha$  means that  $\beta_k \leq \alpha_k$  for all  $k \geq 1$  and  $\binom{\alpha}{\beta} = \prod_{k=1}^{\infty} \binom{\alpha_k}{\beta_k}$ .

LEMMA 2.2. A Borel measurable  $\varphi : E \longrightarrow \mathbb{R}$  is a Wiener map if and only if  $\varphi = \mathcal{I}(h)$  for some  $h \in H$ .

*Proof.* First, suppose that  $\varphi = \mathcal{I}(h)$ . If  $h = h_{x^*}$  for some  $x^* \in E^*$ , then  $\varphi = \langle \cdot, x^* \rangle$  *W*-almost surely and so, because  $\langle \cdot, x^* \rangle$  is linear and therefore a Wiener map, it follows that  $\varphi$  is a Wiener map also. To extend the result

to general  $h \in H$ , simply remember that the set of  $\mathbb{R}$ -valued Wiener maps is closed in  $L^2(\mathcal{W};\mathbb{R})$ .

Now suppose that  $\varphi$  is an  $\mathbb{R}$ -valued Wiener map. Since  $\varphi$  is a centered Gaussian under  $\mathcal{W}, \varphi \in L^2(\mathcal{W}; \mathbb{R})$ . Now let  $\varphi_n$  denote the orthogonal projection of  $\varphi$  onto  $Z^{(n)}(\mathcal{W})$ . We will know that  $\varphi = \mathcal{I}(h)$  for some  $h \in H$  once we know that  $\varphi_n = 0$  for  $n \neq 1$ . Since  $\mathbb{E}^{\mathcal{W}}[\varphi] = 0$ ,  $\varphi_0 = 0$ . Thus, assume that  $n \ge 2$ . To show that  $\varphi_n = 0$ , I will first show that  $\varphi_n$  is a Wiener map. Indeed, from  $\varphi = \sum_{m=0}^{\infty} \varphi_m$ , we know that  $\varphi \circ S = \sum_{m=0}^{\infty} \varphi_m \circ S$ . Moreover, by Lemma 2.1,  $\varphi_m \circ S \in Z^{(m)}(\mathcal{W}^2)$ , and therefore  $\varphi_n \circ S$  is the projection of  $\varphi \circ S$  onto  $Z^{(n)}(\mathcal{W}^2)$ . At the same time, because  $\varphi \circ S = \frac{\varphi \circ \pi_1 + \varphi \circ \pi_2}{\sqrt{2}} \mathcal{W}^2$ almost surely, it is clear that  $\frac{\varphi_n \circ \pi_1 + \varphi_n \circ \pi_2}{\sqrt{2}}$  is also the projection of  $\varphi \circ S$  onto almost surely, it is clear that  $\frac{\sqrt{2}}{\sqrt{2}}$  is also the projection of  $\varphi \circ S$  onto  $Z^{(n)}(\mathcal{W}^2)$ . Hence  $\varphi_n \circ S = \frac{\varphi_n \circ \pi_1 + \varphi_n \circ \pi_2}{\sqrt{2}} \mathcal{W}^2$ -almost surely. From this and Lemma 2.1, it follows that  $\partial_h(\varphi_n \circ \pi_1) = \partial_h(\varphi_n \circ \pi_2) \mathcal{W}^2$ -almost surely. But  $\partial_h(\varphi_n \circ \pi_1)$  is independent of  $\partial_h(\varphi_n \circ \pi_2)$  under  $\mathcal{W}^2$ , and therefore they can be  $\mathcal{W}^2$ -almost surely equal only if  $\partial_h \varphi_n$  is  $\mathcal{W}$ -almost surely constant. Since this means that  $\partial_h \varphi \in Z^{(0)}(\mathcal{W}) \cap Z^{(n-1)}(\mathcal{W})$  and  $n \geq 2$ , we now know that  $\partial_h \varphi = 0$  W-almost surely. In particular, if  $\{h_k : k \ge 1\}$  is an orthonormal basis in H, then, by (1.3),

$$0 = \sum_{k=1}^{\infty} \|\partial_{h_k} \varphi_n\|_{L^2(\mathcal{W};\mathbb{R})}^2 = n \|\varphi_n\|_{L^2(\mathcal{W};\mathbb{R})}^2,$$

and so  $\varphi_n = 0$  *W*-almost surely.

THEOREM 2.3. If  $\Phi: E \longrightarrow F$  is Borel measurable, then  $\Phi$  is a Wiener map if and only if there is a bounded, linear map  $A: H \longrightarrow F$  such that  $\langle \Phi, y^* \rangle = \mathcal{I}(A^\top y^*)$  W-almost surely for each  $y^* \in F^*$ , where  $A^\top: F^* \longrightarrow H$ is the adjoint of A. Moreover, if A exists, then it is unique, it is continuous from the weak\* topology on H into the strong topology on F, and, for any orthonormal basis  $\{h_k: k \geq 1\}$ ,

$$\Phi = \sum_{m=1}^{\infty} \mathcal{I}(h_k) A h_k \quad \mathcal{W}\text{-almost surely},$$

where the convergence is W-almost sure as well as in  $L^p(W; \mathbb{R})$  for each  $p \in [1, \infty)$ . In particular, if  $F_A$  is the closure in F of the range AH of A, then  $\Phi \in F_A$  W-almost surely.

*Proof.* First, suppose that A exists. Then, by Lemma 2.2, for each  $y^* \in E^*$ ,  $\langle \Phi, y^* \rangle$  is an  $\mathbb{R}$ -valued Wiener map and therefore

$$\langle \Phi \circ S, y^* \rangle = \frac{\langle \Phi \circ \pi_1 + \Phi \circ \pi_2, y^* \rangle}{\sqrt{2}}$$
  $\mathcal{W}^2$ -almost surely.

Hence, since  $F^*$  is separable in the weak<sup>\*</sup> topology,  $\Phi$  is a Wiener map. Furthermore, because

$$\left\langle Ah, y^* \right\rangle = \left( h, A^\top y^* \right)_H = \mathbb{E}^{\mathcal{W}} \left[ \mathcal{I}(h) \mathcal{I} \left( A^\top y^* \right) \right] = \mathbb{E}^{\mathcal{W}} \left[ \mathcal{I}(h) \left\langle \Phi, y^* \right\rangle \right]$$

it is clear that there is at most one choice of A.

Now suppose that  $\Phi$  is a Wiener map. I begin by constructing the map  $A^{\top}: F^* \longrightarrow H$  which will be the adjoint of the A for which we are looking. Namely, given  $y^* \in F^*$ , set  $\varphi = \langle \Phi, y^* \rangle$ . Then, because  $\Phi$  is a Wiener map,  $\varphi$  is an  $\mathbb{R}$ -valued Wiener map. Hence, by Lemma 2.2, there exists a necessarily unique  $A^{\top}y^* \in H$  such that  $\langle \Phi, y^* \rangle = \mathcal{I}(A^{\top}y^*)$   $\mathcal{W}$ -almost surely. By uniqueness,  $A^{\top}$  is linear. Furthermore, because, by (2.2),  $\Phi \in L^2(\mathcal{W}; F)$  and

$$\|A^{\top}y^{*}\|_{H} = \|\mathcal{I}(A^{\top}y^{*})\|_{L^{2}(\mathcal{W};\mathbb{R})} \le \|y^{*}\|_{F^{*}}\|\Phi\|_{L^{2}(\mathcal{W};F)},$$

it is clear that  $A^{\top}$  is bounded from  $F^*$  into H. In fact, if  $y_n^* \longrightarrow 0$  in the weak\* topology on  $F^*$ , then, because, by the uniform boundedness principle,  $C = \sup\{\|y_n^*\|_{F^*} : n \ge 1\} < \infty$  and therefore  $|\langle \Phi, y_n^* \rangle| \le C \|\Phi\|_F \in L^2(\mathcal{W}; \mathbb{R})$ , it follows from Lebesgue's Dominated Convergence theorem that

$$\left\|A^{\top}y_{n}^{*}\right\|_{H}^{2} = \int \left\langle\Phi, y_{n}^{*}\right\rangle^{2} d\mathcal{W} \longrightarrow 0.$$

Hence,  $A^{\top}$  is continuous from the weak<sup>\*</sup> topology on  $F^*$  into the strong topology on H. In particular, this means that if  $h \in H$  and  $y^{**} = (A^{\top})^{\top} h \in F^{**}$ , then

$$\left\langle \boldsymbol{y}_{n}^{*},\boldsymbol{y}^{**}\right\rangle =\left(\boldsymbol{h},\boldsymbol{A}^{\top}\boldsymbol{y}_{n}^{*}\right)_{H}\longrightarrow\boldsymbol{0}$$

when  $\{y_n^* : n \ge 1\} \subseteq F^*$  tends to 0 in the weak\* topology. Thus (cf. Theorem 9 on p. 421 of [1]),  $y^{**} \in F$ , and so  $(A^{\top})^{\top}$  determines a bounded linear map, which I will call A, from H into F, and clearly  $A^{\top}$  is the adjoint of A.

The continuity of A with respect to the weak<sup>\*</sup> topology on H into to the strong topology on F is an immediate consequence of the corresponding continuity property of  $A^{\top}$  and the fact that the closed unit ball in  $F^*$  is weak<sup>\*</sup> compact. To prove the concluding assertion, let  $\{h_k : k \ge 1\}$  be an orthonormal basis for H and take  $\mathcal{F}_n$  to be the  $\sigma$ -algebra generated by  $\{\mathcal{I}(h_k) : 1 \le k \le n\}$ . Then, because  $\Phi \in L^p(\mathcal{W}; F)$  for every  $p \in [1, \infty)$  and the  $\mathcal{W}$ -completion of  $\bigvee_{n=1}^{\infty} \mathcal{F}_n$  contains the Borel field over E, we know that  $\mathbb{E}^{\mathcal{W}}[\Phi|\mathcal{F}_n] \longrightarrow \Phi$  both  $\mathcal{W}$ -almost surely as well as in  $L^p(\mathcal{W}; F)$  for every  $p \in [1, \infty)$ . On the other hand,

for each  $y^* \in F^*$ . Hence, since the weak<sup>\*</sup> topology on  $F^*$  is separable,

$$\mathbb{E}^{\mathcal{W}}[\Phi|\mathcal{F}_n] = \sum_{k=1}^n \mathcal{I}(h_k)Ah_k.$$

### **3.** *A*'s which determine Wiener maps

Given a bounded, linear map  $A: H \longrightarrow F$  and a Wiener map  $\Phi: E \longrightarrow F$ , I will say that  $\Phi$  comes from A if  $\langle \Phi, y^* \rangle = \mathcal{I}(A^\top y^*)$  W-almost surely for each  $y^* \in F^*$ , in which case I will say that A determines  $\Phi$ . Again because  $F^*$  is weak<sup>\*</sup> separable, it is obvious that, up to a set of  $\mathcal{W}$ -measure 0, A can determine at most one  $\Phi$ . On the other hand, the problem of deciding whether a given A determines any  $\Phi$  is much more difficult. Indeed, it turns out to be tantamount to finding out whether a certain potential abstract Wiener space can be made into an abstract Wiener space. To explain this, let  $H_A = AH$  be the range of A, turn  $H_A$  into a Hilbert space with norm  $\|\cdot\|_{H_A}$  determined by  $\|Ah\|_{H_A} = \|h\|_H$  when  $h \perp \text{Null}(A)$ . Next, take  $F_A$  to be the closure of  $H_A$  in F, and turn  $F_A$  into a Banach space by restricting  $\|\cdot\|_F$  to  $F_A$ . Obviously,  $(H_A, F_A)$  is a potential Wiener space. Moreover,  $F_A^*$  can be identified as the quotient space  $F^*/\sim$ , where  $z^*\sim y^*$  means that  $\langle y, z^* - y^* \rangle = 0$  for all  $y \in F_A$ . Finally, if  $y^* \in F^*$  and  $[y^*]_A$  is the ~-equivalence class containing  $y^*$ , then  $AA^{\top}y^*$  is the unique  $g_A \in H_A$  with the property that  $\langle h_A, [y^*]_A \rangle = (g_A, h_A)_{H_A}$  for all  $h_A \in H_A$ . Thus,  $(h_A)_{[y^*]_A} = AA^\top y^*$ , and therefore  $||(h_A)_{[y^*]_A}||_{H_A} = ||A^{\top}y^*||_H$ .

THEOREM 3.1. Refer to the preceding discussion. Then the following are equivalent.

- (1) A determines a Wiener map  $\Phi: E \longrightarrow F$ .
- (2) There is an orthonormal basis  $\{h_k : k \ge 1\}$  in H such that the series

$$\sum_{k=1}^{\infty} \xi_k Ah_k \text{ converges in } F \text{ for } \gamma_{0,1}^{\mathbb{Z}^+} \text{-almost every } (\xi_1, \dots, \xi_k, \dots) \in \mathbb{R}^{\mathbb{Z}^+}$$

(3) There is a  $W_A$  on  $F_A$  for which  $(H_A, F_A, W_A)$  is an abstract Wiener space.

Moreover, if (1) holds and  $\{h_k : k \ge 1\}$  is an orthonormal basis in H, then

$$\Phi = \sum_{k=1}^{\infty} \mathcal{I}(h_k) A h_k \quad \mathcal{W}\text{-almost surely},$$

where the series converges in  $L^p(\mathcal{W}; F_A)$  for every  $p \in [1, \infty)$  as well as  $\mathcal{W}$ almost surely. In addition, the  $\mathcal{W}_A$  in (3) equals the restriction to  $F_A$  of  $\Phi_*\mathcal{W}$ . *Proof.* First suppose that (1) holds, and let  $\{h_k : k \ge 1\}$  be an orthonormal basis in H. Then, by the last part of Theorem 2.3,  $\sum_{k=1}^{\infty} \mathcal{I}(h_k)Ah_k$  converges in F to  $\Phi$   $\mathcal{W}$ -almost surely as well as in  $L^p(\mathcal{W}; F)$  for every  $p \in [1, \infty)$ . Hence, we now know that the concluding assertion is true. Also, because  $\gamma_{0,1}^{\mathbb{Z}^+}$  is the  $\mathcal{W}$ -distribution of  $(\mathcal{I}(h_1), \ldots, \mathcal{I}(h_k), \ldots)$ , we know that (1) implies (2).

Next suppose that (2) holds, and denote by  $\Phi$  the sum of the series. Then, because, by Lemma 2.2, each of the summands is a Wiener map, it follows that  $\Phi$  is also a Wiener map. In addition, for any  $y^* \in F^*$ , *W*-almost surely

$$\langle \Phi, y^* \rangle = \sum_{k=1}^{\infty} \mathcal{I}(h_k) \langle Ah_k, y^* \rangle = \sum_{k=1}^{\infty} \mathcal{I}(h_k) (h_k, A^\top y^*)_H = \mathcal{I}(A^\top y^*).$$

Hence, (2) implies (1).

Next, suppose that (1), and therefore (2), holds. By (2),  $\Phi \in F_A \mathcal{W}$ -almost surely, and so  $\Phi_*\mathcal{W}(F_A) = 1$ . In addition,

$$\mathbb{E}^{\mathcal{W}}\left[e^{\sqrt{-1}\langle\Phi,y^*\rangle}\right] = \mathbb{E}^{\mathcal{W}}\left[e^{\sqrt{-1}\mathcal{I}(A^{\top}y^*)}\right] = e^{-\frac{1}{2}\|A^{\top}y^*\|_{H}^2}.$$

Hence, since  $\|(h_A)_{[y^*]_A}\|_{H_A} = \|A^\top y^*\|_H$ ,  $(H_A, F_A, \Phi_*\mathcal{W} \upharpoonright F_A)$  is an abstract Wiener space, and so (1) implies (3) and the  $\mathcal{W}_A$  in (3) equals  $\Phi_*\mathcal{W} \upharpoonright F_A$ . Conversely, if (3) holds, choose an orthonormal basis  $\{h_k : k \ge 1\}$  in H so that, for each  $k \ge 1$ , either  $Ah_k = 0$  or  $h_k \perp \text{Null}(A)$ , and denote by  $\mathcal{K}$  the set  $\{k : h_k \perp \text{Null}(A)\}$ . Then  $\{Ah_k : k \in \mathcal{K}\}$  is an orthonormal basis in  $H_A$ , and so  $\sum_{k \in \mathcal{K}} \mathcal{I}_A(Ah_k)Ah_k$  is  $\mathcal{W}_A$ -almost surely convergent in  $F_A$ , where  $\mathcal{I}_A$ is the Paley–Wiener map for  $(H_A, F_A, \mathcal{W}_A)$ . Since  $Ah_k = 0$  for  $k \notin \mathcal{K}$  and  $\{\mathcal{I}_A(Ah_k) : k \in \mathcal{K}\}$  under  $\mathcal{W}_A$  are mutually independent standard normal random variables, I have proved that (3) implies (2).

Obviously, if  $A: H \longrightarrow F$  is a linear map which is bounded with respect to  $\|\cdot\|_E$  in the sense that  $\|Ah\|_F \leq C \|h\|_E$  for some  $C < \infty$ , then the unique extension of A as a bounded linear map from E to F is a Wiener map determined by A. I will now discuss a more interesting source of A's which determine Wiener maps. Before stating the main result in this direction, I recall the following familiar fact (cf. [3] and [4]).

LEMMA 3.2. Let  $(\Omega, \mathcal{F}, \mathbb{P})$  be a probability space, assume that the random variables  $\{X_i : i \geq 1\} \subseteq L^2(\mathbb{P}; \mathbb{R})$  span a Gaussian family, and set  $C_{i,j} = \mathbb{E}^{\mathbb{P}}[X_iX_j]$ . Then the joint distribution  $\mu$  on  $\mathbb{R}^{\mathbb{Z}^+}$  of  $\{X_i : i \geq 1\}$  under  $\mathbb{P}$  is absolutely continuous with respect to  $\gamma_{0,1}^{\mathbb{Z}^+}$  if and only if the matrix  $((C_{i,j}))_{i,j\in\mathbb{Z}^+}$  determines a bounded operator C on  $\ell^2(\mathbb{Z}^+;\mathbb{R})$ , Null(C) = 0, and  $\sum_{i,j=1}^{\infty} (\delta_{i,j} - C_{i,j})^2 < \infty$ . Equivalently,  $\mu \ll \gamma_{0,1}^{\mathbb{Z}^+}$  if and only if  $((C_{i,j}))_{i,j\in\mathbb{Z}^+}$  determines a bounded, non-degenerate operator C on  $\ell^2(\mathbb{Z}^+;\mathbb{R})$  such that I - C is Hilbert–Schmidt. THEOREM 3.3. Let  $(E, H, W_H)$  and  $(F, G, W_G)$  be a pair of infinite dimensional abstract Wiener spaces, let  $A : H \longrightarrow G$  be a bounded linear map, and use  $A^{\dagger} : G \longrightarrow H$  to denote its adjoint. If, for some  $\lambda > 0$ ,  $\lambda I - AA^{\dagger}$  is a Hilbert–Schmidt operator on G, then A determines a Wiener map from E to F.

*Proof.* Clearly it suffices to handle the case when  $\lambda = 1$  since the general case reduces to this one when A is replaced by  $\lambda^{-\frac{1}{2}}A$ . Thus, assume that  $\lambda = 1$ .

I will first show that there is an orthonormal basis for G such that  $S = \lim_{m \to \infty} S_m$  exists in  $F \ W_H$ -almost surely when  $S_m = \sum_{i=1}^m \mathcal{I}(A^{\dagger}g_i)g_i$ . To this end, note that, because  $I - AA^{\dagger}$  is Hilbert–Schmidt, Null $(A^{\dagger}) =$  Null $(AA^{\dagger})$  is finite dimensional. Now choose an orthonormal basis for G  $\{g_i : i \geq 1\}$  so that  $g_i \in \text{Null}(A^{\dagger})$  if  $1 \leq i \leq N = \dim(\text{Null}(A^{\dagger}))$  and  $g_i \perp \text{Null}(A^{\dagger})$  if i > N. Obviously,  $S_m = 0$  if  $1 \leq m \leq N$  and

$$S_m = \sum_{i=N+1}^m \mathcal{I}(A^{\dagger}g_i)g_i \quad \text{if } m > N.$$

Moreover, because  $AA^{\dagger}$  is non-degenerate on Null $(A^{\dagger})^{\perp}$  and  $I - AA^{\dagger}$  is Hilbert– Schmidt, Lemma 3.2 says that the joint distribution of  $\{\mathcal{I}(A^{\dagger}g_i): i > N\}$ under  $\mathcal{W}_H$  is absolutely continuous with respect to the joint distribution of  $\{\mathcal{J}(g_i): i > N\}$  under  $\mathcal{W}_G$ , where  $\mathcal{J}$  denotes the Paley–Wiener map for  $(G, F, \mathcal{W}_G)$ . Hence, because  $\sum_{i=1}^{\infty} \mathcal{J}(g_i)g_i$  is  $\mathcal{W}_G$ -almost surely convergent in F, it follows that  $\lim_{m\to\infty} S_m$  converges  $\mathcal{W}_H$ -almost surely in F.

To complete the proof, observe that each  $S_m$  is a centered, F-valued, Gaussian random variable under  $\mathcal{W}_H$ , and therefore that S is also. Thus, by Fernique's theorem, the fact that  $S_m \longrightarrow S$  in  $F \mathcal{W}_H$ -almost surely implies that there exists an  $\varepsilon > 0$  such that  $\sup_{m \ge 1} E^{\mathcal{W}_H}[e^{\varepsilon ||S_m||_F^2}] < \infty$ , and therefore that  $S_m \longrightarrow S$  in  $L^1(\mathcal{W}_H; F)$ . Now let  $\{h_k : k \ge 1\}$  be an orthonormal basis for H, and let  $\mathcal{F}_n$  be the  $\sigma$ -algebra generated by  $\{\mathcal{I}(h_k) : 1 \le k \le n\}$ . Then  $\mathbb{E}^{\mathcal{W}_H}[S|\mathcal{F}_n] = \lim_{m \to \infty} \mathbb{E}^{\mathcal{W}_H}[S_m|\mathcal{F}_n]$  in  $L^1(\mathcal{W}_H; F)$ . At the same time,

$$\mathbb{E}^{\mathcal{W}_H}[S_m|\mathcal{F}_n] = \sum_{i=1}^m \mathbb{E}^{\mathcal{W}_H} \left[ \mathcal{I}(A^{\dagger}g_i) | \mathcal{F}_n \right] g_i = \sum_{i=1}^m \left( \sum_{k=1}^n \left( A^{\dagger}g_i, h_k \right)_H \mathcal{I}(h_k) \right) \\ = \sum_{k=1}^n \mathcal{I}(h_k) \left( \sum_{i=1}^m (Ah_k, g_i)_G g_i \right) \longrightarrow \sum_{k=1}^n \mathcal{I}(h_k) Ah_k \quad \text{in } G \text{ as } m \to \infty.$$

Hence,

$$\mathbb{E}^{\mathcal{W}_H}[S|\mathcal{F}_n] = \sum_{k=1}^n \mathcal{I}(h_k)Ah_k \quad \mathcal{W}_H\text{-almost surely},$$

and therefore  $\sum_{k=1}^{n} \mathcal{I}(h_k) Ah_k$  converges in F to S  $\mathcal{W}_H$ -almost surely.

Now apply Theorem 3.1.

As the next result makes explicit, in Theorem 3.3 the image  $\Phi_* \mathcal{W}_H$  of  $\mathcal{W}_H$ under  $\Phi$  will not be absolutely continuous with respect to  $\mathcal{W}_G$  unless  $\lambda = 1$ .

COROLLARY 3.4. Let  $(H, E, W_H)$  and  $(G, F, W_G)$  be as in the proceeding, and suppose that  $A : H \longrightarrow F$  is a bounded, linear map. Then the following are equivalent.

- (1) A determines a Wiener map  $\Phi: E \longrightarrow F$  and  $\Phi_* \mathcal{W}_H \ll \mathcal{W}_G$ .
- (2) A maps H boundedly into G,  $\text{Null}(A^{\dagger}) = 0$ , and  $I AA^{\dagger}$  is Hilbert-Schmidt on G.

Moreover, if (2), and therefore (1), holds, then  $\mathcal{J}(g) \circ \Phi = \mathcal{I}(A^{\dagger}g) \mathcal{W}_{H}$ -almost surely for each  $g \in G$ , where  $\mathcal{I}$  and  $\mathcal{J}$  denote the Paley–Wiener maps for  $(H, E, \mathcal{W}_{H})$  and  $(G, F, \mathcal{W}_{G})$ , respectively.

Proof. First, assume that (1) holds. Choose  $\{y_i^*: i \ge 1\} \subseteq F^*$  so that  $\{g_i: i \ge 1\}$  is an orthonormal basis for G when  $g_i = g_{y_i^*}$ . Because  $\Phi_* \mathcal{W}_H \ll \mathcal{W}_G$  and  $\langle \Phi, y_i^* \rangle = \mathcal{J}(g_i) \circ \Phi \mathcal{W}_H$ -almost surely, the joint distribution of  $\{\langle \Phi, y_i^* \rangle: i \ge 1\}$  under  $\mathcal{W}_H$  is absolutely continuous with respect to the joint distribution of  $\{\mathcal{I}(A^\top y_i^*): i \ge 1\}$  under  $\mathcal{W}_H$  is absolutely continuous with respect to  $\gamma_{0,1}^{\mathbb{Z}^+}$ . Hence, by Lemma 3.2, if  $C_{i,j} = \langle A^\top y_i^*, A^\top y_j^* \rangle$ , then the matrix  $((C_{i,j}))_{i,j \in \mathbb{Z}^+}$  determines a bounded, non-degenerate operator C on  $\ell^2(\mathbb{Z}^+; \mathbb{R})$  and I - C is Hilbert–Schmidt. Starting from this and taking  $M = 1 + \|I - C\|_{\mathrm{H.S.}}$ , it is easy to check that  $\|A^\top y^*\|_H \le M\|g_{y^*}\|_G$  and therefore that there is a unique bounded, linear  $A^\dagger: G \longrightarrow H$  such that  $A^\dagger g_{y^*} = A^\top y^*$ , and clearly  $\|A^\dagger\|_{\mathrm{op}} \le M$ . Since this means that  $|\langle Ah, y^* \rangle| \le M\|h\|_H \|g_{y^*}\|_G$ , it follows that A maps H boundedly into G and that  $A^\dagger$  is the adjoint of A as a map from H to G. Furthermore, from the properties of C, it is clear that  $AA^\dagger$  is non-degenerate and that  $I - AA^\dagger$  is Hilbert–Schmidt on H.

Now assume that (2) holds. By Theorem 3.3, we know that A determines a Wiener map  $\Phi: E \longrightarrow F$ . Moreover, in the notation used above, the operator C is non-degenerate and I - C is Hilbert–Schmidt. Hence, by Lemma 3.2, the joint distribution of  $\{\mathcal{I}(A^{\top}y_i^*): i \geq 1\}$  under  $\mathcal{W}_H$  is absolutely continuous with respect to  $\gamma_{0,1}^{\mathbb{Z}^+}$ , and so the joint distribution of  $\{\langle \Phi, y_i^* \rangle: i \geq 1\}$  under  $\mathcal{W}_H$  is absolutely continuous with that of  $\{\mathcal{J}(g_i): i \geq 1\}$  under  $\mathcal{W}_G$ . Since this means that the distribution of  $\Phi$  under  $\mathcal{W}_H$  is absolutely continuous with respect to  $\mathcal{W}_G$ ,  $\Phi_*\mathcal{W}_H \ll \mathcal{W}_G$ .

We now know that  $(1) \iff (2)$ . Assume (2) and therefore (1) hold. Given  $g \in G$ , choose  $\{y_n^* : n \ge 1\} \subseteq F^*$  so that  $g_n \equiv g_{y_n^*} \longrightarrow g$  in G. Then,  $\mathcal{J}(g_n) \longrightarrow \mathcal{J}(g)$  in  $L^2(\mathcal{W}_G; \mathbb{R})$ , and so I may and will assume that  $\mathcal{J}(g_n) \longrightarrow \mathcal{J}(g) \mathcal{W}_G$ -almost surely, which, by absolute continuity, means that  $\mathcal{J}(g_n) \circ \Phi \longrightarrow \mathcal{J}(g) \circ \Phi \mathcal{W}_H$ -almost surely. At the same time,  $A^{\dagger}g_n \longrightarrow A^{\dagger}g$  in H, and so I may and

1353

will also assume that  $\mathcal{I}(A^{\dagger}g_n) \longrightarrow \mathcal{I}(A^{\dagger}g) \mathcal{W}_H$ -almost surely. Hence, since  $\mathcal{J}(g_n) \circ \Phi = \mathcal{I}(A^{\dagger}g_n) \mathcal{W}_H$ -almost surely, it follows that  $\mathcal{J}(g) \circ \Phi = \mathcal{I}(A^{\dagger}g) \mathcal{W}_H$ -almost surely.  $\Box$ 

COROLLARY 3.5. Again let  $(H, E, W_H)$  and  $(G, F, W_G)$  be as in Theorem 3.3, and assume that  $\Phi : E \longrightarrow F$  is a Wiener map for which  $\Phi_*W_H \ll W_G$ . In addition, let  $\Psi$  be a Wiener map from F into a third real, separable Banach space K. Then  $\Psi \circ \Phi$  is a Wiener map from E into K. Moreover, if  $\Phi$  and  $\Psi$  are determined by A and B, respectively, then  $\Psi \circ \Phi$  is determined by  $B \circ A$ .

*Proof.* Because  $\Phi_* \mathcal{W}_H \ll \mathcal{W}_G$ , the set of  $(x_1, x_2) \in E^2$  for which

$$\Psi\left(\frac{\Phi(x_1) + \Phi(x_2)}{\sqrt{2}}\right) \neq \frac{\Psi \circ \Phi(x_1) + \Psi \circ \Phi(x_2)}{\sqrt{2}}$$

has  $\mathcal{W}_H^2$ -measure 0. Hence, since

$$\Phi\left(\frac{x_1+x_2}{\sqrt{2}}\right) = \frac{\Phi(x_1) + \Phi(x_2)}{\sqrt{2}}$$

for  $\mathcal{W}_{H}^{2}$ -almost every  $(x_{1}, x_{2}) \in E^{2}$ , it follows that  $\Psi \circ \Phi$  is also a Wiener map.

Next, let  $A: H \longrightarrow F$  be the bounded, linear map which determines  $\Phi$ . By Corollary 3.4, A is a bounded map from H to G and  $\mathcal{J}(g) \circ \Phi = \mathcal{I}(A^{\dagger}g)$  $\mathcal{W}_{H}$ -almost surely for each  $g \in G$ . Hence, if  $B: F \longrightarrow K$  determines  $\Psi$ , then, for any  $z^{*} \in K^{*}$ ,

$$\langle \Psi \circ \Phi, z^* \rangle = \mathcal{J}(B^\top z^*) \circ \Phi = \mathcal{I}(A^\dagger B^\top z^*) \quad \mathcal{W}_H$$
-almost surely.

But  $A^{\dagger}B^{\top} = (B \circ A)^{\top}$ , and so  $B \circ A$  determines  $\Psi \circ \Phi$ .

COROLLARY 3.6. Let  $(H, E, W_H)$  and  $(G, F, W_G)$  be as in Theorem 3.3, suppose that  $\Phi : E \longrightarrow F$  is a Wiener map for which  $\Phi_*W_H \ll W_G$ , and let  $A : H \longrightarrow G$  be the associated map described in (2) of Corollary 3.4. If Ais non-degenerate and  $I - A^{\dagger}A$  is Hilbert–Schmidt on H, then there exists a Wiener map  $\Psi : F \longrightarrow E$  such that  $\Psi_*W_G \ll W_H$  and  $\Psi$  is the inverse of  $\Phi$ in the sense that  $\Psi \circ \Phi(x) = x$  for  $W_H$ -almost every  $x \in E$  and  $\Phi \circ \Psi(y) = y$ for  $W_G$ -almost every  $y \in F$ .

*Proof.* I begin by showing that A has a bounded inverse  $A^{-1}: G \longrightarrow H$ and that  $I - A^{-1}(A^{-1})^{\dagger}$  is Hilbert–Schmidt on H. Indeed, because A is non-degenerate and  $I - A^{\dagger}A$  is Hilbert–Schmidt,  $AA^{\dagger}$  has a pure point spectrum  $\{\lambda_i : i \ge 1\} \subseteq (0, \infty)$  for which 1 is the only possible accumulation point. Hence, there exists an  $\varepsilon > 0$  such that  $\lambda_i \ge \varepsilon$  for all  $i \ge 1$ , and therefore A has a bounded inverse  $A^{-1}: G \longrightarrow H$ . Furthermore, the Hilbert–Schmidt norm of  $I - A^{-1}(A^{-1})^{\dagger}$  is

$$\sum_{i=1}^{\infty} \left(1 - \frac{1}{\lambda_i}\right)^2 \le \varepsilon^{-2} \sum_{i=1}^{\infty} (1 - \lambda_i)^2 < \infty,$$

and so  $I - A^{-1}(A^{-1})^{\dagger}$  is Hilbert–Schmidt on H.

Given the preceding, Corollary 3.4 says that  $A^{-1}$  determines a Wiener map  $\Psi: F \longrightarrow E$  such that  $\Psi_* \mathcal{W}_G \ll \mathcal{W}_H$ , and Corollary 3.5 says that  $\Psi \circ \Phi$  is a Wiener map which is determined by  $A^{-1} \circ A = I$ . Hence, by uniqueness,  $\Psi \circ \Phi$  is  $\mathcal{W}_H$ -almost surely equal to the identity map on E. Similarly,  $\Phi \circ \Psi$  must be  $\mathcal{W}_G$ -almost surely equal to the identity map on F.

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