SCALE-INVARIANT MEASURABILITY IN WIENER SPACE

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The fact that change of scale is a pathological transformation in Wiener space has long been known. For many problems, this pathology causes no special difficulties. However it is sometimes necessary to consider functions of the form $F(\lambda x)$ where λ varies over the positive reals and xvaries over Wiener space. In this setting a variety of conceptual subtleties arise. In this paper we give a framework and several results which prove useful in dealing with these difficulties. In the last section of this paper we discuss several papers in the recent literature in the light of this framework.

Notation and terminology; introduction. Let T = [0, 1]1. and let $C_0(T)$ denote Wiener space, that is, the space of real-valued continuous functions on T which vanish at t = 0 (The notation $C_0(T)$) will never be abbreviated to C_0 . The latter notation will be introduced latter for a certain proper subsets of $C_0(T)$.). Let \mathscr{B} denote the Borel measurable subsets of $C_0(T)$ and let m_1 denote Wiener measure. One can complete $(C_0(T), \mathcal{B}, m_1)$ in the usual way to obtain $(C_0(T), \mathscr{S}_1, m_1)$ where \mathscr{S}_1 is the class of all Wiener measurable sets. Let σ_n be the partition $0 = t_0 < t_1 < \cdots < t_{2^n} = 1$ where $t_k = k/2^n$ for $k = 0, 1, \dots, 2^n$. Given x in $C_0(T)$, let $S_{\sigma_m}(x) = \sum_{k=1}^{2^n} [x(t_k) - \sum_{k=1}^{2^n}$ $x(t_{k-1})]^2$. For $\lambda \ge 0$, let $C_{\lambda} \equiv \{x \text{ in } C_0(T) \colon \lim_{n \to \infty} S_{\sigma_n}(x) = \lambda^2 \}$ and let $D \equiv \{x \text{ in } C_0(T) \colon \lim_{n \to \infty} S_{\sigma_n}(x) \text{ fails to exist} \}.$ Note that $\lambda C_{\mu} = C_{\lambda \mu}$. Clearly D and the sets C_{λ} , $\lambda \geq 0$, are all Borel sets and $C_0(T)$ is the disjoint union of this family of sets.

The key to our discussion is the following result due to Lévy [35] and independently, but later, to Cameron and Martin [8].

THEOREM 1. $m_1(C_1) = 1$.

Lévy actually established a deeper result. He showed that if $\{\pi_n\}$ is any nested sequence of partitions of [0, 1] whose norm approaches zero, then $\lim_{n\to\infty} S_{\pi_n}(x) = 1$ for m_1 -a.e. x in $C_0(T)$. Cameron and Martin [8] and Cameron [2] showed that Theorem 1 has some surprising implications for the scale change and translation transformations in Wiener space.

A subset A of $C_0(T)$ is said to be scale-invariant measurable provided λA is in \mathscr{S}_1 for all $\lambda > 0$. A scale-invariant measurable set N is said to be scale-invariant null provided $m_1(\lambda N) = 0$ for all $\lambda > 0$. A property which holds except on a scale-invariant null set will be said to hold s-almost everywhere (denoted by s-a.e.). Cameron and Martin introdduce the concept of scale-invariant null sets in [8]. The classes of scale-invariant measurable and scale-invariant null sets will be denoted by \mathscr{S} and \mathscr{N} respectively. A function F is said to be scale-invariant measurable provided F is defined on a scale-invariant measurable set and $F(\lambda x)$ is Wiener measurable for every $\lambda > 0$. Two functions F and G on Wiener space are said to be equivalent ($F \approx G$) if they agree s-a.e.. This notion of equivalence was introduced by Cameron and Storvick in [16]. Closely related notations of equivalence were introduced earlier in [4 and 27].

The present authors (and others) have often been perplexed about the exact nature of the above concepts. As long as one considers a fixed scaling in Wiener space, the scale change pathologies pointed out in [8 and 2] cause no problems. However it is often necessary [3, 4, $9 \sim 19$, 21, 24 \sim 34, 38, 39] to consider functions of the form $F(\lambda x)$ where x is in $C_0(T)$ and λ varies over the positive reals. (The natural functions to consider are often somewhat more complicated than $F(\lambda x)$; for example, $F(\lambda^{-1/2}x + \xi)$ where $\lambda > 0$, $x \in C_0(T)$ and ξ is a real number. We will focus our attention on the functions $F(\lambda x)$ since, once, this is understood, it is not difficult to make adjustments to fit the other cases.) Some of the problems that arise with scale changes can be avoided by considering only Borel measurable functions F on the uncompleted Wiener measure space $(C_0(T), \mathcal{Q}, m_1)$. Unfortunately however, one cannot avoid all the scale change pathologies by restricting attention to Borel measurable F. One striking illustration of this is provided by an example and theorem due to Cameron and Storvick [16, pp. 5-7]. They exhibit two Borel measurable functions F and G which agree except on an m_1 -null Borel set and yet their "Fourier-Feynman transforms" are unequal m_1 -a.e.. In contrast, they show that if F and G are equivalent (equal s-a.e.) and if the "Fourier-Feynman transform" of F exists, then the "Fourier-Feynman transform" of G exists and is equivalent to it. In the last section of this paper we will see that related phenomena occur in the setting of the "Feynman integral".

The results in §§2 and 3 on scale change and translations turn out to be quite simple when looked at from the right point of view; in spite of this (perhaps because of this) they provide valuable insight into exactly what scale-invariant measurable sets and scaleinvariant null sets are really like and exactly how they compare to Wiener measurable sets and Wiener null sets respectively. These results allow us to expand on and understand better some of the old results in [8] and [2]. In addition, as we will discuss in the last section of this paper, they help considerably in understanding several recent papers.

2. Scale-invariant measurable sets. Let m_{λ} be the Borel measure given by $m_{\lambda}(B) = m_1(\lambda^{-1}B)$ for B in \mathscr{B} . Since $\lambda^{-1}C_{\lambda} = C_1$, we see by Theorem 1 that m_{λ} is concentrated on the Borel set C_{λ} ; i.e., $m_{\lambda}(C_{\lambda}) = 1$. Let \mathscr{S}_{λ} denote the σ -algebra obtained by completing $(C_0(T), \mathscr{B}, m_{\lambda})$ and let \mathscr{N}_{λ} denote the m_{λ} -null sets. Note that every subset of $C_0(T) \backslash C_{\lambda}$ is in \mathscr{N}_{λ} .

PROPOSITION 2. (i) N is in η_{λ} if and only if $\lambda^{-1}N$ is in \mathcal{N}_{1} ; equivalently, $\mathcal{N}_{\lambda} = \lambda \mathcal{N}_{1}$. (ii) E is in \mathcal{S}_{λ} if and only if $\lambda^{-1}E$ is in \mathcal{S}_{1} ; equivalently, $\mathcal{S}_{2} = \lambda \mathcal{S}_{1}$. (iii) We have $m_{\lambda}(E) = m_{1}(\lambda^{-1}E)$ for E in \mathcal{S}_{λ} .

Proof. (i) Let N be in \mathscr{N}_{λ} . Then $N \subset M$ where M is an m_{λ} -null Borel set. Hence $m_1(\lambda^{-1}M) = m_{\lambda}(M) = 0$ and so $\lambda^{-1}M$ is an m_1 -null Borel set. But then $\lambda^{-1}N \subset \lambda^{-1}M$ is in \mathscr{N}_1 . The converse can be shown in essentially the same way.

(ii) Let E be in \mathscr{S}_{λ} . Then $E = B \cup N$ where B is in \mathscr{B} and N is in \mathscr{N}_{λ} . Then $\lambda^{-1}N$ is in \mathscr{N}_{1} by (i) and so $\lambda^{-1}E = \lambda^{-1}B \cup \lambda^{-1}N$ is in \mathscr{S}_{1} . The rest of (ii) is easily checked.

(iii) Let E be in \mathscr{S}_{λ} . Then $E = B \cup N$ where B is in \mathscr{B} and N is m_{λ} -null. Then

$$m_{\mathfrak{d}}(E) = m_{\mathfrak{d}}(B \cup N) = m_{\mathfrak{d}}(B) = m_{\mathfrak{d}}(\lambda^{-1}B) = m_{\mathfrak{d}}(\lambda^{-1}B \cup \lambda^{-1}N) = m_{\mathfrak{d}}(\lambda^{-1}E)$$
.

PROPOSITION 3. $\mathscr{S} = \bigcap_{\lambda>0} \mathscr{S}_{\lambda}$; $\mathscr{N} = \bigcap_{\lambda>0} \mathscr{N}_{\lambda}$; \mathscr{S} is a σ -algebra of subsets of $C_0(T)$.

Proof. Suppose A is in \mathscr{S} . Then $\lambda^{-1}A$ is in \mathscr{S}_1 for every $\lambda > 0$. Hence A is in $\lambda \mathscr{S}_1 = \mathscr{S}_{\lambda}$ for every $\lambda > 0$. Hence $\mathscr{S} \subset \bigcap_{\lambda>0} \mathscr{S}_{\lambda}$. Reversing the argument we see that $\bigcap_{\lambda>0} \mathscr{S}_{\lambda} \subset \mathscr{S}$. Hence $\mathscr{S} = \bigcap_{\lambda>0} \mathscr{S}_{\lambda}$. Similarly $\mathscr{N} = \bigcap_{\lambda>0} \mathscr{N}_{\lambda}$. \mathscr{S} is a σ -algebra since the intersection of σ -algebras is a σ -algebra.

PROPOSITION 4. (i) *E* is in \mathscr{S} if and only if $E \cap C_{\lambda}$ is in \mathscr{S}_{λ} for every $\lambda > 0$. (ii) *E* is in \mathscr{N} if and only if $E \cap C_{\lambda}$ is in \mathscr{N}_{λ} for every $\lambda > 0$.

Proof. Suppose E is in \mathscr{S} . Let $\lambda > 0$ be given. E is in \mathscr{S}_{λ} and C_{λ} is in $\mathscr{B} \subset \mathscr{S}_{\lambda}$ and so $E \cap C_{\lambda}$ is in \mathscr{S}_{λ} . Conversely, suppose $E \cap C_{\lambda}$ is in \mathscr{S}_{λ} for every $\lambda > 0$. We wish to show that E is in \mathscr{S}_{λ} for every $\lambda > 0$. But $C_0(T) \setminus C_{\lambda}$ is m_{λ} -null and so $E = (E \cap C_{\lambda}) \cup$ $(E \cap (C_0(T) \setminus C_{\lambda}))$ is in \mathscr{S}_{λ} . (ii) is proved in somewhat similar fashion. The next theorem is quite simple, but it gives a very useful characterization of \mathscr{S} and \mathscr{N} in that it shows rather well what scale-invariant measurable sets and scale-invariant null sets are really like and how they compare to Wiener measurable sets and Wiener null sets respectively.

THEOREM 5. (i) E is in S if and only if E has the form

$$(1) E = \left(\bigcup_{\lambda>0} E_{\lambda}\right) \cup L$$

where each E_{λ} is an m_{λ} -measurable subset of C_{λ} and L is an arbitrary subset of $C_0 \cup D$. Further, for E written in this manner, $m_{\lambda}(E) = m_{\lambda}(E_{\lambda})$ for all $\lambda > 0$. (ii) N is in \mathcal{N} if and only if N has the form

$$(2) N = \left(\bigcup_{\lambda>0} N_{\lambda}\right) \cup L$$

where each N_{λ} is an m_{λ} -null subset of C_{λ} and L is an arbitrary subset of $C_0 \cup D$.

Proof. (i) Suppose E is in \mathscr{S} . Let $E_{\lambda} = E \cap C_{\lambda}$ and let $L = E \cap (C_0 \cup D)$. The decomposition follows from Proposition 4 and the fact that $C_0(T)$ is the disjoint union of D and the sets $C_{\lambda}, \lambda \geq 0$. Conversely, suppose that E has the form (1). To show that E is in \mathscr{S} , it suffices by Proposition 3 to show that E is in \mathscr{S}_{λ_0} for every $\lambda_0 > 0$. Now $E \cap C_{\lambda_0}$ is in \mathscr{S}_{λ_0} by assumption and $(\bigcup_{\substack{\lambda \geq 0 \\ \lambda \neq \lambda_0}} E_{\lambda}) \subset L$ is m_{λ_0} -null. Hence E is in \mathscr{S}_{λ_0} as desired. The formula $m_{\lambda}(E) = m_{\lambda}(E_{\lambda})$ follows from the fact that m_{λ} is concentrated on C_{λ} .

(ii) If N is in \mathcal{N} , then N is in \mathcal{S} and so by (i) can be written in the form $N = (\bigcup_{\lambda>0} N_{\lambda}) \cup L$ where each N_{λ} is in \mathcal{S}_{λ} and $L \subset (D \cup C_0)$. We only need to show that $m_{\lambda}(N_{\lambda}) = 0$ for every $\lambda > 0$. But N in \mathcal{N} implies, by Proposition 3, that N is in \mathcal{N}_{λ} . Hence $m_{\lambda}(N) = m_{\lambda}(N_{\lambda}) = 0$ as desired. The converse of (ii) is easily checked.

REMARK. The preceding theorem shows rather strikingly that there are many more Wiener measurable sets than scale-invariant measurable sets: A set E is Wiener measurable if and only if it has the form $E_1 \cup L$ where E_1 is an m_1 -measurable subset of C_1 and L is an arbitrary subset of $(\bigcup_{0 < \lambda \neq 1} C_{\lambda}) \cup D \cup C_0$. Similarly a set is Wiener null if and only if it has the form $N_1 \cup L$ where N_1 is a m_1 -null subset of C_1 and L is an arbitrary subset of $(\bigcup_{0 < \lambda \neq 1} C_2) \cup D \cup C_0$.

Let $0 < t_1 < \cdots < t_n \leq 1$ and let G be a Lebesgue measurable

subset of *n*-dimensional Euclidean space R^n . It is well known that sets of the form

(3)
$$E \equiv \{x \text{ in } C_0(T): (x(t_1), \dots, x(t_n)) \in G\}$$

are Wiener measurable. (In fact E is Wiener measurable if and only if G is Lebesgue measurable [39].) And that they are not Borel measurable if G is not Borel measurable. It is easy to see that such sets E are scale-invariant measurable since, for any $\lambda > 0$, $\lambda E = \{x \in C_0(T): (x(t_1), \dots, x(t_n)) \in \lambda^{-1}G\}$ is Wiener measurable.

PROPOSITION 6. For every $\lambda_0 > 0$, $\mathscr{B} \subseteq \mathscr{S} \subseteq \mathscr{S}_{\lambda_0}$.

Proof. The containments are clear from the fact that $\mathscr{B} \subset \mathscr{S}_{\lambda}$ for every $\lambda > 0$ and the equality $\mathscr{S} = \bigcap_{\lambda>0} \mathscr{S}_{\lambda}$. Let $G \subset \mathbb{R}^n$ be Lebesgue measurable but not Borel measurable. Then E as in (3) is in \mathscr{S} but not in \mathscr{B} . To see that $\mathscr{S} \neq \mathscr{S}_{\lambda_0}$, take $\lambda_1 \neq \lambda_0$ and let E_{λ_1} be a subset of C_{λ_1} which is not m_{λ_1} -measurable. Then E_{λ_1} is in \mathscr{S}_{λ_0} but not in \mathscr{S}_{λ_1} and so not in \mathscr{S} .

The following striking result of Cameron and Martin [8] becomes rather transparent using Theorem 5.

COROLLARY 7. Let f be any function (however nasty) from $(0, \infty)$ to [0, 1]. Then there exists E in S such that $m_1(\lambda E) = f(\lambda)$ for all $\lambda > 0$.

Proof. For each $\lambda > 0$, pick $E_{\lambda} \subset C_{\lambda}$ such that E_{λ} is in \mathscr{S}_{λ} and $m_{\lambda}(E_{\lambda}) = f(\lambda^{-1})$. Then $E \equiv \bigcup_{\lambda>0} E_{\lambda}$ is the desired set since by Proposition 2 and Theorem 5 we have

$$m_1(\lambda E) = m_{\lambda^{-1}}(E) = m_{\lambda^{-1}}(E_{\lambda^{-1}}) = f(\lambda) \;.$$

Our sets C_{λ} , $\lambda \geq 0$, and D depend on the particular sequence of partitions that we choose. If $\pi = \{\pi_1, \pi_2, \dots\}$ denotes another nested sequence of partitions whose norms go to zero, we may let $C_{\lambda}^{\pi} \equiv \{x \in C_0(T): \lim_{n \to \infty} S_{\pi_n}(x) = \lambda^2\}$ and $D^{\pi} \equiv \{x \in C(T): \lim_{n \to \infty} S_{\pi_n}(x) \text{ fails to}$ exist}. Essentially because of Lévy's more general form of Theorem 1, all of the results obtained up to this point, with changes in notation where appropriate, go through. Note however that \mathscr{S}_{λ} , \mathscr{N}_{λ} , m_{λ} , \mathscr{S} and \mathscr{N} are all independent of the sequence of partitions. A set E in \mathscr{S} now has two decompositions according to the two versions of Theorem 5:

$$(4) E = \left(\bigcup_{\lambda>0} E_{\lambda}\right) \cup L = \left(\bigcup_{\lambda>0} E_{\lambda}^{\pi}\right) \cup L^{\pi}$$

where $E_{\lambda}^{\pi} = E \cap C_{\lambda}^{\pi}$ and $L^{\pi} = E \cap (C_0^{\pi} \cup D^{\pi})$. How do these two decompositions relate to one another? The next proposition shows that they agree up to a scale-invariant null set.

PROPOSITION 8. Let E be in \mathcal{S} . Then the two decompositions of E given by (4) and corresponding respectively to our original sequence of partitions and to π have the property that the set

$$(5) \qquad \left(\bigcup_{\lambda>0} E_{\lambda} \bigtriangleup E_{\lambda}^{\pi}\right) \cup (L \bigtriangleup L^{\pi})$$

is scale-invariant null.

Proof. First note that for all $\lambda > 0$

$$egin{aligned} m_\lambda(E_\lambdaackslash E_\lambda^{\pi})&=m_\lambda[(E\cap C_\lambda)ackslash(E\cap C_\lambda^{\pi})]\ &=m_\lambda[E\cap (C_\lambdaackslash C_\lambda^{\pi})]\ &\leq m_\lambda[C_\lambdaackslash C_\lambda^{\pi}]\ &\leq m_\lambda[C_0(T)ackslash C_\lambda^{\pi}]=0\;. \end{aligned}$$

Thus by Theorem 5, the set $\bigcup_{\lambda>0} (E_{\lambda} \setminus E_{\lambda}^{\pi}) \cup (L \setminus L^{\pi})$ is scale-invariant null. In similar fashion one can show that the set $\bigcup_{\lambda>0} (E_{\lambda}^{\pi} \setminus E_{\lambda}) \cup (L^{\pi} \setminus L)$ is scale-invariant null which concludes the proof since

$$egin{aligned} &igcup_{\lambda > 0} \left(E_\lambda igta E_\lambda^{ au}
ight) \cup \left(L igta L^{ au}
ight) \ &= \left\{ igcup_{\lambda > 0} \left(E_\lambda igrak E_\lambda^{ au}
ight) \cup \left(L igla L^{ au}
ight)
ight\} \ \cup \ \left\{ igcup_{\lambda > 0} \left(E_\lambda^{ au} igrak E_\lambda
ight) \cup \left(L^{ au} igrak L
ight)
ight\} \end{aligned}$$

3. Translations. In [2] Cameron used the pathology of scale change transformations in Wiener space to show that almost no translations preserve Wiener measurability. Specifically, he obtained a set E in \mathscr{S}_1 such that $T_yE \equiv E + y$ is not in \mathscr{S}_1 for m_1 -a.e. y in $C_0(T)$. We obtain several facts below which fill in this picture. For example, we will see that if E is in $\mathscr{S}_{\sqrt{2}}$ then the set E + y is in \mathscr{S}_1 for m_1 -a.e. y. More generally, we will see that if E is in $\mathscr{S}_{\sqrt{p^2+q^2}}$, then E + y is in \mathscr{S}_p for m_q -a.e. y. We need the following result.

THEOREM 9. Let p and q be positive numbers. The following assertions are equivalent:

- (a) $f(\sqrt{p^2 + q^2} z)$ is an m_1 -measurable function of z.
- (b) f(z) is an $m_{\sqrt{p^2+q^2}}$ -measurable function of z.
- (c) f(x + y) is an $m_p \times m_q$ -measurable function of x and y.

(d) f(px + qy) is an $m_1 \times m_1$ -measurable function of x and y. If any one (and hence all) of $(a) \sim (d)$ holds, then

$$\int_{C_0(T)} f(\sqrt{p^2 + q^2} z) dm_1(z) \stackrel{*}{=} \int_{C_0(T)} f(z) dm_{\sqrt{p^2 + q^2}}(z)$$

$$\stackrel{*}{=} \int_{C_0(T) \times C_0(T)} f(x + y) d(m_p \times m_q)(x, y)$$

$$\stackrel{*}{=} \int_{C_0(T) \times C_0(T)} f(px + qy) d(m_1 \times m_1)(x, y)$$

where $by \stackrel{\circ}{=} we$ mean that if either side exists, both sides exist and they are equal.

Comments on the proof of Theorem 9. A simple use of the change of variables formula [20, p. 163] shows the equivalence of (a) and (b) and the equivalence of (c) and (d) and the corresponding integration formulas. The integration formulas for Borel measurable f were known to Lévy. The fact that (a) implies (d) and the corresponding integration formula is a corollary of a more general result due to Bearman [1]; this corollary was specifically pointed out in [15; Lemma 2, p. 239]. The fact that (d) implies (a) requires some work and may not have appeared in the literature; it was proved in [23]. We omit the proof since we do not need this particular implication below.

REMARK. There is a natural extension of Theorem 9 involving n positive numbers instead of 2.

The next result follows immediately from Theorem 9, the fact that for Wiener integrals $\int_{C_0(T)} F(-x)dm_1(x) = \int_{C_0(T)} F(x)dm_1(x)$, and the Fubini theorem, if we take f to be the characteristic function of E where E is in $\mathscr{S}_{\sqrt{n^2+a^2}}$.

THEOREM 10. Let p and q be positive numbers. Let E belong to $\mathscr{S}_{\sqrt{p^2+q^2}}$. Then E + y and E - y are in \mathscr{S}_p for m_q -a.e. y and $m_p(E+y)$ and $m_p(E-y)$ are m_p -measurable functions of y. Similarly E+x and E-x are in \mathscr{S}_q for m_p -a.e. x and $m_q(E+x)$ and $m_q(E-x)$ are m_p -measurable functions of x. Furthermore

$$egin{aligned} &\int_{\mathcal{C}_0(T)} m_p(E+y) dm_q(y) = \int_{\mathcal{C}_0(T)} m_p(E-y) dm_q(y) \ &= (m_p imes m_q)(\{(x, \, y) \colon x+y \ is \ in \ E\}) \ &= m_{\sqrt{p^2+q^2}}(E) \ &= \int_{\mathcal{C}_0(T)} m_q(E-x) dm_p(x) \ &= \int_{\mathcal{C}_0(T)} m_q(E+x) dm_p(x) \ . \end{aligned}$$

COROLLARY 11. Let E be in $\mathscr{S}_{\sqrt{2}}$; then E + y is in \mathscr{S}_1 for m_1 -a.e. y.

When p = q = 1, the next corollary is Cameron's earlier result.

COROLLARY 12. Let p and q be positive numbers. The translation map T_y from $(C_0(T), \mathcal{S}_q, m_q)$ to $(C_0(T), \mathcal{S}_p, m_p)$ is m_q -almost never measurability preserving.

Proof. Applying Theorem 10 with $E = C_{\sqrt{p^2+q^2}}$, we see that $m_p(y + C_{\sqrt{p^2+q^2}}) = 1$ for m_q -a.e. y. Let M be a m_p -nonmeasurable subset of C_p . For each y such that $m_p(y + C_{\sqrt{p^2+q^2}}) = 1$, let $M^* = M \cap (y + C_{\sqrt{p^2+q^2}})$. Then M^* is m_p -nonmeasurable. Let $A = M^* - y$. Since $m_q(C_{\sqrt{p^2+q^2}}) = 0$ and $A \subset C_{\sqrt{p^2+q^2}}$, we see that A is in \mathcal{N}_q and hence in \mathcal{S}_q . But $y + A = M^*$ is not in \mathcal{S}_p .

COROLLARY 13. $(m_p \times m_q)(\{(x, y) \text{ in } C_0(T) \times C_0(T): x + y \text{ is in } C_{\sqrt{p^2+q^2}}\}) = 1;$ in particular, x + y is in $C_{\sqrt{2}}$ for $m_1 \times m_1$ -a.e. (x, y). In contrast, $(m_p \times m_q)(\{(x, y) \text{ in } C_0(T) \times C_0(T): x + y \text{ is in } C_2, \lambda \neq \sqrt{p^2 + q^2}\}) = 0.$

Next we give some positive results concerning the translation of scale-invariant measurable sets and scale-invariant null sets.

COROLLARY 14. Let E be in \mathcal{S} . Then for each p > 0, E + y is in \mathcal{S}_p with the exception of at most a scale-invariant null sets of y's.

Proof. Let E be in \mathscr{S} and let p > 0 be given. It suffices to show that for each $\lambda > 0$, E + y is in \mathscr{S}_p for m_{λ} -a.e. y. But E in \mathscr{S} implies that E is in $\mathscr{S}_{\sqrt{p^2+\lambda^2}}$ and so the result follows from Theorem 10.

COROLLARY 15. Let N be in \mathcal{N} . Then for each p > 0, N + y is in \mathcal{N}_p with the exception of at most a scale-invariant null set of y's.

Proof. Let N be in \mathscr{N} and let p > 0 be given. It suffices to show that for each $\lambda > 0$, $m_p(N+y) = 0$ for m_2 -a.e. y. But N in \mathscr{N} implies that

$$\int_{C_0(T)} m_p(N+y) dm_{\lambda}(y) = m_{\sqrt{p^2+\lambda^2}}(N) = 0 \ .$$

Open questions. (i) Suppose E is in \mathcal{S} . Is it the case that

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E + y is in \mathscr{S} for s-a.e. y? (ii) Suppose N is in \mathscr{N} . Is it the case that N + y is in \mathscr{N} for s-a.e. y?

Remark on (ii). Let N be in \mathscr{N} . For each $\lambda > 0$ let $A_{\lambda} = \{y \in C_0(T): m_{\lambda}(N+y) > 0\}$. We know that A_{λ} is in \mathscr{N} for each $\lambda > 0$. The answer to question (ii) would be yes if one could show that $A \equiv \bigcup_{\lambda>0} A_{\lambda}$ was in \mathscr{N} .

Although it will not concern us in this paper, we should mention that there is a small $(m_1$ -measure 0) but useful set of translators for which the translation map on $(C_0(T), \mathcal{S}_1, m_1)$ is measurable. The relevant result is known as the Cameron-Martin translation theorem. For information on this topic, see the papers of Cameron and Martin [6, 7], Maruyama [36] and Segal [37].

4. Scale-invariant measurable functions. In this section we give some simple but useful results about measurable functions. Let $\mathscr{B}(F)(\mathscr{S}(F), \mathscr{S}(F); \lambda > 0)$ denote the class of all real-valued functions which are defined on a Borel (scale-invariant measurable, m_{λ} -measurable respectively) subset of $C_0(T)$ and which are measurable with respect to the σ -algebra $\mathscr{B}(\mathscr{S}, \mathscr{S}_{\lambda}$ respectively). We will only discuss real-valued functions for convenience. However, all the results hold for complex-valued functions as well. The next three propositions follows from Proposition 3, Theorem 9, and Proposition 6 respectively.

PROPOSITION 16. $\mathscr{S}(F) = \bigcap_{\lambda>0} \mathscr{S}_{\lambda}(F).$

PROPOSITION 17. Let F be a function defined on a subset Δ of $C_0(T)$. Given $\lambda > 0$, let F_{λ} be defined on $\lambda^{-1}\Delta$ by $F_{\lambda}(x) = F(\lambda x)$. (i) F_{λ} is in $\mathscr{S}_1(F)$ if and only if F is in $\mathscr{S}_{\lambda}(F)$. (ii) F is in $\mathscr{S}(F)$ if and only if F_{λ} is in $\mathscr{S}_1(F)$ for every $\lambda > 0$; that is, $\mathscr{S}(F)$ is exactly the class of scale-invariant measurable functions defined in §1.

PROPOSITION 18. For every $\lambda_0 > 0$,

$$\mathscr{B}(F) \subsetneq \mathscr{S}(F) \subsetneq \mathscr{S}_{\lambda_0}(F)$$
.

THEOREM 19. Let F be a function with domain Δ . F is s-a.e. defined and in $\mathscr{S}(F)$ if and only if, for each $\lambda > 0$, $F^{2} \equiv F|_{\mathcal{A} \cap C_{\lambda}}$ is m_{2} -a.e. defined and in $\mathscr{S}_{2}(F)$.

Proof. Suppose F is s-a.e. defined and in $\mathscr{S}(F)$. Then by Theorem 5, F is defined except on a scale-invariant null set N =

 $(\bigcup_{\lambda>0} N_{\lambda}) \cup L$ where $N_{\lambda} \subset C_{\lambda}$ with $m_{\lambda}(N_{\lambda}) = 0$ and where $L \subset D \cup C_{0}$. Hence F^{ι} is defined on C_{λ} except on N_{λ} and so F^{λ} is m_{λ} -a.e. defined. Since F is in $\mathscr{S}(F)$, given a Borel subset B of the reals, $F^{-1}(B)$ is in $\mathscr{S} = \bigcap_{\lambda>0} \mathscr{S}_{\lambda}$. Hence $(F^{\iota})^{-1}(B) = F^{-1}(B) \cap C_{\lambda}$ is in \mathscr{S}_{λ} . Hence F^{λ} is in $\mathscr{S}_{\lambda}(F)$.

Conversely, suppose that for each $\lambda > 0$, F^{2} is defined except on an m_{λ} -null set $N_{\lambda} \subset C_{\lambda}$ and F^{2} is in $\mathscr{S}_{\lambda}(F)$. Then F must be defined except on some subset of the scale-invariant null set $(\bigcup_{\lambda>0} N_{\lambda}) \cup (D \cup C_{0})$. Hence $C_{0}(T) \setminus \Delta$ is scale-invariant null and F is s-a.e. defined. Let B be a Borel subset of the real line. To show that $F^{-1}(B)$ is in \mathscr{S} it suffices by Proposition 4 to show that $F^{-1}(B) \cap C_{\lambda}$ is in \mathscr{S}_{λ} for every $\lambda > 0$. But this is so since $F^{-1}(B) \cap C_{\lambda} = (F^{\lambda})^{-1}(B)$ which is in \mathscr{S}_{λ} as desired.

If F is a bounded, Borel measurable function on the reals and if $\{\lambda_n\}$, λ are positive numbers such that $\lambda_n \to \lambda$, then one easily sees that $\int_a^b F(\lambda_n z) dz \to \int_a^b F(\lambda z) dz$. It is tempting to try to make the same argument for functions on Wiener space, but the following example shows that this cannot be done and further illustrates the care that must be taken in dealing with scale changes.

EXAMPLE 20. Let $\lambda > 0$ be given and suppose that $\lambda_n \uparrow \lambda$. Let $F(z) = 1 - \chi_{C_{\lambda}}(z)$. F is bounded and Borel measurable. Now z is in C_1 for m_1 -a.e. z and so $\lambda_n z$ is in $\lambda_n C_1 = C_{\lambda_n} \subset C_0(T) \backslash C_{\lambda}$ for m_1 -a.e. z. Hence $\int_{C_0(T)} F(\lambda_n z) dm_1(z) = 1$ for every n. On the other hand, $\int_{C_0(T)} F(\lambda z) dm_1(z) = 0$ and so

$$\int_{C_0(T)} F(\lambda_n z) dm_1(z) \xrightarrow{\hspace{0.1cm} \# \to \hspace{0.1cm}} \int_{C_0(T)} F(\lambda z) dm_1(z) \ .$$

A positive result along these lines can be obtained by assuming that F is continuous s-a.e..

The next two propositions are easy consequences of Theorem 9.

PROPOSITION 21. Let F be a scale-invariant measurable function. Then for each p > 0, F(px + y) is an m_1 -measurable function of x for s-a.e. y.

Proof. Given p we must show that

{y: F(px + y) fails to be m_1 -measurable as a function of x} is in \mathcal{N} . Hence by Proposition 4, given q > 0 we must show that

 $\{y \text{ in } C_q: F(px + y) \text{ fails to be } m_1\text{-measurable as a function of } x\}$ is in \mathcal{N}_q . Since $C_q = qC_1$ it suffices to show that {y in C_1 : F(px + qy) fails to be m_1 -measurable as a function of x} is in \mathcal{N}_1 . But $F(\sqrt{p^2 + q^2} z)$ is an m_1 -measurable function of z by assumption. Hence by Theorem 9, F(px + qy) is an $m_1 \times m_1$ -measurable function of (x, y). Hence for m_1 -a.e. y or, equivalently, for m_1 a.e. y in C_1 , F(px + qy) is an m_1 -measurable function of x. Hence

 $\{y \text{ in } C_1: F(px + qy) \text{ fails to be } m_1\text{-measurable as a function of } x\}$ is in \mathcal{N}_1 as desired.

PROPOSITION 22. Let F and G be scale-invariant measurable functions that are equal s-a.e.; that is $F \approx G$. Then for every p, q > 0, F(px + qy) = G(px + qy) for $m_1 \times m_1$ -a.e. (x, y).

One can see that several useful functions are s-a.e. defined and scale-invariant measurable by starting with the simple proposition that follows.

PROPOSITION 23. Let $0 < t_1 < t_2 < \cdots < t_n \leq 1$ and let f be a Lebesgue measurable function on \mathbb{R}^n . Let $F(x) = f(x(t_1), \cdots, x(t_n))$. Then F is s-a.e. defied and scale-invariant measurable.

Proof. It is well-known that if f is a.e. defined and Lebesgue measurable, then F is m_1 -a.e. defined and Wiener measurable. The result follows since, for every $\lambda > 0$, multiplication by λ followed by f is a. e. defined and Lebesgue measurable just as f is.

5. Some related papers. In this section we indicate some ways in which the work in this paper provides insight into several recent papers. Except for the first topic below, we will assume that the reader is familiar with the basic definitions. We begin with a discussion of the analytic Wiener and analytic Feynman "integrals" since in this case the necessary definitions can be easily and quickly given. These integrals were introduced by Cameron in [3] and have played a key role in certain later work [4, 16, 34].

Let C^+ denote the set of complex numbers with positive real part. Let F be a function such that the Wiener integral $J(\lambda) = \int_{C_0(T)} F(\lambda^{-1/2}x) dm_1(x)$ exists for a.e. $\lambda > 0$. If there exists a function $J^*(\lambda)$ analytic in C^+ such that $J^*(\lambda) = J(\lambda)$ for a.e. $\lambda > 0$, then we define this "essential analytic extension" of J to be the analytic Wiener integral of F over $C_0(T)$ with parameter λ . For λ in C^+ , we write

$$\int_{C_0(T)}^{anw_\lambda} F(x) dm_1(x) = J^*(\lambda) \; .$$

Let q be a real number $(q \neq 0)$ and let F be a function whose analytic Wiener integral exists in C^+ . Then if the following limit exists, we call it the *analytic Feyman integral* of F over $C_0(T)$ with parameter q, and we write

$$\int_{C_0(T)}^{anf_q} F(x) dm_1(x) = \lim_{\substack{\substack{\lambda \to -iq \\ \lambda \to in \ C^+}}} \int_{C_0(T)}^{anw_\lambda} F(x) dm_1(x) \ .$$

Let $G \equiv 0$. Clearly the analytic Wiener and Feynman integrals of G are 0 for all values of the parameters.

THEOREM 24. There exists F such that $F = G(G \equiv 0)$ m_1 -a.e. but

$$(6) \qquad \qquad \int_{c_0(T)}^{anw_\lambda} F(x) dm_1(x) \neq 0 \quad for \ all \ \lambda \ in \ C^+ \ dm_1(x) \neq 0$$

and

(7)
$$\int_{a_0(T)}^{a_nf_q} F(x)dm_1(x) \neq 0 \quad \text{for all real } q \neq 0.$$

In fact, given any subset Λ of $(0, \infty)$ of Lebesgue measure 0 (for example, the rationals), there exists F such that for every λ in Λ , $F(\lambda^{-1/2}x) = 0$ m₁-a.e. but (6) and (7) hold.

Proof. For the first assertion we take $F(x) \equiv \chi_{C_0(T)\setminus C_1}(x)$. Then F(x) = G(x) m_1 -a.e. since F = G on C_1 . However for $\lambda \neq 1$, $F(\lambda^{-1/2}x) = 1$ m_1 -a.e. and so

(8)
$$\int_{c_0(T)}^{anw_{\lambda}} F(x) dm_1(x) = 1 = \int_{c_0(T)}^{anf_q} F(x) dm_1(x)$$

for all λ in C^+ and all real $q \neq 0$.

To prove the second assertion, let $F(x) \equiv \chi_{C_0(T) \setminus \bigcup \lambda \in A} C_{\lambda^{-1/2}}(x)$. For λ_0 in Λ , $F(\lambda_0^{-1/2}x) = 0$ for m_1 -a.e. x since $\lambda_0^{-1/2}x$ is in $\lambda_0^{-1/2}C_1 = C_{\lambda_0^{-1/2}}$ for m_1 -a.e. x. For λ_0 not in Λ , $F(\lambda_0^{-1/2}x) = 1$ for m_1 -a.e. x. Since Λ has measure 0, we again have (8) for all λ and q.

The following positive result follows easily from the definitions and our earlier considerations.

THEOREM 25. Suppose that Λ is a subset of $(0, \infty)$ of Lebesgue measure 0 and that for λ not in Λ , $H(\lambda^{-1/2}x) = F(\lambda^{-1/2}x)$ for m_1 -a.e. x. Then if $\int_{C_0(T)}^{anw_2} F(x) dm_1(x)$ exists throughout C^+ , $\int_{C_0(T)}^{anw_2} H(x) dm_1(x)$ exists throughout C^+ and equality holds. Also if $\int_{C_0(T)}^{anf_q} F(x) dm_1(x)$ exists, then $\int_{C_0(T)}^{anf_q} H(x) dm_1(x)$ exists and equality holds. COROLLARY 26. If $H \approx F$, the conclusions of Theorem 25 holds.

When they exists, $\int_{C_0(T)}^{anw_\lambda} F(x)dm_1(x)$ is an analytic function in C^+ and the Feynman integrals $\int_{C_0(T)}^{anf_q} F(x)dm_1(x)$ are the boundary values of this analytic function. It is natural to ask if arbitrary analytic functions and their boundary values can arise as analytic Wiener and Feynman integrals. They can as our next result shows.

THEOREM 27. Given an arbitrary analytic function g on C^+ , there exists a scale-invariant measurable function F such that for all λ in C^+

$$\int_{C_0(T)}^{anw_\lambda} F(x) dm_1(x) = g(\lambda) \; .$$

Proof. Let $F(x) = g(\lambda)$ for x in $C_{\lambda^{-1/2}}$, $\lambda > 0$ and let F(x) = 0 for x in $C_0 \cup D$. F is scale-invariant measurable by Theorem 19. The result now follows since for $\lambda > 0$,

$$\int_{C_0(T)} F(\lambda^{-1/2} x) dm_1(x) = \int_{C_1} F(\lambda^{-1/2} x) dm_1(x) = \int_{C_1} g(\lambda) dm_1(x) = g(\lambda) \; .$$

The next result shows that one can get arbitrary analytic functions g (and so arbitrary boundary values) even with very "nasty" F's.

THEOREM 28. Given any subset Λ of $(0, \infty)$ of Lebesgue measure 0 and any analytic function g on C^+ there exists a function F such that $F(\lambda^{-1/2}x)$ fails to be m_1 -measurable for all λ in Λ but $\int_{c_0(T)}^{anw_2} F(x) dm_1(x) = g(\lambda)$ for all λ in C^+ .

Proof. We define F as in Theorem 27 except that for F in $C_{\lambda^{-1/2}}$, λ in Λ , we let $F(x) = \chi_{A_{\lambda}}(x)$ where A_{λ} is an $m_{\lambda^{-1/2}}$ -nonmeasurable subset of $C_{\lambda^{-1/2}}$. Since Λ has measure 0, $\int_{C_0(T)}^{anw_{\lambda}} F(x) dm_1(x) = g(\lambda)$ just as before, but of course $F(\lambda^{-1/2}x)$ is m_1 -nonmeasurable for all λ in Λ .

We have just seen that arbitrary analytic functions g can arise as analytic Wiener integrals. The next theorem is along opposite lines.

THEOREM 29. Given an arbitrary function f on $(0, \infty)$ (however nasty), there exists a scale-invariant measurable function Fsuch that for all $\lambda > 0$

$$\int_{_{G_0(T)}} F(\lambda^{_{-1/2}}x) dm_{_1}(x) = f(\lambda) \; .$$

Proof. Let $F(x) = f(\lambda)$ for λ in $C_{\lambda^{-1/2}}$, $\lambda > 0$ and let F(x) = 0 for x in $C_0 \cup D$ and argue as in Theorem 27.

We next turn to a discussion of the papers [16] and [34] on the Fourier-Feynman transform. We commented in §1 on Cameron and Storvick's example [16] of Borel measurable functions F and G such that $F = G \ m_1$ -a.e. but their Fourier-Feynman transforms T_qF and T_qG are unequal m_1 -a.e. The functions used by Cameron and Storvick were $F \equiv 1$ and $G = \chi_{c_1}$. We remark that the same functions provide a similar kind of example in the setting of the L_p Fourier-Feynman transform [34].

In [16] and [34] all the functions considered are required to be Borel measurable. However a study of the proofs shows that for the class of functions dealt with in detail by those papers, namely functions F as in Proposition 23 and others "built" out of such functions, Lebesgue measurability rather than Borel measurability is the "right" assumption on the functions f. On the other hand, it is clear from the Cameron-Storvick example and other considerations that Wiener measurability is not the appropriate setting for these papers. It turns out that the concept of scale-invariant measurability is precisely the correct one for the theory developed in [16] and [34]. In I ~ IV below we describe the main results of [16] in the scale-invariant measurability setting.

I. If F is scale-invariant measurable and if T_qF (the L_2 analytic Fourier-Feynman transform of F) exists, then T_qF is scale-invariant measurable. If, in addition, G is scale-invariant measurable and $G \approx F$, then T_qG exists and $T_qG \approx T_qF$.

II. Let F be as in Proposition 23 with f Lebesgue measurable and in $L_2(\mathbb{R}^n)$. Then T_qF exists and is scale-invariant measurable.

III. Let $\Delta_n \equiv \{(t_1, \dots, t_n): 0 < t_1 < \dots < t_n \leq 1\}$. Let f be a Lebesgue measurable function on $\Delta_n \times \mathbb{R}^n$ such that $f(t_1, \dots, t_n; \cdot, \dots, \cdot)$ is in $L_2(\mathbb{R}^n)$ for a.e. (t_1, \dots, t_n) in Δ_n and $\sup_{\Delta_n} ||f(t_1, \dots, t_n; \cdot, \dots, \cdot)||_2 < \infty$. Let

$$F(x) = \int \cdots_{t_n} \int f(t_1, \cdots, t_n; x(t_1), \cdots, x(t_n)) dt_1 \cdots dt_n$$

Then F is defined s-a.e. and is scale-invariant measurable. Furthermore $T_{g}F$ exists and is scale-invariant measurable.

IV. One can proceed to build a larger space by taking certain sums of functions as in III. One can show that such functions F are scale-invariant measurable and that T_qF exists and is scale-invariant measurable. We remark that this final class of functions includes functions of the form. $F(x) = \exp\left(\int_0^1 \theta(t, x(t))dt\right)$, a class of functions on Wiener space which is of considerable interest.

We will not give proofs of $I \sim IV$ above. We note however that Proposition 21 plays a key role in the proof of I and Propositions 21 and 23 play major roles in the proof of II. II is then used to prove III which, in turn, is used to prove IV. Finally we remark that improvements similar to $I \sim IV$ can be made throughout [34].

Next we give some discussion and results associated with the operator-valued Feynman integral or Cameron-Storvick function space integral as studied in $[9 \sim 13, 18, 19, 21, 24 \sim 33]$. These results clarify the basic definitions of the theory in a variety of ways and show the necessity of the equivalences introduced in [27, 29, 30, 33]. The results below hold in the $L(L_p, L_{p'})$ theory [33] in general and for all allowable dimensions. However, for simplicity, we give the results just in the $L(L_2, L_2)$ case with dimension 1. See [9] or [33] for the basic definitions.

One considers complex-valued functions F on C(T). Any y in C(T) has a unique decomposition $y = x + \xi$ where x is in $C_0(T)$ and ξ is in \mathbf{R} . Hence it is natural to regard two functions F and G on C(T) as equivalent if $F(x + \xi) = G(x + \xi)$ for a.e. (x, ξ) in $C_0(T) \times \mathbf{R}$. However we will see that this is not the right notion of equivalence for the operator-valued function space integrals $I_{\lambda}(F)$. Let G be identically 0 on C(T). Given any p > 0, $p \neq 1$, we will show in Corollary 31 below that there exists a function F on C(T) such that $F(x + \xi) = G(x + \xi)$ for a.e. (x, ξ) but $I_p(G)$ is the 0 operator whereas $I_p(F)$ is the operator K_p in $L(L_2, L_2)$ defined by

$$(K_p\psi)(\xi) = (p/2\pi)^{1/2} \int_{-\infty}^{\infty} \psi(x) \exp{(-p(x-\xi)^2/2)} dx \; .$$

(See [33, Lemma 1.1] and [24, p. 776] for some discussion concerning this operator; in those papers it was denoted by $C_{2.}$)

THEOREM 30. Let Z be an arbitrary subset of $(0, \infty)$ and let $W = (0, \infty) \setminus Z$. Then there exists a function F_Z on C(T) such that $I_{\lambda}(F_Z)$ exists as a bounded linear operator on $L_2(\mathbf{R})$ for every $\lambda > 0$ and

$$I_{\lambda}(F_{Z}) = egin{pmatrix} 0 & operator, \ \lambda & in \ Z \ K_{\lambda} & , \ \lambda & in \ W \ . \end{cases}$$

Also for every λ in Z and for every ξ in R, $F_z(\lambda^{-1/2}x + \xi) = 0$ for m_1 -a.e. x in $C_0(T)$ while for every λ in W and for every ξ in R, $F_z(\lambda^{-1/2}x + \xi) = 1$ for m_1 -a.e. x.

$$\begin{array}{ll} \textit{Proof.} \quad \text{Let} \ F_{z}(y) = \chi_{\cup_{W} C_{\lambda} - 1/2}(y - y(0)). \quad \text{Then for every } \lambda_{0} > 0, \\ F_{z}(\lambda_{0}^{-1/2}x + \xi) = \chi_{\bigcup_{W} C_{\lambda} - 1/2}(\lambda_{0}^{-1/2}x + \xi - \lambda_{0}^{-1/2}x(0) - \xi) = \chi_{\bigcup_{W} C_{\lambda} - 1/2}(\lambda_{0}^{-1/2}x) \\ &= \chi_{\bigcup_{W} \lambda_{0}^{1/2} C_{\lambda} - 1/2}(x) = \chi_{\bigcup_{W} C_{\lambda} \lambda_{0}^{1/2} \lambda_{\lambda} - 1/2}(x) \ . \end{array}$$

Now for λ_0 in W, C_1 is one of the sets $\{C_{\lambda_0^{1/2}\lambda^{-1/2}}: \lambda \text{ in } W\}$ and so, for every ξ in R, $F_z(\lambda_0^{-1/2}x + \xi) = 1$ for m_1 -a.e. x. Thus for $\psi \in L_2(R)$

$$egin{aligned} &(I_{\lambda_0}(F_Z)\psi)(\xi) = \int_{C_0(T)} F_Z(\lambda_0^{-1/2}x+\xi)\psi(\lambda_0^{-1/2}x(1)+\xi)dm_{ ext{i}}(x) \ &= \int_{C_0(T)}\psi(\lambda_0^{-1/2}x(1)+\xi)dm_{ ext{i}}(x) \ &= (K_{\lambda_0}\psi)(\xi) \end{aligned}$$

as was to be shown. Now let λ_0 be in Z. Then for every ξ in R, $F_Z(\lambda_2^{-1/2}x + \xi) = 0$ for m_1 -a.e. x. Thus

$$(I_{\lambda_0}(F_Z)\psi)(\xi) = \int_{C_0(T)} (0) \psi(\lambda_0^{-1/2} x(1) + \xi) dm_1(x) = 0 \; .$$

Note that $F_Z|_{C_0(T)}$ is \mathscr{S} -measurable. In fact, if either Z or W is countable, it is clear that $F_Z|_{C_0(T)}$ is Borel measurable.

COROLLARY 31. Let Z be any subset of $(0, \infty)$ such that 1 is in Z. Then $F_Z(x + \xi) = G(x + \xi)$ for a.e. (x, ξ) . However for every λ in W, $I_{\lambda}(F_Z) = K_{\lambda}$ whereas $I_{\lambda}(G)$ is the 0 operator. Moreover, if $W \neq \phi$, $I_{\lambda}^{an}(F_Z)$ and $J_q^{an}(F_Z)$ fail to exist whereas $I_{\lambda}^{an}(G)$ and $J_q^{an}(G)$ are all the 0 operators.

Proof. We just need to be comment on the last assertion. It is well-known that K_{λ} is analytic in λ throughout C^+ [24] and is never the 0 operator. For $I_{\lambda}^{an}(F_Z)$ to exist it would have to be an operator-valued analytic function on C^+ such that, (i) $I_{\lambda}^{an}(F_Z) = I_{\lambda}(F_Z)$ is the zero operator for λ in Z and, (ii) $I_{\lambda}^{an}(F_Z) = I_{\lambda}(F_Z) = K_{\lambda}$ for λ in W. But this cannot happen by the Identity theorem [22] since at least one of Z, W has a limit point in $(0, \infty)$. Since $I_{\lambda}^{an}(F_Z)$ fails to exist, the operators $J_{\alpha}^{an}(F_Z)$ cannot possibly exist.

In the next corollary, we explicitly point out that it is possible to have F_z equal to G in a very strong sense and yet fail to have $I_{\lambda}(F_z) = I_{\lambda}(G)$ for some $\lambda > 0$. COROLLARY 32. Let p > 0, $p \neq 1$ and let $W = \{p\}$. Then for every $\lambda > 0$, $\lambda \neq p$, and for every ξ in R, $F_z(\lambda^{-1/2}x + \xi) = G(\lambda^{-1/2}x + \xi)$ for m_1 -a.e. x. However $I_p(F_z) = K_p$ whereas $I_p(G) = 0$. Further the operators $I_{\lambda}^{an}(F_z)$ and $J_q^{an}(F_z)$ fail to exist.

REMARK. If one defined $I_{\lambda}^{an}(F)$ in a slightly different but reasonable way, $I_{\lambda}^{an}(F_Z)$ would exist and be the 0 operator for the F_Z from Corollary 32. This could be accomplished by defining $I_{\lambda}^{an}(F)$ as an operator-valued analytic function on C^+ which agrees with $I_{\lambda}(F)$ for a.e. λ in $(0, \infty)$. Even with this altered definition, one could still get $I_{\lambda}^{an}(F_Z)$ failing to exist simply by taking Z so that both Z and W have nonzero measure.

The following positive results are easily established.

THEOREM 33. Let F and H be functions on C(T).

(a) Fix $\lambda > 0$ and suppose that $I_{\lambda}(F)$ exists and that $H(\lambda^{-1/2}x + \xi) = F(\lambda^{-1/2}x + \xi)$ for a.e. (x, ξ) . Then $I_{\lambda}(H)$ exists and equals $I_{\lambda}(F)$.

(b) Suppose $I_{\lambda}^{an}(F)$ exists for λ in C^+ and that, for every $\lambda > 0$, $H(\lambda^{-1/2}x + \xi) = F(\lambda^{-1/2}x + \xi)$ for a.e. (x, ξ) . Then $I_{\lambda}^{an}(H)$ exists throughout C^+ and equals $I_{\lambda}^{an}(F)$.

(c) Suppose $J_q^{an}(F)$ exists for $q \neq 0$ in R and that for every $\lambda > 0$, $H(\lambda^{-1/2}x + \xi) = F(\lambda^{-1/2}x + \xi)$ for a.e. (x, ξ) . Then $J_q^{an}(H)$ exists and equals $J_q^{an}(F)$.

 $I_{\lambda}^{an}(F)$, when it exists, is an operator-valued analytic function of λ for λ in C^+ . Is there anything special about the analytic functions that arise in this way? For example, can one put any limits on the growth of $||I_{\lambda}^{an}(F)||$ as $\lambda \to -iq$? The following result answers some such questions in the negative.

THEOREM 34. Let $g(\lambda)$ be an arbitrary scalar-valued analytic function on C^+ . Then there exists F such that $I_{\lambda}^{an}(F)$ exists and $I_{\lambda}^{an}(F) = g(\lambda)K_{\lambda}$ for λ in C^+ . For this F, $||I_{\lambda}^{an}(F)|| = |g(\lambda)|$ for all λ in C^+ .

Proof. $g(\lambda)K_{\lambda}$ is an operator-valued analytic function of λ for λ in C^{+} and so it suffices to find F such that $I_{\lambda}(F) = g(\lambda)K_{\lambda}$ for $\lambda > 0$. Let

$$F(y) = egin{cases} g(\lambda) & ext{when } y-y(0) ext{ is in } C_{\lambda^{-1/2}}, \ \lambda>0 \ 0 & ext{when } y-y(0) ext{ is in } C_0 \cup D \ . \end{cases}$$

Now for $\lambda > 0$, $\lambda^{-1/2}x + \xi - \lambda^{-1/2}x(0) - \xi = \lambda^{-1/2}x$ which is in $C_{\lambda^{-1/2}}$

for all ξ and m_1 -a.e. x. Hence for each $\lambda > 0$, $(\lambda^{-1/2}x + \xi) = g(\lambda)$ for every ξ and m_1 -a.e. x. Hence

$$egin{aligned} & (I_{\lambda}(F)\psi)(\xi) = \int_{\mathcal{C}_0(T)} F(\lambda^{-1/2}x+\xi)\psi(\lambda^{-1/2}x(1)+\xi)dm_1(x) \ &= g(\lambda)\int_{\mathcal{C}_0(T)}\psi(\lambda^{-1/2}x(1)+\xi)dm_1(x) \ &= g(\lambda)(K_\lambda\psi)(\xi) \;. \end{aligned}$$

The result now follows since $||K_{\lambda}|| = 1$ [24, p. 776].

The next result shows that $I_{\lambda}(F)$, $\lambda > 0$, can depend on λ in a rather arbitrarily pathological way.

THEOREM 35. Let $f(\lambda)$ be any function (however nasty) on $(0, \infty)$. Then there exists F such that $I_{\lambda}F = f(\lambda)K_{\lambda}$ for $\lambda > 0$.

Proof. Simply let

$$F(y) = egin{cases} f(\lambda) & ext{when } y-y(0) ext{ is in } C_{\lambda}^{-1/2} ext{, } \lambda > 0 ext{ ,} \ 0 & ext{when } y-y(0) ext{ is in } C_0 \cup D ext{ ,} \end{cases}$$

and argue as in the preceding proof.

We will conclude this section with some relatively brief comments on the papers [14, 15, 17]. Although we are less familiar with these papers, it appears very likely that some of the ideas and results of this paper (as well as the extensions of [16] found in [34] in the case p=2) will allow one to better understand and extend these papers as well. As a test of this, we tried to see if we could improve on Theorem 1 of [15] with reference to the assumptions on ψ . In [15], it is assumed that ψ is Borel measurable and that for every $\lambda > 0$ and every η in $C_0(T)$, $\psi(\lambda y + \eta)$ is Wiener integrable as a function of y. We get a related result under the simpler and more general assumption that ψ is scale-invariant integrable; i.e., ψ is *s*-a.e. defined and $\int_{C_0(T)} |\psi(\lambda x)| dm_1(x) < \infty$ for every $\lambda > 0$.

THEOREM 36. Let F be bounded continuous in the uniform topology on $C_2^*[R]$ and let $\lambda > 0$. Then the operator $I_{\lambda}(F)$ (now defined in terms of a Yeh-Wiener integral) carries the space of scale-invariant integrable functions into itself.

We will not carry out the proof but we mention that the key step involves a use of Proposition 21 and an extension of Theorem 1 of [14] to a situation summarized by the formula

$$egin{aligned} &\int_{C_2[\mathcal{R}] imes C_1[lpha,eta]} f(x(b, \ \cdot), \ y(\ \cdot)) d(m \ imes \ m_1)(x, \ y) \ &= \int_{C_1[lpha,eta] imes C_1[lpha,eta]} f(((b \ - a)/2)^{1/2} z(\ \cdot), \ y(\ \cdot)) d(m_1 \ imes \ m_1)(z, \ y) \end{aligned}$$

where m denotes Yeh-Wiener measure on $C_2[R]$.

In addition to the likelihood of being able to extend [17], it appears that one should be able to give examples and results for the operator-valued Yeh-Feynman integral which parallel the examples and results given earlier in this section for the operator-valued Cameron-Storvick (or Feynman) integral.

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