Weak measurement takes a simple form for cumulants

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A weak measurement on a system is made by coupling a pointer weakly to the system and then measuring the position of the pointer. If the initial wavefunction for the pointer is real, the mean displacement of the pointer is proportional to the so-called weak value of the observable being measured. This gives an intuitively direct way of understanding weak measurement. However, if the initial pointer wavefunction takes complex values, the relationship between pointer displacement and weak value is not quite so simple, as pointed out recently by R. Jozsa [1]. This is even more striking in the case of sequential weak measurements [2]. These are carried out by coupling several pointers at different stages of evolution of the system, and the relationship between the products of the measured pointer positions and the sequential weak values can become extremely complicated for an arbitrary initial pointer wavefunction. Surprisingly, all this complication vanishes when one calculates the cumulants of pointer positions. These are directly proportional to the cumulants of sequential weak values. This suggests that cumulants have a fundamental physical significance for weak measurement.

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I. INTRODUCTION

In physics, formal simplicity is often a reliable guide to the significance of a result. The concept of weak measurement, due to Aharonov and his coworkers [3, 4], derives some of its appeal from the formal simplicity of its basic formulae. One can extend the basic concept to a sequence of weak measurements carried out at a succession of points during the evolution of a system [2], but then the formula relating pointer positions to weak values turns out to be not quite so simple, particularly if one allows arbitrary initial conditions for the measuring system. I show here that the complications largely disappear if one takes the cumulants of expected values of pointer positions; these are related in a formally satisfying way to weak values, and this form is preserved under all measurement conditions.

The goal of weak measurement is to obtain information about a quantum system given both an initial state $|\psi_i\rangle$ and a final, post-selected state $|\psi_f\rangle$. Since weak measurement causes only a small disturbance to the system, the measurement result can reflect both the initial and final states. It can therefore give richer information than a conventional (strong) measurement, including in particular the results of all possible strong measurements [5, 6]. To carry out the measurement, a measuring device is coupled to the system in such a way that the system is only slightly perturbed; this can be achieved by having a small coupling constant $g$. After the interaction, the pointer’s position $q$ is measured (or possibly some other pointer observable; e.g. its momentum $p$). Suppose that, following the standard von Neumann paradigm, [7], the interaction between measuring device and system is taken to be $H_{\text{int}} = g\delta(t)pA$, where $p$ is the momentum of a pointer and the delta function indicates an impulsive interaction at time $t$. It can be shown [4] that the expectation of the pointer position, ignoring terms of order $g^2$ or higher, is

$$\langle q \rangle = gReA_w,$$

where $A_w$ is the weak value of the observable $A$ given by

$$A_w = \langle \psi_f | A | \psi_i \rangle / \langle \psi_f | \psi_i \rangle.$$ (2)

As can be seen, (1) has an appealing simplicity, relating the pointer shift directly to the weak value. However, this formula only holds under the rather special assumption that the initial pointer wavefunction $\phi$ is a gaussian, or, more generally, is real and has zero mean. When $\phi$ is a completely general wavefunction, i.e. is allowed to take complex values and have any mean value [1, 2], equation (1) is replaced by

$$\langle q \rangle = \langle q \rangle_i + gReA_w + gImA_w (\langle pq + qp \rangle_i - 2\langle q \rangle_i \langle p \rangle_i),$$ (3)
where, for any pointer variable $x$, $\langle x \rangle_i$ denotes the initial expected value $\langle \phi | x | \phi \rangle$ of $x$; so for instance $\langle q \rangle_i$ and $\langle p \rangle_i$ are the means of the initial pointer position and momentum, respectively. (Again, this formula ignores terms of order $g^2$ or higher.)

Equation 3 seems to have lost the simplicity of 1, but we can rewrite it as

$$\langle q \rangle = \langle q \rangle_i + g \text{Re}(\xi_A),$$

where

$$\xi = -2i (\langle qp \rangle_i - \langle q \rangle_i \langle p \rangle_i),$$

and equation 4 is then closer to the form of 1. As will become clear, this is part of a general pattern.

One can also weakly measure several observables, $A_1, \ldots, A_n$, in succession 2. Here one couples pointers at several locations and times during the evolution of the system, taking the coupling constant $g_k$ at site $k$ to be small. One then measures each pointer, and takes the product of the positions $q_k$ of the pointers. For two observables, and in the special case where the initial pointer distributions are real and have zero mean, e.g. a gaussian, one finds 2

$$\langle q_1 q_2 \rangle = \frac{g_1 g_2}{2} \text{Re} \left[ (A_2, A_1)_w + (A_1)_w (A_2)_{w} \right],$$

ignoring terms in higher powers of $g_1$ and $g_2$. Here $(A_2, A_1)_w$ is the sequential weak value defined by

$$(A_2, A_1)_w = \frac{\langle \psi_f W_A^2 V_A U \psi_i \rangle}{\langle \psi_f | W V U | \psi_i \rangle},$$

where $U$ is a unitary taking the system from the initial state $| \psi_i \rangle$ to the first weak measurement, $V$ describes the evolution between the two measurements, and $W$ takes the system to the final state. (Note the reverse order of operators in $(A_2, A_1)$, which reflects the order in which they are applied.) If we drop the assumption about the special initial form of the pointer distribution and allow an arbitrary $\phi$, then the counterpart of 6 becomes extremely complicated: see Appendix, equation A1.

Even the comparatively simple formula 6 is not quite ideal. By analogy with 1, we would hope for a formula of the form $\langle q_1 q_2 \rangle \propto \text{Re}(A_2, A_1)_w$, but there is an extra term $(A_1)_w (A_2)_{w}$. What we seek, therefore, is a relationship that has some of the formal simplicity of 1 and furthermore preserves its form for all measurement conditions. It turns out that this is possible if we take the cumulant of the expectations of pointer positions. As we shall see in the next section, this is a certain sum of products of joint expectations of subsets of the $q_i$, which we denote by $\langle q_1 \ldots q_n \rangle^c$. For a set of observables, we can define a formally equivalent expression using sequential weak values, which we denote by $(A_n, \ldots, A_1)_w^c$. Then the claim is that, up to order $n$ in the coupling constants $g_k$ (assumed to be all of the same approximate order of magnitude):

$$\langle q_1 \ldots q_n \rangle^c = g_1 \ldots g_n \text{Re} \{ \xi(A_n, \ldots, A_1)_w^c \},$$

where $\xi$ is a factor dependent on the initial wavefunctions for each pointer. Equation 8 holds for any initial pointer wavefunction, though different wavefunctions produce different values of $\xi$. The remarkable thing is that all the complexity is packed into this one number, rather than exploding into a multiplicity of terms, as in $(A1)$.

Note also that 4 has essentially the same form as 8 since, in the case $n = 1$, $(A)_w^c = A_w$. However, there is an extra term $\langle q \rangle_i$ in 4; this arises because the cumulant for $n = 1$ is anomalous in that its terms do not sum to zero.

II. CUMULANTS

Given a collection of random variables, such as the pointer positions $q_i$, the cumulant $\langle q_1 \ldots q_n \rangle^c$ is a polynomial in the expectations of subsets of these variables 3, 9; it has the property that it vanishes whenever the set of variables $q_i$ can be divided into two independent subsets. One can say that the cumulant, in a certain sense, picks out the maximal correlation involving all of the variables.

We introduce some notation to define the cumulant. Let $x$ be a subset of the integers $\{1, \ldots, n\}$. We write $\prod_{x(i)} q$ for $\prod_{i=1}^{|x|} q_{x(i)}$, where $|x|$ is the size of $x$ and the indices of the $q$’s in the product run over all the integers $x(i)$ in $x$. Then the cumulant is given by

$$\langle q_1 \ldots q_n \rangle^c = \sum_{b = \{b_1, \ldots, b_k\}} a_k \prod_{j=1}^k \left\langle \prod_{b_j} q \right\rangle,$$
where \( b = \{ b_1, \ldots, b_k \} \) runs over all partitions of the integers \( \{1, \ldots, n\} \) and the coefficient \( a_k \) is given by
\[
a_k = (k-1)!(-1)^{k-1}. \tag{10}
\]

For \( n = 1 \) we have \( \langle q \rangle^c = \langle q \rangle \), and for \( n = 2 \)
\[
\langle q_1 q_2 \rangle^c = \langle q_1 q_2 \rangle - \langle q_1 \rangle \langle q_2 \rangle. \tag{11}
\]

There is an inverse operation for the cumulant \( \text{[8,10]} \):

**Proposition II.1.**

\[
\langle q_1 \ldots q_n \rangle = \sum_{b = \{ b_1, \ldots, b_k \}} \prod_{j=1}^k \left( \prod_{b_j} q \right)^c. \tag{12}
\]

**Proof.** To see that this equation holds, we must show that the term \( \prod_{j=1}^k \prod_{b_j} q \) obtained by expanding the right-hand side is zero unless \( b \) is the partition consisting of the single set \( \{1, \ldots, n\} \). Replacing each subset \( b_j \) by the integer \( j \), this is equivalent to \( \sum a_k, \ldots a_k = 0 \), where the sum is over all partitions of \( \{1, \ldots, k\} \) by subsets of sizes \( k_1, \ldots, k_r \) and the \( a_k \)'s are given by \( \text{[11]} \). In this sum we distinguish partitions with distinct integers; e.g. \( \{1,2\}, \{3,4\} \) and \( \{1,3\}, \{2,4\} \). There are \( (k_1 \ldots k_r)(l_1 \ldots l_k)^{-1} \) such distinct partitions with subset sizes \( k_1, \ldots, k_r \), where \( i \) is the number of \( k \)'s equal to \( i \), so our sum may be rewritten as \( k! \sum (-1)^{k-1} \cdots (-1)^{k-1} (l_1 ! l_2 ! k_1 ! \ldots k_r !)^{-1} \), where the sum is now over partitions in the standard sense \( \text{[11]} \). This is \( k! \) times the coefficient of \( x^k \) in
\[
\left(1 + x + \frac{x^2}{2!} + \ldots\right) \left(1 + (-x^2/2)^2 + \ldots\right) \left(1 + (x^3/3)^2 + \ldots\right) \ldots \tag{13}
\]
\[
= e^{x-x^2/2+x^3/3\ldots} = e^{\log_e(1+x)} = 1 + x. \tag{14}
\]

Thus the sum is zero except for \( k = 1 \), which corresponds to the single-set partition \( b \). \( \square \)

**Definition II.2.** If \( \{1, \ldots, n\} \) can be written as the disjoint union of two subsets \( S_1 \) and \( S_2 \), we say the variables corresponding to these subsets are independent if
\[
\langle \prod_{S_1} q \prod_{S_2} q \rangle = \langle \prod_{S_1} q \rangle \langle \prod_{S_2} q \rangle, \tag{15}
\]
for any subsets \( S'_1 \subseteq S_1 \).

We now prove the characteristic property of cumulants:

**Proposition II.3.** The cumulant vanishes if its arguments can be divided into two independent subsets.

**Proof.** For \( n = 2 \) this follows at once from \( \text{[11]} \) and \( \text{[15]} \), and we continue by induction. From \( \text{[12]} \) and the inductive assumption for \( n - 1 \), we have
\[
\langle q_1 \ldots q_n \rangle = \langle q_1 \ldots q_n \rangle^c + \sum_{b = \{ b_1, \ldots, b_k \} \subset S_1} \prod_{j=1}^k \left( \prod_{b_j} q \right)^c \sum_{c = \{ c_1, \ldots, c_l \} \subset S_2} \prod_{j=1}^l \left( \prod_{c_j} q \right)^c. \tag{16}
\]

This holds because any term on the right-hand side of \( \text{[12]} \) vanishes when any subset of the partition \( b \) includes elements of both \( S_1 \) and \( S_2 \). Using \( \text{[12]} \) again, this implies
\[
\langle q_1 \ldots q_n \rangle = \langle q_1 \ldots q_n \rangle^c + \prod_{S_1} q \prod_{S_2} q, \tag{17}
\]
and by independence, \( \langle q_1 \ldots q_n \rangle^c = 0 \). Thus the inductive assumption holds for \( n \). \( \square \)

In fact, the coefficients \( a_k \) in \( \text{[9]} \) are uniquely determined to have the form \( \text{[10]} \) by the requirement that the cumulant vanishes when the variables form two independent subsets \( \text{[12,13]} \).
For $n = 2$, the cumulant is just the covariance, $\langle q_1 q_2 \rangle^c = \langle (q_1 - \langle q_1 \rangle)(q_2 - \langle q_2 \rangle) \rangle$, and the same is true for $n = 3$, namely $\langle q_1 q_2 q_3 \rangle^c = \langle (q_1 - \langle q_1 \rangle)(q_2 - \langle q_2 \rangle)(q_3 - \langle q_3 \rangle) \rangle$. For $n = 4$, however, there is a surprise. The covariance is given by

$$\sum_{i=1}^{4} \langle q_i - \langle q_i \rangle \rangle = \langle q_1 q_2 q_3 q_4 \rangle - \sum \langle q_i q_j q_k q_l \rangle \langle q_i \rangle + \sum \langle q_i q_j \rangle \langle q_k \rangle \langle q_l \rangle - 3 \langle q_1 \rangle \langle q_2 \rangle \langle q_3 \rangle \langle q_4 \rangle,$$  \hspace{1cm} (18)

where the sums include all distinct combinations of indices, but the cumulant is

$$\langle q_1 q_2 q_3 q_4 \rangle^c = \langle q_1 q_2 q_3 q_4 \rangle - \sum \langle q_i q_j q_k q_l \rangle \langle q_i \rangle - \sum \langle q_i q_j \rangle \langle q_k \rangle \langle q_l \rangle + 2 \sum \langle q_i q_j \rangle \langle q_k \rangle \langle q_l \rangle - 6 \langle q_1 \rangle \langle q_2 \rangle \langle q_3 \rangle \langle q_4 \rangle,$$ \hspace{1cm} (19)

which includes terms like $\langle q_1 q_2 \rangle \langle q_3 q_4 \rangle$ that do not occur in the covariance. Note that, if the subsets $\{1, 2\}$ and $\{3, 4\}$ are independent, the covariance does not vanish, since independence implies we can write the first term in (18) as $\langle q_1 q_2 q_3 q_4 \rangle = \langle q_1 q_2 \rangle \langle q_3 q_4 \rangle$ and there is no cancelling term. However, as we have seen, the cumulant does contain such a term, and it is a pleasant exercise to check that the whole cumulant vanishes.

### III. SEQUENTIAL WEAK VALUES AND CUMULANTS

To carry out a sequential weak measurement, one starts a system in an initial state $|\psi_i\rangle$, then weakly couples pointers at several times $t_k$ during the evolution of the system, and finally post-selects the system state $|\psi_f\rangle$. One then measures the pointers and finally takes the product of the values obtained from these pointer measurements. It is assumed that one can repeat the whole process many times to obtain the expectation of the product of pointer values. If one measures pointer positions $q_k$, for instance, one can estimate $\langle q_1 \ldots q_n \rangle$, but one could also measure the momenta of the pointers to estimate $\langle p_1 \ldots p_n \rangle$.

If the coupling for the $k$th pointer is given by $H_{int} = \delta(t - t_k) r_k p$, and if the individual initial pointer wavefunctions are gaussian, or, more generally, are real with zero mean, then it turns out [2] that these expectations can be expressed in terms of sequential weak values of order $n$ or less. Here the sequential weak value of order $n$, $(A_n, \ldots A_1)_w$, is defined by

$$A_n, \ldots A_1) = \frac{|\psi_f,U_{n+1}A_nU_n \ldots A_1U_1|\psi_i\rangle}{\langle \psi_f,U_{n+1} \ldots U_1|\psi_i\rangle},$$ \hspace{1cm} (20)

where $U_i$ defines the evolution of the system between the measurements of $A_{i-1}$ and $A_i$.

When the $A_k$ are projectors, $A_k = |x_k\rangle\langle x_k|$, we can write the sequential weak value as [2]

$$A_n, \ldots A_1 = \frac{\langle \psi_f,U_{n+1}|x_n\rangle \langle x_n,U_n|x_{n-1}\rangle \ldots \langle x_1,U_1|\psi_i\rangle}{\sum_{y} \langle \psi_f,U_{n+1}|y_n\rangle \langle y_n,U_n|y_{n-1}\rangle \ldots \langle y_1,U_1|\psi_i\rangle} = \frac{\text{amplitude}(x)}{\sum_{y} \text{amplitude}(y)},$$ \hspace{1cm} (21)

which shows that, in this case, the weak values have a natural interpretation as the amplitude for following the path defined by the $x_k$. Figure 1 shows an example taken from [2] where the path (labelled by '1' and '2' successively) is a route taken by a photon through a pair of interferometers, starting by injecting the photon at the top left (with state $|\psi_i\rangle$) and ending with post-selection by detection at the bottom right (with final state $|\psi_f\rangle$).

![Diagram of interferometers](image.png)

**FIG. 1:**

In the last section, the cumulant was defined for expectations of products of variables. One can define the cumulant for other entities by formal analogy; for instance for density matrices [10], or hypergraphs [9]. We can do the same.
for sequential weak values, defining the cumulant by (20) with $\langle \prod_{b_j} \rangle_q$ replaced by $\langle \hat{A}_{b_1(1)} \cdots \hat{A}_{b_j(1)} \rangle_w$, where the arrow indicates that the indices, which run over the subset $b_j$, are arranged in ascending order from right to left. For example, for $n = 1$, $(A_w)^c = A_w$, and for $n = 4$

$$(A_4, A_3, A_2, A_1)^c_w = (A_4, A_3, A_2, A_1)_w - \sum (A_i, A_j, A_k)_{w}(A_i)_{w} - \sum (A_i, A_j)_{w}(A_k, A_l)_{w}$$

$$+ 2 \sum (A_i, A_j)_{w}(A_k)_{w}(A_l)_{w} - 6(A_1)_{w}(A_2)_{w}(A_3)_{w}(A_4)_{w}.$$  

There is a notion of independence that parallels (15): given a disjoint partition $S_1 \cup S_2 = \{1, \ldots, n\}$ such that

$$(\hat{A}_{S_1^c \cup S_2^c})_w = (\hat{A}_{S_1^c})_w(\hat{A}_{S_2^c})_w,$$

for any subsets $S_1' \subseteq S_1$, then we say the observables labelled by the two subsets are weakly independent. There is then an analogue of Lemma 1.3

**Lemma III.1.** The cumulant $(A_n, \ldots, A_1)^c_w$ vanishes if the $A_k$ are weakly independent for some subsets $S_1$, $S_2$.

As an example of this, if one is given a bipartite system $\mathcal{H}^A \otimes \mathcal{H}^B$, and initial and final states that factorise as $|\psi_i\rangle = |\psi_i^A\rangle \otimes |\psi_i^B\rangle$ and $|\psi_f\rangle = |\psi_f^A\rangle \otimes |\psi_f^B\rangle$, then observables on the $A$- and $B$-parts of the system are clearly weakly independent. Another class of examples comes from what one might describe as a “bottleneck” construction, where, at some point the evolution of the system is divided into two parts by a one-dimensional projector (the bottleneck) and its complement, and the post-selection excludes the complementary part. Then, if all the measurements before the projector belong to $S_1$ and all those after the projector belong to $S_2$, the two sets are weakly independent. This follows because we can write

$$(\hat{A}_{S_1^c \cup S_2^c})_w = (\hat{A}_{S_1^c})_w(\hat{A}_{S_2^c})_w,$$

where $W_k|\psi_b\rangle\langle \psi_b|V_k$ is the part of $U_k$ lying in the post-selected subspace. As an illustration of this, suppose we add a connecting link (Figure 2, “L”) between the two interferometers in Figure 1 so $|\psi_b\rangle\langle \psi_b|$, the bottleneck, is the projection onto $L$, and post-selection discards the part of the wavefunction corresponding to the path $L'$. Then measurements at ‘1’ and ‘2’ are weakly independent; in fact $(A_1)_w = 1/2$, $(A_2)_w = 1/2$ and $(A_2, A_1)_w = 1/4$. Note that the same measurements are not independent in the double interferometer of Figure 1 where $(A_1)_w = 0$, $(A_2)_w = 0$, and yet, surprisingly, $(A_2, A_1)_w = -1/2$, [2].

**IV. THE MAIN THEOREM**

Consider $n$ system observables $A_1, \ldots, A_n$. Suppose $s_k$, for $k = 1, \ldots, n$, are observables of the $k$th pointer, namely Hermitian functions $s_k(q_k, p_k)$ of pointer position $q_k$ and momentum $p_k$, and the interaction Hamiltonian for the weak
measurement of system observable $A_k$ is $H_k = g_k s_k A_k$, where $g_k$ is a small coupling constant (all $g_k$ being assumed of the same order of magnitude $g$). Suppose further that the pointer observables $r_k$ are measured after the coupling. Let $\phi_k$ be the $k$-th pointer’s initial wave-function. For any variable $x_k$ associated to the $k$-th pointer, write $\langle x_k \rangle_i$ for $\langle \phi_k | x_k | \phi_k \rangle$.

We are now almost ready to state the main theorem, but first need to clarify the measurement procedure. When we evaluate expectations of products of the $r_k$ for different sets of pointers, for instance when we evaluate $\langle r_1 r_2 \rangle$, we have a choice. We could either couple the entire set of $n$ pointers and then select the data for pointers 1 and 2 to give $\langle r_1 r_2 \rangle$. Or we could carry out an experiment in which we couple just pointers 1 and 2 to give $\langle r_1 r_2 \rangle$. These procedures give different answers. For instance, if we couple three pointers and measure pointers 1 and 2 to get $\langle r_1 r_2 \rangle$, in addition to the terms in $g_1$, $g_2$ and $g_1 g_2$ we also get terms in $g_2 g_3$ and $g_1 g_3$ involving the observable $A_3$. This means we get a different cumulant $\langle r_1 \ldots r_n \rangle^c$, depending on the procedure used. In what follows, we regard each expectation as being evaluated in a separate experiment, with only the relevant pointers coupled. It will be shown elsewhere that, with the alternative definition, the theorem still holds but with a different value of the constant $\xi$.

**Theorem IV.1 (Cumulant theorem).** For $n \geq 2$, for any pointer observables $r_k$ and $s_k$, and for any initial pointer wavefunctions $\phi_k$, up to total order $n$ in the $g_k$,

$$
\langle r_1 \ldots r_n \rangle^c = g_1 \ldots g_n Re \{ \xi \langle A_1 \ldots , A_n \rangle \wedge \},
$$

where $\xi$ (sometimes written more explicitly as $\xi_{r_1 \ldots r_n}$) is given by

$$
\xi = 2(-i)^n \left( \prod_{k=1}^n \langle r_k s_k \rangle - \prod_{k=1}^n \langle r_k \rangle \langle s_k \rangle \right).
$$

For $n = 1$ the same result holds, but with the extra term $\langle r \rangle_i$:

$$
\langle r \rangle = \langle r \rangle_i + g Re(\xi A_w).
$$

**Proof.** We use the methods of $[2]$ to calculate the expectations of products of pointer variables for sequential weak measurements. Let the initial and final states of the system be $|\psi_i\rangle$ and $|\psi_f\rangle$, respectively. Consider some subset $b = \{b_1, \ldots, b_\nu\}$ of $\{1, \ldots, n\}$, with $b_1 \leq b_2 \leq \ldots \leq b_\nu$. The state of the system and the pointers $b_1, \ldots, b_\nu$ after the coupling of those pointers is

$$
\Psi_{S,M} = U_{n+1} \ldots U_{b_\nu+1} e^{-ig_{b_\nu} s_{b_\nu} A_{b_\nu} U_{b_\nu-1} \ldots e^{-ig_{b_1} s_{b_1} A_{b_1}} U_{b_1} \ldots U_1 |\psi_i\rangle \phi_{b_1}(r_{b_1}) \ldots \phi_{b_\nu}(r_{b_\nu}),
$$

and following post-selection by the system state $|\psi_f\rangle$, the state of the pointers is

$$
\Psi_{M} = |\psi_f\rangle U_{n+1} \ldots U_{b_\nu+1} e^{-ig_{b_\nu} s_{b_\nu} A_{b_\nu} U_{b_\nu-1} \ldots e^{-ig_{b_1} s_{b_1} A_{b_1}} U_{b_1} \ldots U_1 |\psi_i\rangle \phi_{b_1}(r_{b_1}) \ldots \phi_{b_\nu}(r_{b_\nu}).
$$

Expanding each exponential, we have

$$
\langle r_{b_1} \ldots r_{b_\nu} \rangle = \frac{\int \Psi_{M}^{\dagger} r_{b_1} \ldots r_{b_\nu} \Psi_{M} dr_{b_1} \ldots dr_{b_\nu}}{\int |\Psi_{M}|^2 dr_{b_1} \ldots dr_{b_\nu}},
$$

where $i_k \geq 0$ are integers, $i_1, \ldots, i_n \in b$ means that $i_l = 0$ for $l \notin b$, and

$$
\alpha_{i_1 \ldots i_n} = \left( \prod_{k=1}^n g_k^{i_k} \right) (A^n_{i_1} \ldots A^1_{i_n}) \wedge,
$$

$$
u_{i_1 \ldots i_n} = \int (m!)^{-1} (-is_k)^m \phi_k(r_k) r_k(!)^{-1} (-is_k)^l \phi_k(r_k) dr_k,
$$

$$
u_{i_1 \ldots i_n} = \int (m!)^{-1} (-is_k)^m \phi_k(r_k) (!)^{-1} (-is_k)^l \phi_k(r_k) dr_k.
$$

Let us write (30) as

$$
\langle r_{b_1} \ldots r_{b_\nu} \rangle = \frac{\sum_{i \in b, j \in b} x_{ij}}{\sum_{i \in b, j \in b} y_{ij}}.
$$
where
\[ x_{ij} = \alpha_{i_1 \ldots i_n} \overline{a}_{j_1 \ldots j_n} b_{i_{k_1} j_{k_1}} \ldots b_{i_{k_n} j_{k_n}}, \]
\[ y_{k,i} = \alpha_{k_1 \ldots k_n} \overline{a}_{i_1 \ldots i_n} b_{k_{i_1} j_{i_1}} \ldots b_{k_{i_n} j_{i_n}}, \]
and \( i \) denotes the index set \( \{i_1 \ldots i_n\} \), etc. Define
\[ X_b = \sum_{i \in b, j \in b} x_{ij}, \quad Y_b = \sum_{i \in b, j \in b} y_{k,i}. \]
Then
\[ \langle r_1, \ldots, r_n \rangle^c = \sum_{b_1, \ldots, b_k} (k-1)! \sum_{l=1}^{k} \prod_{i=1}^{k} \langle r_{b_i(1)}, \ldots, r_{b_i(|b_i|)} \rangle \]
\[ = \sum_{b_1, \ldots, b_k} (k-1)! \sum_{l=1}^{k} \prod_{l=1}^{k} \frac{X_{b_l}}{Y_{b_l}}. \]
Set \( \mathcal{Y} = \prod_{b \subseteq \{1, \ldots, n\}} Y_b \), where \( b \) in the product ranges over all distinct subsets of the integers \( \{1, \ldots, n\} \). Then \( \mathcal{Y}(r_1 \ldots r_n)^c \) is an (infinite) weighted sum of terms
\[ z_\mathcal{I} = (x_{i(1)j(1)} \ldots x_{i(m)j(m)}) (y_{k(1)i(1)} \ldots y_{k(m')i(m')}), \]
where
\[ \mathcal{I} = \mathcal{I}_1 \cup \mathcal{I}_2 \cup \mathcal{I}_3 \]
\[ = \{i(1), \ldots, i(m)\} \cup \{j(1), \ldots, j(m)\} \cup \{k(1), \ldots, k(m')\} \cup \{l(1), \ldots, l(m')\} \]
denotes the set of all the index sets that occur in \( z_\mathcal{I} \). The strategy is to show that, when the size of the index set \( \mathcal{I} \) is less than \( n \), the coefficient of \( z_\mathcal{I} \) vanishes; by \((31)\) this implies that all coefficients of order less than \( n \) in \( g \) vanish. We then look at the index sets of size \( n \), corresponding to terms of order \( g^n \), and show that the relevant terms sum up to the right-hand side of \((24)\). But if \( \mathcal{Y}(r_1 \ldots r_n)^c = g^nx + O(g^{n+1}) \) for some \( x \), then we also have \( \langle r_1 \ldots r_n \rangle^c = g^nx + O(g^{n+1}) \), since \( \mathcal{Y} = 1 + O(g) \).
Let \( b = \{b_1, \ldots, b_s\} \) be a partition of \( \{1, \ldots, n\} \). We say that \( b \) is a valid partition for \( \mathcal{I} \) if
(i) For each \( r \) with \( 1 \leq r \leq m \), \( i(r) + j(r) \in b_l \), for some \( b_l \), and we can associate a distinct \( b_l \) to each \( r \). (Here \( i+j \) means the index set \( \{i_1 + j_1, \ldots, i_n + j_n\} \).)
(ii) For each \( r \) with \( 1 \leq r \leq m' \), \( k(r) + l(r) \in S \), for some subset \( S \subseteq \{1, \ldots, n\} \) that is not in the partition \( b \), i.e. for which \( S \neq b_l \) for any \( l \), and we can associate a distinct \( S \) to each \( r \). Let \( \gamma(\mathcal{I}, b) \) be the number of ways of associating a subset \( S \) to each \( r \).

**Lemma IV.2.** The coefficient of \( z_\mathcal{I} \) in \( \mathcal{Y}(r_1 \ldots r_n)^c \) is zero if all the index sets in \( \mathcal{I} \) have a zero at some position \( r \).

*Proof.* If we expand \( \mathcal{Y}(r_1 \ldots r_n)^c \) using \((39)\), each term in this expansion is associated with a partition \( b \) of \( \{1, \ldots, n\} \). Let \( b \) be a valid partition for \( \mathcal{I} \), and let \( c = \{c_1, \ldots, c_s\} \) denote the partition derived from \( b \) by removing \( r \) from the subset \( b_l \) that contains it, and deleting that subset if it contains only \( r \). Then the following partitions include \( b \) and are all valid:
\[ c^{(1)} = \{(rc_1), c_2, \ldots, c_s\} \]
\[ c^{(2)} = \{c_1, (rc_2), \ldots, c_s\} \]
\[ \ldots \ldots \]
\[ c^{(s)} = \{c_1, c_2, \ldots, (rc_s)\} \]
\[ c^{(s+1)} = \{r, c_1, c_2, \ldots, c_s\}. \]
Each partition \( c^{(i)} \), for \( 1 \leq i \leq s + 1 \) contributes \( \gamma(\mathcal{I}, b) \) to the coefficient of \( z_\mathcal{I} \) in \( \mathcal{Y} \prod_{l=1}^{k} X_{c^{(i)}} / Y_{c^{(i)}} \), and since this term has coefficient \( (s-1)!(-1)^{(s-1)} \) in \((39)\) for partitions \( c^{(1)}, c^{(2)}, \ldots, c^{(s)} \), and \( s!(-1)^s \) for \( c^{(s+1)} \), the sum of all contributions is zero. \( \Box \)
From equations (31) and (41), the power of $g$ in the term $z_T$ is $|I| = |I_1| + |I_2| + |I_3| + |I_4|$. This, together with the preceding Lemma, implies that the lowest order non-vanishing terms in $\mathcal{V}(r_1 \ldots r_n)^c$ are $z_T$ that have a '1' occurring once and once only in each position; we call these complete lowest-degree terms.

**THEOREM IV.3.** The coefficient of a complete lowest-degree term $z_T$ in $\mathcal{V}(r_1 \ldots r_n)^c$ is zero unless only one of the four classes of indices in $\mathcal{I}$, viz. $I_1, I_2, I_3, I_4$, has non-zero terms.

**Proof.** Consider the first case where the indices in $I_2$ and $I_3$ are zero, and where both $I_1$ and $I_4$ have some non-zero indices. Let $b = \{b_1, \ldots, b_r\}$ be the partition whose subsets consist of the non-zero positions in index sets $I_3$ in $I_1$, and let $c = \{c_1, \ldots, c_s\}$ be some partition of the remaining integers in $\{1, \ldots, n\}$. Suppose $s \leq r$. Then we can construct a set of partitions by mixing $b$ and $c$; these have the form

$$d^{(w)} = \{c_1, \ldots, c_{i1}, (x_1 b_1), \ldots, (x_r b_r)\},$$

(43)

where each $x_i$ is either empty or consists of some $c_i$, and all the subsets $b_i$ are present once only in the partition. If any $d^{(w)}$ is eligible, all the other mixtures will also be eligible. Furthermore, the set of all eligible partitions can be decomposed into non-overlapping subsets of mixtures obtained in this way.

Any mixture $d^{(w)}$ gives the same value of $\gamma(I, d^{(w)})$, which we denote simply by $\gamma$; so to show that all the contributions to the coefficient of $z_T$ cancel, we have only to sum over all the mixtures, weighting a partition with $t$ subsets by $(t - 1)!(-1)^{t-1}$. This gives

$$\text{Coefficient of } z_T = \gamma \sum_{i=0}^{s} (s + r - 1)!(-1)^{s+r-i} \binom{s}{i} \binom{r}{i}!$$

$$= \gamma(-1)^{s+r-1}s! \sum_{i=0}^{s} (s + r - i - 1) \ldots (s - i + 1) \binom{r}{i} (-1)^i,$$

$$= \gamma(-1)^{s+r-1}s! \frac{\partial^{-1}}{\partial x^{-1}} \{x^{s-1}(x-1)^r\} |_{x=1} = 0.$$

The above argument applies equally well to the situation where $I_1$ and $I_4$ both have some non-zero indices and indices in $I_2$ and $I_3$ are zero. If the non-zero indices are present in $I_1$ and $I_3$, we can take any eligible partition $a = \{a_1, \ldots, a_t\}$ and divide each subset $a_k$ into two subsets $b_k$ and $c_k$ with the indices from $I_3$ in $b_k$ and those from $I_2$ in $c_k$. All the mixtures of type (43) are eligible, and they include the original partition $a$. By the above argument, the coefficients of $z_T$ arising from them sum to zero. Other combinations of indices are dealt with similarly.

Note that, for $n = 4$ and for the index sets $(1, 1, 0, 0) \in I_1$ and $(0, 0, 1, 1) \in I_2$, the “mixture” argument shows that coefficient of $z_T$ coming from $(r_1 r_2 r_3 r_4)$ cancels that coming from $(r_1 r_2) (r_3 r_4)$ to give zero. This cancellation occurs with the cumulant (19), but not with the covariance (15), where the term $(r_1 r_2) (r_3 r_4)$ is absent.

The only terms that need to be considered, therefore, are complete lowest-degree terms with non-zero indices only in one of the sets $I_1, I_2, I_3, I_4$ and $I_t$. It is easy to calculate the coefficients one gets for such terms. Consider the case of $I_t$. We only need to consider the single partition $b$ whose subsets are the index sets of $I_t$. For this partition, by (40), (45) and (39),

$$z_T = \prod_{e=1}^{t} \alpha_i(e) \prod_{k=1}^{n} u_1^{k} u_0^{k} = g_1 \ldots g_n \prod_{e=1}^{t} \{A_i(e)(|i\rangle\langle i|)\} \sum_{k=1}^{n} \langle r_k s_k \rangle_i$$

(44)

From (39), $z_T$ appears in $\mathcal{V}(r_1 \ldots r_n)^c$ with a coefficient $(t - 1)!(-1)^{t-1}$. So, summing over all $z_T$ with indices in $I_t$, one obtains $g_1 \ldots g_n \{A_i, \ldots, A_1\} \sum_{k=1}^{n} \langle -i\langle r_k s_k \rangle \rangle$. Similarly, from (51), (52) and (39), summing over the $z_T$ with indices in $I_3$ gives the complex conjugate of $g_1 \ldots g_n \{A_n, \ldots, A_1\} \sum_{k=1}^{n} \langle -i\langle r_k s_k \rangle \rangle$. Thus $I_t$ and $I_3$ together give

$$g_1 \ldots g_n (2 \prod_{k=1}^{n} \langle -i\langle r_k s_k \rangle \rangle) \Re \{\langle A_n, \ldots, A_1\rangle \}.$$
and from the above argument, this appears appears in $\mathcal{Y}(r_1 \ldots r_n)^c$ with coefficient $-(t-1)!(-1)^{t-1}$. Again, the index sets in $I_k$ give the complex conjugate of those in $I_k$. Thus we obtain the remaining half of $\xi$, which proves (24) for $n \geq 2$. For $n = 1$ the constant terms (of order zero in $g$) in $\mathcal{Y}(r)$ do not vanish, but the proof goes through if we consider $\mathcal{Y}(\langle r \rangle - \langle r \rangle_i)$ instead.

\[ \square \]

V. EXPLORING THE THEOREM

Consider first the simplest case, where $n = 1$ and $r = q$. We take $H_{int} = g\delta(t)pA$ throughout this section, so $s = p$. Then (26) and (25) give

\[ \langle q \rangle = \langle q \rangle_i + gRe(\xi_q A_w) \quad \text{with} \quad \xi_q = -2i(\langle qp \rangle_i - \langle q \rangle_i \langle p \rangle_i), \quad (46) \]

which we have already seen as equations (4) and (5). If we measure the pointer momentum, so $r = p$, we find

\[ \langle p \rangle = \langle p \rangle_i + gRe(\xi_p A_w) \quad \text{with} \quad \xi_p = -2i(\langle p^2 \rangle_i - \langle p \rangle_i^2), \quad (47) \]

which is equivalent to the result obtained in [1].

For two variables, our theorem for $r_1 = q_1, r_2 = q_2$, is

\[ \langle q_1 q_2 \rangle^c = g_1g_2Re(\xi_{qp}(A_2, A_1)_w^c), \quad (48) \]

with

\[ \xi_{qp} = 2(\langle q_1 \rangle_i \langle p_1 \rangle_i \langle q_2 \rangle_i \langle p_2 \rangle_i - \langle q_1 \rangle_i \langle q_2 \rangle_i \langle p_1 \rangle_i \langle p_2 \rangle_i). \quad (49) \]

The calculations in the Appendix allow one to check (48) and (49) by explicit evaluation; see (A3). Note in passing that, if one writes $\Delta q = \sqrt{(\langle q_1 - \langle q \rangle \rangle^2)$, the Cauchy-Schwarz inequality

\[ \{ \langle q_1 q_2 \rangle^c \} \leq (\langle q_1 - \langle q \rangle \rangle^2) \langle q_2 - \langle q \rangle \rangle^2 \]

implies a Heisenberg-type inequality

\[ \Delta q_1 \Delta q_2 \geq g_1g_2Re\{\xi_{qp}(A_2, A_1)_w^c \}, \]

relating the pointer noise distributions of two weak measurements carried out at different times during the evolution of the system.

When one or both of the $q_k$ in (48) is replaced by the pointer momentum $p_k$, we get

\[ \langle q_1 p_2 \rangle^c = g_1g_2Re(\xi_{qp}(A_2, A_1)_w^c), \quad (50) \]

\[ \langle p_1 p_2 \rangle^c = g_1g_2Re(\xi_{pp}(A_2, A_1)_w^c), \quad (51) \]

with

\[ \xi_{qp} = -2(\langle q_1 \rangle_i \langle p_2 \rangle_i - \langle q_1 \rangle_i \langle p_1 \rangle_i \langle p_2 \rangle_i^2), \quad (52) \]

\[ \xi_{pp} = -2(\langle p_1 \rangle_i \langle p_2 \rangle_i - \langle p_1 \rangle_i^2 \langle p_2 \rangle_i^2). \quad (53) \]

Consider now the special case where $\phi$ is real with zero mean. Then the very complicated expression for $\langle q_1 q_2 \rangle$ in (A1) reduces to

\[ \langle q_1 q_2 \rangle = \frac{g_1g_2}{2} Re \left[ (A_2, A_1)_w + (A_1)_w(\bar{A}_2)_w \right], \quad (54) \]

as shown in [2]. Two further examples from [2] are

\[ \langle q_1 q_2 q_3 \rangle = \frac{g_1g_2g_3}{4} Re \left[ (A_3, A_2, A_1)_w + (A_3, A_2)_w(\bar{A}_1)_w + (A_3, A_1)_w(\bar{A}_2)_w + (A_2, A_1)_w(\bar{A}_3)_w \right], \quad (55) \]

\[ \langle q_1 q_2 q_3 q_4 \rangle = \frac{g_1g_2g_3g_4}{8} Re \left[ (A_4, A_3, A_2, A_1)_w + (A_4, A_3, A_2)_w(\bar{A}_1)_w + \ldots + (A_4, A_3, A_2)_w(\bar{A}_2)_w + \ldots \right]. \quad (56) \]
We can use these formulae to calculate the cumulant \( \langle q_1 \ldots q_n \rangle \), and thus check Theorem [V.4] for this special class of wavefunctions \( \psi \). Each formula contains on the right-hand side a leading sequential weak value, but there are also extra terms, such as \((A_1)_w(A_2)_w\) in (53) and \((A_2, A_1)_w(A_3)_w\) in (55). All these extra terms are eliminated when the cumulant is calculated, and we are left with (24) with \( \xi_{q_1 \ldots q_n} = (1/2)^{n-1} \).

This gratifying simplification depends on the fact that the cumulant is a sum over all partitions. For instance, it does not occur if one uses the covariance instead of the cumulant. To see this, look at the case

\[ \langle (X) \rangle \]

The simultaneous weak value can be carried through unchanged for simultaneous measurement. Thus if we replace the sequential weak value by \( \xi \) the cumulant is calculated, and we are left with (24) with \( \xi_{q_1 \ldots q_n} = (1/2)^{n-1} \).

We have treated the interactions between each pointer and the system individually, the Hamiltonian for the \( k \)th pointer and system being \( H_k = g_k \delta(t - t_k) s_k A_k \), but of course we can equivalently describe the interaction between all the pointers and the system by \( H = \sum_k g_k \delta(t - t_k) s_k A_k \). For sequential measurements we implicitly assume that all the times \( t_k \)'s are distinct. However, the limiting case where there is no evolution between coupling of the pointers and all the \( t_k \)'s are equal is of interest, and is the simultaneous weak measurement considered in [14, 15, 16]. In this case, the state of the pointers after post-selection is given by

\[ \Psi_M = (\psi_f | e^{-i(g_{s_1} A_1 + g_{s_2} A_2 + \ldots + g_{s_n} A_n)} | \psi_i) \phi_1(r_1) \ldots \phi_n(r_n). \]  

(57)

The exponential \( e^{-i(g_{s_1} A_1 + g_{s_2} A_2 + \ldots + g_{s_n} A_n)} \) here differs from the sequential expression \( e^{-i g_{s_1} A_1} \ldots e^{-i g_{s_n} A_n} \) in (28) in that each term in the expansion of the latter appears with the operators in a specific order, viz. the arrow order \( \leftarrow \) as in (22), whereas in the expansion of the former the same term is replaced by a symmetrised sum over all orderings of operators. For instance, for arbitrary operators \( X, Y \) and \( Z \), the third degree terms in \( e^X e^Y e^Z \) include \( X^3/3!, X^2 Y/2! \) and \( X Y Z \), whose counterparts in \( e^{X+Y+Z} \) are, respectively, \( X^3/3!, \{X^2 Y + X Y X + Y X^2\}/3! \) and \( \{X Y Z + X Z Y + Y X Z + Y Z X + Z X Y + Z Y X\}/3! \). Apart from this symmetrisation, the calculations in Section [V] can be carried through unchanged for simultaneous measurement. Thus if we replace the sequential weak value by the simultaneous weak value [14, 15, 16]

\[ (A_{i_1} \ldots A_{i_1})_{\text{ws}} = \frac{1}{k!} \sum_{\pi \in S_k} (A_{i_{\pi(k)}} \ldots A_{i_{\pi(1)}})_{\text{ws}}, \]  

(58)

where the sum on the right-hand side includes all possible orders of applying the operators, we obtain a version of Theorem [V.3] for simultaneous weak measurement:

\[ \langle r_1 \ldots r_n \rangle^c = g_1 \ldots g_n Re \{ \xi(A_n, \ldots, A_1)^c \}_{\text{ws}}. \]  

(59)

Likewise, relations such [54], [55], etc., hold with simultaneous weak values in place of the sequential weak values; indeed, these relations were first proved for simultaneous measurement [14, 15].

1

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FIG. 3:

From (58) we see that, when the operators \( A_k \) all commute, the sequential and simultaneous weak values coincide. One important instance of this arises when the operators \( A_k \) are applied to distinct subsystems, as in the case of the simultaneous weak measurements of the electron and positron in Hardy’s paradox [17, 18].
When the operators do not commute, the meaning of simultaneous weak measurement is not so obvious. One possible physical interpretation follows from the well-known formula

$$e^{X+Y} = \lim_{N \to \infty} \left( e^{X/N} e^{Y/N} \right)^N$$

and its analogues for more operators. Suppose two pointers, one for $A_1$ and one for $A_2$, are coupled alternately in a sequence of $N$ short intervals (Figure 3, top diagram) with coupling strength $g_k/N$ for each interval. This is an enlarged sense of sequential weak measurement [2] in which the same pointer is used repeatedly, coherently preserving its state between couplings. The state after post-selection is

$$\Psi_M = \langle \psi_f | e^{-i\left(\frac{g_2}{N}\right) s_2 A_2} e^{-i\left(\frac{g_1}{N}\right) s_1 A_1} | \psi_i \rangle \phi_1(r_1) \phi_2(r_2).$$

From (60) we deduce that

$$\Psi_M \approx \langle \psi_f | e^{-i(g_2 s_2 A_2 + g_1 s_1 A_1)} | \psi_i \rangle \phi_1(r_1) \phi_2(r_2).$$

This picture readily extends to more operators $A_k$. One can also simulate a simultaneous measurement by averaging the results of a set of sequential measurements with the operators in all orders; in effect, one carries out a set of experiments that implement the averaging in (58). There is then no single act that counts as simultaneous measurement, but weak measurement in any case relies on averaging many repeats of experiments in order to extract the signal from the noise. In a certain sense, therefore, sequential measurement includes and extends the concept of simultaneous measurement. However, if we wish to accomplish simultaneous measurement in a single act, then we need a broader concept of weak measurement where pointers can be re-used; indeed, we can go further, and consider generalised weak coupling between one time-evolving system and another, followed by measurement of the second system. However, even in this case, the measurement results can be expressed algebraically in terms of the sequential weak values of the first system [2].

**VII. LOWERING OPERATORS**

Lundeen and Resch [16] showed that, for a gaussian initial pointer wavefunction, if one defines an operator $a$ by

$$a_{LR} = \langle p^2 \rangle^{1/2} \left( q + \frac{ip}{2\langle p^2 \rangle} \right),$$

then the relationship

$$\langle a_{LR} \rangle = g \langle p^2 \rangle^{1/2} A_w$$

holds. They argued that $a_{LR}$ can be interpreted physically as a lowering operator, carrying the pointer from its first excited state $|1\rangle$, in number state notation, to the gaussian state $|0\rangle$ (despite the fact that the pointer is not actually in a harmonic potential). Although $a_{LR}$ is not an observable, $\langle a_{LR} \rangle$ can be regarded as a prescription for combining expectations of pointer position and momentum to get the weak value.

If instead of $a_{LR}$ one takes

$$a = q + \frac{ip}{2\langle p^2 \rangle},$$

then the even simpler relationship

$$\langle a \rangle = g A_w,$$

holds. We refer to $a$ as a generalised lowering operator.

Lundeen and Resch also extended their lowering operator concept to simultaneous weak measurement of several observables $A_k$. Rephrased in terms of our generalised lowering operators $a_k$ defined by (63), their finding [16] can be stated as

$$\langle a_1 \ldots a_n \rangle = g_1 \ldots g_n (A_1 \ldots A_n)_{ws}. $$

(65)
This is of interest for two reasons. First, the entire simultaneous weak value appears on the right-hand side, not just its real part; and second, the “extra terms” in the simultaneous analogues of (54), (55) and (56) have disappeared. The lowering operator seems to relate directly to weak values.

We can generalise these ideas in two ways. First, we extend them from simultaneous to sequential weak measurements. Secondly, instead of assuming the initial pointer wavefunction is a gaussian, we allow it be arbitrary; we do this by defining a generalised lowering operator

$$a = q + i\frac{p}{\eta}, \quad \text{with} \quad \eta = -\overline{\xi_p/\xi_q};$$

(66)

For a gaussian $\phi$, $\eta = 2(p^2)_i$, so the above definition reduces to (63) in this case. In general, however, $\phi$ will not be annihilated by $a$ and is therefore not the number state $|0\rangle$ (this state is a gaussian with complex variance $\eta^{-1}$). Nonetheless, there is an analogue of Theorem IV.1 in which the whole sequential weak value, rather than its real part, appears:

**Theorem VII.1** (Cumulant theorem for lowering operators). For $n > 1$

$$\langle a_1 \ldots a_n \rangle^c = g_1 \ldots g_n \vartheta (A_n \ldots A_1)^c_w,$$

where $\vartheta$ is given by

$$\vartheta = \sum_{(i_1, \ldots, i_n) \in \{0, 1\}^n} \frac{(-1)^i}{2} \xi_{r_1 \ldots r_n} \left( \xi_{r_1 \ldots r_n} \right) \left( \frac{\xi_{p_1} \ldots \xi_{p_n}}{\xi_{q_1} \ldots \xi_{q_n}} \right)^c,$$

(68)

For $n = 1$ the same result holds, but with the extra term $\langle a \rangle$:

$$\langle a \rangle = \langle a \rangle + \vartheta g A_w.$$

(69)

**Proof.** Put $r_0 = q$, $r_1 = p$. Then

$$\langle a_1 \ldots a_n \rangle^c = \langle (q_1 + i p_1/\eta_1) \ldots (q_n + i p_n/\eta_n) \rangle^c,$$

$$= \sum_{(i_1, \ldots, i_n) \in \{0, 1\}^n} \frac{(-1)^i}{2} \xi_{r_1 \ldots r_n} \left( \xi_{r_1 \ldots r_n} \right) \left( \frac{\xi_{p_1} \ldots \xi_{p_n}}{\xi_{q_1} \ldots \xi_{q_n}} \right)^c,$$

where we used Theorem IV.1 to get the last line, and where $\vartheta$ is given by (68) and $\vartheta$ by

$$\vartheta = \sum_{(i_1, \ldots, i_n) \in \{0, 1\}^n} \frac{(-1)^i}{2} \xi_{r_1 \ldots r_n} \left( \xi_{r_1 \ldots r_n} \right) \left( \frac{\xi_{p_1} \ldots \xi_{p_n}}{\xi_{q_1} \ldots \xi_{q_n}} \right)^c,$$

(note the bar over $\xi_{r_1 \ldots r_n}$ that is absent in the definition of $\vartheta$ by (68)).

We want to prove $\vartheta = 0$, and to do this it suffices to prove that the complex conjugate of the numerator is zero, i.e.

$$\vartheta' = \sum_{(i_1, \ldots, i_n) \in \{0, 1\}^n} (-1)^i \xi_{r_1 \ldots r_n} \left( \xi_{r_1 \ldots r_n} \right) = 0.$$

Let $a_k = \langle q_k s_k \rangle$, $b_k = \langle q_k \rangle \langle s_k \rangle$, $c_k = \langle p_k s_k \rangle$, $d_k = \langle p_k \rangle \langle s_k \rangle$. Using the definition of $\xi$ in (26), the above equation can be written

$$\vartheta'/(2^{n+1}(-1)^n) = \prod_{k=1}^{n} \{ a_k (c_k - d_k) - c_k (a_k - b_k) \} - \prod_{k=1}^{n} \{ b_k (c_k - d_k) - d_k (a_k - b_k) \}$$

$$= \prod (b_k c_k - a_k d_k) - \prod (b_k c_k - a_k d_k) = 0.$$

□

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Suppose the interaction Hamiltonian has the standard von Neumann form $H_{int} = g p A$, so $s = p$ in the definition of $\xi$ by equation (24). Then for $n = 1$, since $\xi_p = \xi_p$ and $(qp)_i = \langle p q \rangle_i$, $\vartheta = (-i)(\xi_q - \xi_q) = (-i)(\langle q p \rangle_i - \langle p q \rangle_i) = 1$, so we get the even simpler result

$$\langle a \rangle = \langle a \rangle_i + g A_w.$$  

(70)

This is valid for all initial pointer wavefunctions, and therefore extends Lundeen and Resch’s equation (64). It seems almost too simple: there is no factor corresponding to $\xi$ in equation (66). However, a dependency on the initial pointer wavefunction is of course built into the definition of a through $\eta$.

For $n > 1$ it is no longer true that $\vartheta = 1$, even with the standard interaction Hamiltonian. However, if in addition $\langle p \rangle_i = 0$, then

$$\vartheta = (-i)^n \prod_{k=1}^n (\langle q_k p_k \rangle_i - \langle p_k q_k \rangle_i) = (-i)^n(i)^n = 1.$$  

Thus $\langle a_1 \ldots a_n \rangle^c = g_1 \ldots g_n (A_{n+1}^c)^c$ for all $n$. Applying the inverse operation for the cumulant, given by Proposition II.4 we deduce:

**Corollary VII.2.** If $\langle p \rangle_i = 0$, e.g. if the initial pointer wavefunction $\phi$ is real, then for $n > 1$

$$\langle a_1 \ldots a_n \rangle = g_1 \ldots g_n (A_n \ldots A_1)_w.$$   

(71)

This is the sequential weak value version of the result for simultaneous measurements, (65), but is more general than the gaussian case treated in (10).

We might be tempted to try to repeat the above argument for pointer positions $q_k$ instead of the lowering operators $a_k$ by applying the anti-cumulant to both sides of (24). This fails, however, because of the need to take the real part.

Recall that the cumulant vanishes when its variables belong to two independent sets. The product of the pointer observables is the best we can do to achieve this goal. However, this brings along extra terms, such as $(A_1)_w (A_2)_w$ in (24), which are an artefact of this method of extracting information. From this perspective, the cumulant extracts the information we really want.

In [2], a somewhat idealised measuring device was being considered, where the pointer position distribution is real and has zero mean. When the pointer distribution is allowed to be arbitrary, the expressions for $\langle q_1 \ldots q_n \rangle$ become wildly complicated (see for instance (A1)). Yet the cumulant of these terms condenses into the succinct equation (24) with all the complexity hidden away in the one number $\xi$. Why does the cumulant have this property?

Recall that the cumulant vanishes when its variables belong to two independent sets. The product of the pointer positions $q_1 \ldots q_n$ will include terms that come from products of disjoint subsets of these pointer positions, and the cumulant of these terms will be sent to zero, by Lemma II.3. For instance, with $n = 2$, the pointers are deflected in proportion to their individual weak values, according to (4), and the cumulant subtracts this component leaving only the component that arises from the $O(q^2)$-influence of the weak measurement of $A_1$ on that of $A_2$. The subtraction of this component corresponds to the subtraction of the term $(A_1)_w (A_2)_w$ from (54). In general, the cumulant of pointer positions singles out the maximal correlation involving all the $q_i$, and the theorem tells us that this is directly related to the corresponding “maximal correlation” of sequential weak values, $(A_{n+1} \ldots A_1)_w^c$, which involves all the operators.

In fact, the theorem tells us something stronger: that it does not matter what pointer observable $r(p, q)$ we measure, e.g. position, momentum, or some Hermitian combination of them, and that likewise the coupling of the pointer with the system can be via a Hamiltonian $H_{int} = gs(p, q)A$ with any Hermitian $s(p, q)$. Different choices of $r$ and $s$ lead only to a different multiplicative constant $\xi$ in front of $(A_{n+1} \ldots A_1)_w^c$ in (24). We always extract the same function of sequential weak values, $(A_{n+1} \ldots A_1)_w^c$, from the system. This argues both for the fundamental character of sequential weak values and also for the key role played by their cumulants.
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APPENDIX A: AN EXPLICIT CALCULATION

To calculate $\langle q_1 q_2 \rangle$ for arbitrary pointer wavefunctions $\phi_1$ and $\phi_2$, we use (28) to determine the state of the two pointers after the weak interaction, and then evaluate the expectation using (29), keeping only terms up to order $g^2$.

We define

$$\mu_k = \langle q_k \rangle_i, \quad \nu_k = \langle p_k \rangle_i, \quad \zeta_k = \langle p_k^2 \rangle_i, \quad \rho_k = \langle q_k p_k \rangle_i, \quad \sigma_k = \langle q_k p_k^2 \rangle_i, \quad \tau_k = \langle p_k q_k p_k \rangle_i,$$

Then, expanding the exponential in (28) and substituting $\Psi$ in (29) gives, up to order $g^2$,

$$\langle q_1 q_2 \rangle = \mu_1 \mu_2 - ig_1 \left\{ \left( (A_1)_w - \bar{(A_1)}_w \right) \mu_1 \nu_1 \mu_2 - \bar{(A_1)}_w \bar{\rho}_1 \mu_2 + (A_1)_w \rho_1 \mu_2 \right\}$$

$$+ g_2 \left\{ \left( (A_2)_w - \bar{(A_2)}_w \right) \mu_1 \mu_2 \nu_2 - (A_2)_w \mu_1 \rho_2 + \bar{(A_2)}_w \mu_1 \mu_2 \nu_2 \right\}$$

$$+ g_2 \left\{ \left( (A_1)_w - \bar{(A_1)}_w \right)^2 \mu_1 \nu_1^2 \mu_2 + \bar{(A_1)}_w^2 \frac{\mu_1 \nu_1^2}{2} \mu_2^2 + (A_1)_w^2 \frac{\mu_1 \nu_1^2}{2} \mu_2^2 \right\}$$

$$- g_2 \left\{ \left( (A_2)_w - \bar{(A_2)}_w \right)^2 \mu_1 \mu_2 \nu_2^2 + (A_2)_w^2 \frac{\mu_1 \nu_1^2}{2} \mu_2^2 + (A_2)_w^2 \frac{\mu_1 \nu_1^2}{2} \mu_2^2 \right\}$$

$$+ g_1 g_2 \left\{ (A_1)_w (A_2)_w \rho_1 \bar{\rho}_2 + (A_1)_w (A_2)_w \bar{\rho}_1 \mu_2 - (A_2, A_1)_w \rho_1 \rho_2 - (A_2, A_1)_w \bar{\rho}_1 \bar{\rho}_2 \right\}$$

$$- g_1 g_2 \left\{ 2 (A_1)_w - \bar{(A_1)}_w \right\} \left( (A_2)_w - \bar{(A_2)}_w \right) \mu_1 \nu_1 \mu_2 \nu_2 \right\}$$

$$+ g_1 g_2 \left\{ \left( (A_2)_w + A_1 \mu_2 \right) (A_1)_w \mu_1 \nu_1 \mu_2 \nu_2 \right\}$$

$$+ g_2 \left\{ \left( (A_1)_w - \bar{(A_1)}_w \right) (A_1)_w \nu_1 \rho_1 \mu_2 - (A_1)_w - \bar{(A_1)}_w \right\} \left( (A_1)_w \nu_1 \bar{\rho}_1 \mu_2 \right)$$

Substituting from (A1) and (A2) a radical simplification occurs:

$$\langle q_1 q_2 \rangle_c = g_1 g_2 \left\{ (A_2, A_1)_w - (A_1)_w (A_2)_w \right\} \left( \mu_1 \nu_1 \mu_2 \nu_2 - \rho_1 \rho_2 \right) + \text{complex conjugate.}$$

This, of course, is what Theorem IV.1 tells us.

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