4

RESOLUTION

4.1 TIME-FREQUENCY RESOLUTION

The famous "uncertainty principle" of quantum mechanics resulted from observations that subatomic particles behave like waves with wave frequency proportional to particle momentum. The classical laws of mechanics enable prediction of the future of a mechanical system by extrapolation from presently known position and momentum. But because of the wave nature of matter with momentum proportional to frequency, such prediction requires simultaneous knowledge of both the location and the frequency of a wave. A sinusoidal wave has a perfectly clearly determined frequency, but it is spread over the infinitely long time axis. At the other extreme is a delta function, which is nicely compressed to a point on the time axis but contains a mixture of all frequencies. A mathematical analysis of the uncertainty principle is thus an analysis relating functions to their Fourier transforms.

Such an analysis begins by definitions of time duration and spectral bandwidth. The time duration of a damped exponential function is infinite if by duration you mean the span of nonzero function values. However, for nearly all practical purposes the time span is chosen as the time required for the amplitude to decay to $e^{-1}$ of its original value. For many functions the span is defined by the span between points on the time or frequency axis where the curve (or its envelope) drop to half of the maximum value. The main idea is that the time span $\Delta t$ or the frequency span $\Delta \omega$ should be able to include most of the total energy but need not contain all of it. The precise definition of $\Delta t$ and $\Delta \omega$ is somewhat arbitrary and may be chosen to simplify analysis. The general statement is that for any function the time duration $\Delta t$ and the spectral bandwidth $\Delta \omega$ are related by

$$\Delta \omega \Delta t \geq 2\pi$$

(4-1.1)

Although it is easy to verify (4-1-1) in many special cases, it is not very easy to deduce (4-1-1) as a general principle. This has, however, been done by D. Gabor [Ref. 17]. He chose to define $\Delta t$ and $\Delta \omega$ by second moments.

A similar and perhaps more basic concept than the product of time and frequency spreads is the relationship between spectral bandwidth and rise time of a system response function. The rise time $\Delta t$ of a system response is also defined somewhat arbitrarily, often as the time span between the time of excitation and the time at which the system response is half its ultimate value. In principle, a broad frequency response can result from a rapid decay time as well as from a rapid rise time. Tightness in the inequality (4-1-1) may be associated with situations in which a certain rise time is quickly followed by an equal decay time. Slackness in the inequality (4-1-1) may be associated with increasing inequality between rise time and decay time. Slackness could also result from some combination of rises and falls such as random combinations. Many systems respond very rapidly compared to the rate at which they subsequently decay. Focusing our attention on such systems, we can now seek to derive the inequality (4-1-1) applied to rise time and bandwidth. The first step is to choose a definition for rise time. The

In locating an earthquake or a petroleum drilling site there will be an uncertainty in location, say $(\Delta x, \Delta y, \Delta z)$ caused by measurement errors and the physical size of the target. In measuring a voltage there will be a measuring accuracy $\Delta v$. The frequency of useful seismic waves will have a bandwidth $\Delta \omega$. The time at which an earthquake occurs will have an uncertainty given by the duration of shaking $\Delta t$. A telescope of diameter $\Delta d$ has at best a resolving power measured by a certain angular range $\Delta \theta$. It is often desirable to make measurements in such a way as to reduce the quantities $\Delta x, \Delta y, \Delta z, \Delta t, \Delta \omega, \Delta \alpha, \Delta d$, and $\Delta \theta$ to values as small as possible. These measures of resolution (which are called variances, tolerances, uncertainties, bandwidths durations, spreads, spans, etc.) sometimes interact with one another in such a way that any experimental modification which reduces one must necessarily increase another or some combination of the others.

The purpose of this chapter is to discuss some of the commonly occurring situations where such conflicting interactions occur.

In this chapter we use $\Delta t$ to denote the time duration of a signal. We use $\tau$ to denote the amount of time which passes between sample points. In other chapters, $\Delta t$ is synonymous with $\tau$, the sample interval.
choice is determined not only for clarity and usefulness but also by the need to ensure tractability of the subsequent analysis. I have found a reasonable definition of rise time to be
\[
\Delta t = \int_0^\infty \frac{1}{t} \frac{b(t)^2 \, dt}{t}
\]
(4-1-2)
where \(b(t)\) is the response function under consideration. The numerator is just a normalizing factor. The denominator says we have defined \(\Delta t\) by the first negative moment. For example, if \(b(t)\) is a step function, then the denominator integral diverges, giving the desired \(\Delta t = 0\) rise time. If \(b(t)^2\) grows linearly from zero to \(t_0\) and then vanishes, the rise time \(\Delta t\) is \(t_0/2\), again a reasonable definition.

Although the Z transform method is a great aid in studying situations where divergence (as \(1/t\)) plays a key role, it does have the disadvantage that it destroys the formal identity between the time domain and the frequency domain. Presumably this disadvantage is not fundamental since we can always go to a limiting process in which the discretized time domain tends to a continuum. In order to utilize the analytic simplicity of the Z transform we now consider the dual to the rise-time problem. Instead of a time function whose square vanishes identically at negative time, we now consider a spectrum \(\overline{B(Z)}\), which vanishes at negative frequencies. We measure how fast this spectrum can rise after \(\omega = 0\). We will find this to be related to the time duration \(\Delta \omega\) of the complex time function \(b_t\). More precisely, we will now define the lowest significant frequency component \(\omega = \Delta \omega\) in the spectrum analogously to (4-1-3) to be
\[
\Delta \omega = \int_0^\infty \frac{BB \, dw}{B \, dw}
\]
(4-1-3)
Without loss of generality we can assume that the spectrum has been normalized so that the numerator integral is unity. In other words, the zero lag of the autocorrelation of \(b_t\) is \(+1\). Then
\[
\frac{1}{\Delta \omega} = \int_0^\infty \frac{BB \, dw}{B \, dw}
\]
(4-1-4)
Now we recall the bilinear transform which gives us various Z transform expressions for \((-\omega)^{-1}\). The one we ordinarily use is the integral \((\ldots 0, 0, 0.5, 1, 1, \ldots)\). We could also use \((-\omega)^{-1} = \ldots 0, 0, 0.5, 0, 0, \ldots)\). The pole right on the unit circle at \(Z = 1\) causes some nonuniqueness. Because \(1/\omega\) is an imaginary odd frequency function we will take the desired expansion to be the odd function of time given by
\[
(-\omega)^{-1} = \frac{1}{2} \left(\cdots - Z^{-2} - Z^{-1} + 0 + Z + Z^2 + \cdots\right)
\]
(4-1-5)
Converting (4-1-4) to an integral on the unit circle in Z transform notation we have
\[
\frac{1}{\Delta \omega} = \int_{-\infty}^{\infty} -i \frac{1}{2\pi} \overline{B(Z)} \, d\omega
\]
(4-1-6)
But since this integral selects the coefficient of \(\omega^0\) of its argument we have
\[
\frac{1}{\Delta \omega} = \frac{1}{2} \sum_{r_1} \text{Im} \left(-r_1 \leq \sum_{r_1} \frac{1}{r_1^2} |r_1|\right)
\]
(4-1-7)
where \(r_1\) is the autocorrelation function of \(b_t\). This may be further expressed as
\[
\frac{1}{\Delta \omega} = \sum_{r_1} \text{Re} \left(-r_1 \leq \sum_{r_1} \frac{1}{r_1^2} |r_1|\right)
\]
(4-1-8)
The sum in (4-1-8) is like an integral representing area under the \(|r_1|\) function. Imagine the \(|r_1|\) function replaced by a rectangle function of equal area. This would define a \(\Delta \omega\) for the \(|r_1|\) function. Any autocorrelation function satisfies \(|r_1| < r_0\) and we have normalized \(r_0 = 1\). Thus, we extend the inequality (4-1-8) by
\[
\frac{1}{\Delta \omega} \leq \sum_{r_1} \frac{1}{r_1^2} |r_1| < \Delta \omega\)
(4-1-9)
Finally, we must relate the duration of a time function \(\Delta t\) to the duration of its autocorrelation \(\Delta \omega\). Generally speaking, it is easy to find a long time function which has short autocorrelation. Just take an arbitrary short time function and convolve it by a long and tortuous all-pass filter. The new function is long, but its autocorrelation is short. If a time function has \(n\) nonzero points, then its autocorrelation has only \(2n - 1\) nonzero points. It is obviously impossible to get a long autocorrelation function out of a short time function. It is not even fair to say that the autocorrelation is twice as long as the original time function because the autocorrelation must lie under some tapering function. To construct a time function with as long an autocorrelation as possible, the best thing to do is to concentrate the energy in two lumps, one at each end of the time function. Even from this extreme example, we see that it is not unreasonable to assert that
\[
\text{inserting into (4-1-9) we have the uncertainty relation}
\]
\[
\Delta t \geq \Delta \omega\]
(4-1-10)
The more usual form of the uncertainty principle uses the frequency variable \(f = 2\omega\) \(\rightarrow \omega\) and a different definition of \(\Delta t\), namely time duration rather than rise time. It is
\[
\Delta t \Delta \omega \geq 1 (\Delta t \text{ is duration})
\]
(4-1-12)
The choice of a 2\(\pi\) scaling factor to convert rise time to duration is indicative of the approximate nature of the inequalities.
EXERCISES

1. Consider \( B(Z) = (1 - (Z/Z_0)^a)(1 - Z/Z_0) \) in the limit \( Z_0 \) goes to the unit circle. Sketch the time function and its squared amplitude. Sketch the frequency function and its squared amplitude. Choose \( \Delta f \) and \( \Delta t \).

2. A time series made up of two frequencies may be written as
\[
B(t) = A \cos \omega_0 t + B \sin \omega_1 t + C \cos \omega_2 t + D \sin \omega_3 t
\]

Given \( \omega_0, \omega_1, \omega_2, \omega_3, B_0, B_1, B_2, B_3 \) show how to calculate the amplitude and phase angles of the two sinusoidal components.

Consider the frequency function graphed below.

![Frequency Function Graph]

FIGURE E 4-1-3

Describe the time function in rough terms indicating the times corresponding to \( 1/f_1 \), \( 1/f_2 \), and \( 1/f_3 \). Try to avoid algebraic calculation. Sketch an approximate result.

PROBLEM FOR RESEARCH

Can you find a method of defining \( \Delta \omega \) and \( \Delta t \) of one-sided wavelets in such a way that for minimum-phase wavelets only the uncertainty principle takes on the equality sign?

4.2 TIME-STATISTICAL RESOLUTION

If you flipped a coin 100 times, it is possible that you would get exactly 50 "heads" and 50 "tails." More likely it would be something between 60–40 and 40–60. Typically, how much deviation from 50 would you expect to see? The average (mean) value should be 50, but some other value is almost always obtained from a random sample. The other value is called the sample mean. We would like to know how much difference to expect between the sample mean and the true mean. The average squared difference is called the variance of the sample mean. For a very large sample, the sample mean should be proportionately much closer to the true mean than for a smaller sample. This idea will lead to an uncertainty relation between the probable error in the estimated mean and the size of the sample. Let us be more precise.

The "true value" of the mean could be defined by flipping the coin \( n \) times and conceiving of \( n \) going to infinity. A more convenient definition of "true value" is that the experiment could be conceived of as having been done separately under identical conditions by an infinite number of people (an ensemble). Such an artifice will enable us to define a time-variable mean for coins which change with time.

The utility of the concept of an ensemble is often subjected to serious attack both from the point of view of the theoretical foundations of statistics and from the point of view of experimentalists applying the techniques of statistics. Nonetheless a great body of geophysical literature uses the artifice of assuming the existence of an unobservable ensemble. The advocates of using ensembles (the Gibbsians) have the advantage over their adversaries (the Bayesians) in that their mathematics is more tractable (and more explainable). So, let us begin!

A conceptual average over the ensemble, called an expectation, is denoted by the symbol \( E \). The index for summation over the ensemble is never shown explicitly; every random variable is presumed to have one. Thus, the true mean at time \( t \) may be defined as

\[
m_t = E(x_t)
\]

If the mean does not vary with time, we may write

\[
m_t = E(x_t)
\]

(4-2-1)

Likewise, we may be interested in a property of \( x_t \) called its variance which is a measure of variability about the mean defined by

\[
\sigma_t^2 = E[(x_t - m)^2]
\]

(4-2-3)

The \( x_t \) random numbers could be defined in such a way that \( \sigma \) or \( m \) or both is either time-variable or constant. If both are constant, we have

\[
\sigma^2 = E[(x_t - m)^2]
\]

(4-2-4)

When manipulating algebraic expressions the symbol \( E \) behaves like a summation sign, namely

\[
E = (\lim n \to \infty) \frac{1}{n} \sum_{i=1}^{n} x_i
\]

(4-2-5)

Notice that the summation index is not given, since the sum is over the ensemble, not time.

Now let \( x_t \) be a time series made up from (identically distributed, independently chosen) random numbers in such a way that \( m \) and \( \sigma \) do not depend on time. Suppose we have a sample of \( n \) points of \( x_t \) and are trying to determine the value of \( m \). We could make an estimate \( \hat{m} \) of the mean \( m \) with the formula

\[
\hat{m} = \frac{1}{n} \sum_{i=1}^{n} x_i
\]

(4-2-6)
A somewhat more elaborate method of estimating the mean would be to take a weighted average. Let \( w_i \) define a set of weights normalized so that
\[
\sum w_i = 1 \quad (4-2-7)
\]
With these weights the more elaborate estimate \( \bar{m} \) of the mean is
\[
\bar{m} = \sum w_i x_i \quad (4-2-8)
\]
Actually (4-2-6) is just a special case of (4-2-8) where the weights are \( w_i = 1/n \).

Our objective in this section is to determine how far the estimated mean \( \bar{m} \) is likely to be from the true mean \( m \) for a sample of length \( n \). One possible definition of this excursion \( \Delta m \) is
\[
(\Delta m)^2 = E[(\bar{m} - m)^2] \quad (4-2-9)
\]
\[
= E[(\sum w_i x_i - m)^2] \quad (4-2-10)
\]
Now utilize the fact that \( m = m \sum w_i = \sum w_i m \)
\[
(\Delta m)^2 = E\left[\left(\sum w_i x_i - m\right)^2\right] \quad (4-2-11)
\]
\[
= E\left[\left(\sum w_i x_i - m\right)\left(\sum w_i x_i - m\right)\right] \quad (4-2-12)
\]
\[
= E\left[\sum w_i x_i (x_i - m)\right] \quad (4-2-13)
\]
Now the expectation symbol \( E \) may be regarded as a summation sign and brought inside the sums on \( t \) and \( s \).
\[
(\Delta m)^2 = \sum_{t} \sum_{s} w_t w_s E[(x_t - m)(x_s - m)] \quad (4-2-14)
\]
By the randomness of \( x \), and \( x_s \), the expectation on the right, that is, the sum over the ensemble, gives zero unless \( s = t \). If \( s = t \), then the expectation is the variance defined by (4-2-4). Thus we have
\[
(\Delta m)^2 = \sum_{t} \sum_{s} w_t w_s \sigma^2 \quad (4-2-15)
\]
\[
= \sum w_t^2 \sigma^2 \quad (4-2-16)
\]
or
\[
\Delta m = \sigma \left(\sum w_t^2\right)^{1/2} \quad (4-2-17)
\]
Now let us examine this final result for \( n \) weights each of size \( 1/n \). For this case, we get
\[
\Delta m = \sigma \left(\frac{1}{n}\right)^{1/2} = \frac{\sigma}{\sqrt{n}} \quad (4-2-18)
\]
This is the most important property of random numbers which is not intuitively obvious. For a zero mean situation it may be expressed in words: “\( n \) random numbers of unit magnitude add up to a magnitude of about the square root of \( n \).”

When one is trying to estimate the mean of a random series which has a time-variable mean, one faces a basic dilemma. If one includes a lot of numbers in the sum to get \( \Delta m \) small, then \( m \) may be changing while one is trying to measure it. In contrast, \( \bar{m} \) measured from a short sample of the series might deviate greatly from the true \( m \) (defined by an infinite sum over the ensemble at any point in time).

This is the basic dilemma faced by a stockbroker when a client tells him, “Since the market fluctuates a lot I’d like you to sell my stock sometime when the price is above the mean selling price.”

If we imagine that a time series is sampled every \( \tau \) seconds and we let \( \Delta t = nt \) denote the length of the sample then (4-2-18) may be written as
\[
(\Delta m)^2 \Delta t = \sigma^2 \tau \quad (4-2-19)
\]
It is clearly desirable to have both \( \Delta m \) and \( \Delta t \) as small as possible. If the original random numbers \( x_i \) were correlated with one another, for example, if \( x_i \) were an approximation to a continuous function, then the sum of the \( n \) numbers would not cancel to root \( n \). This is expressed by the inequality
\[
(\Delta m)^2 \Delta t \geq \sigma^2 \tau \quad (4-2-20)
\]
The inequality (4-2-20) may be called an uncertainty relation between accuracy and time resolution.

In considering other sets of weights one may take a definition of \( \Delta t \) which is more physically sensible than \( \tau \) times the number of weights. For example, the weights \( w_i \) are given by a sampled gaussian function as shown in Fig. 4-1, then \( \Delta t \) could be taken as the separation of half-amplitude points, 1/e points, the time span which includes 95 percent of the area, or it could be given many other “sensible” interpretations. Given a little slop in the definition of \( \Delta m \) and \( \Delta t \), it is clear that the inequality of (4-2-20) is not to be strictly applied.

Given a sample of a zero mean random time series \( x_i \), we may define another series \( y_i \) by \( y_i = x_i^2 \). The problem of estimating the variance \( \sigma^2 = \beta \) of \( x_i \) is identical to the problem of estimating the mean \( m \) of \( y_i \). If the sample is short, we may expect an error \( \delta \beta \) in our estimate of the variance. Thus, in a scientific paper one would like to write for the mean
\[
m = \bar{m} \pm \Delta m \quad (4-2-21)
\]
\[
= \bar{m} \pm \sigma / \sqrt{n} \quad (4-2-22)
\]
but since the variance \( \sigma^2 \) often is not known either, it is necessary to use the estimated \( \delta \), that is
\[
m = \bar{m} \pm \delta / \sqrt{n} \quad (4-2-23)
\]
Correlation is the same sort of thing, except $x$ and $y$ are scalar random variables, so instead of having a vector subscript their subscript is the implicit ensemble subscript. Correlation is defined

$$c = \frac{E(xy)}{\sqrt{E(x^2)E(y^2)}}$$

In practice one never has an ensemble. There is a practical problem when the ensemble average is simulated by averaging over a sample. The problem arises with small samples and is most dramatically illustrated for a sample with only one element. Then the sample correlation is

$$\hat{c} = \frac{xy}{|x| |y|} = \pm 1$$

regardless of what value the random number $x$ or the random number $y$ should take. In fact, it turns out that the sample correlation will always scatter away from zero.

No doubt this accounts for many false "discoveries." The topic of bias and variance of coherency estimates is a complicated one, but a rule of thumb seems to be to expect bias and variance of $\hat{c}$ on the order of $1/\sqrt{n}$ for samples of size $n$.

**EXERCISES**

1. Suppose the mean of a sample of random numbers is estimated by a triangle weighting function, i.e.,

$$\bar{x} = \frac{1}{2n} \sum_{i=0}^{n-1} (n-i)x_i$$

Find the scale factor $s$ so that $E(\bar{x}) = m$. Calculate $\Delta \bar{x}$. Define a reasonable $\Delta t$.

2. A random series $x$ with a possibly time-variable mean may have the mean estimated by the feedback equation

$$\bar{x}_t = (1 - e)\bar{x}_{t-1} + bx_t$$

(a) Express $\bar{x}_t$ as a function of $x_1, x_2, \ldots, x_T$, and not $\bar{x}_{t-1}$.

(b) What is $\Delta t$, the effective averaging time?

(c) Find the scale factor $b$ so that if $m = m_0$, then $E(\bar{x}_T) = m$.

(d) Compute the random error $\Delta \bar{x} = [E(\bar{x} - m)^2]^{1/2}$ [answer goes to $\sigma(1/2)^{1/2}$ as $e$ goes to zero].

(e) What is $\Delta \bar{x}/\Delta t$ in this case?

3. Show that

$$\frac{\Delta \bar{x}}{\Delta t} = \frac{1}{\Delta t} \left[ E(x^2) - \sigma^2 \right]$$

4. Define the behavior of an independent zero-mean-time series $x$, by defining the probabilities that various amplitudes will be attained. Calculate $E(x_t), E(x_t^3), (\Delta P)^3$. 

$$\sum_{i=1}^{\Delta t} c_i$$

5. Verify the formula. 

$$\frac{1}{n} \sum_{i=1}^{n} c_i$$
If you have taken a course in probability theory, use a gaussian probability density function for \( x_i \). HINT:

\[
P(x) = \frac{1}{\sigma \sqrt{2\pi}} e^{-x^2/2\sigma^2}
\]

and

\[
\int_{-\infty}^{\infty} e^{x^2} e^{-x^2} \, dx = \frac{1}{\sqrt{\pi}} \int_{-\infty}^{\infty} e^{-x^2/a} \, dx
\]

4-3 FREQUENCY-STATISTICAL RESOLUTION

Observations of sea level for a long period of time can be summarized in terms of a few statistical averages such as the mean height \( m \) and the variance \( \sigma^2 \). Another important kind of statistical average for use on such geophysical time series is the power spectrum. Some mathematical models explain only statistical averages of data and not the data themselves. In order to recognize certain pitfalls and understand certain fundamental limitations on work with power spectra, we first consider an idealized example.

Let \( x_i \) be a time series made up of independently chosen random numbers. Suppose we have \( n \) of these numbers. We can then define the data sample polynomial \( X(Z) \)

\[
X(Z) = x_0 + x_1 Z + x_2 Z^2 + \cdots + x_{n-1} Z^{n-1}
\]

We can now make up a power spectral estimate \( \hat{R}(Z) \) from this sample of random numbers by

\[
\hat{R}(Z) = \frac{1}{n} \overline{X(Z) X(Z)}
\]

The difference between this and our earlier definition of spectrum is that a power spectrum has the divisor \( n \) to keep the expected result from increasing linearly with the somewhat arbitrary sample size \( n \).

The definition of power spectrum is the expected value of \( \hat{R} \), namely

\[
R(Z) = E[\hat{R}(Z)]
\]

It might seem that a practical definition would be to let \( n \) tend to infinity in (4-3-2).

Such a definition would lead us into a pitfall which is the main topic of the present section. Specifically, from Fig. 4-2 we conclude that \( R(Z) \) is a much fuzzier function than \( R(Z) \), so that

\[
R(Z) \neq \lim_{n \to \infty} \hat{R}(Z)
\]

To understand why this is so, we identify coefficients of like powers of \( Z \) in (4-3-2).

\[
\hat{R}_k = \frac{1}{n} \sum_{i=0}^{n-k-1} x_i x_{i+k} \quad k = 0 \text{ to } n - 1
\]

enabling us to write (4-3-2) for real time series \( x_i = \tilde{x}_i \), as

\[
\hat{R} = \hat{R}_0 + 2 \sum_{k=1}^{n-1} \hat{R}_k \cos ku
\]

Let us examine (4-3-6) for large \( n \). To do this, we will need to know some of the statistical properties of the random numbers. Let them have zero mean \( m = E(x_i) = 0 \) and let them have known constant variance \( \sigma^2 = E(x_i^2) \) and recall our assumption of independence which means that \( E(x_i x_{i+k}) = 0 \) if \( k \neq 0 \). Because of random fluctuations, we have learned to expect that \( \hat{R}_0 \) will come out to be \( \sigma^2 \) plus a random fluctuation component which decreases with sample size as \( 1/\sqrt{n} \), namely

\[
\hat{R}_0 = \sigma^2 \pm \frac{\sigma^2}{\sqrt{n}}
\]

Likewise, \( \hat{R}_k \) should come out to be zero but the definition (4-3-5) leads us to expect a fluctuation component

\[
\hat{R}_k = \pm \frac{n-1}{n} \frac{\sigma^2}{\sqrt{n}}
\]
For the kth correlation value \( k > 1 \) we expect a fluctuation of order

\[
\hat{r}_k = \pm \frac{n-k}{\sqrt{n}} \frac{\sigma^2}{n}
\]

Equation (4-3-7c) for a particular set of random numbers is displayed in Fig. 4-3. Now one might imagine that as \( n \) goes to infinity the fluctuation terms vanish and (4-3-2) takes the limiting form \( \hat{R} = \sigma^2 \). Such a conclusion is false. The reason is that although the individual fluctuation terms go as \( 1/\sqrt{n} \) the summation in (4-3-6) contains \( n \) such terms. Luckily, these terms are randomly canceling one another so the sum does not diverge as \( \sqrt{n} \). We recall that the sum of \( n \) random signed numbers of unit magnitude is expected to add up to a random number in the range \( \pm \sqrt{n} \). Thus the sum (4-3-6) adds up to

\[
\hat{R} \approx (1 \pm \sqrt{n}) \frac{\sigma^2}{\sqrt{n}} = (1 \pm 1)\sigma^2
\]

This is the basic result that a power spectrum estimated from the energy density of a sample of random numbers has a fluctuation from frequency to frequency and from sample to sample which is as large as the expected spectrum.

It should be clear that letting \( n \) go to infinity does not take us to the theoretical result \( \hat{R} = \sigma^2 \). The problem is that, as we increase \( n \), we increase the frequency resolution but not the statistical resolution. To increase the statistical resolution we need to simulate ensemble averaging. There are two ways to do this: (1) Take the sample of \( n \) points and break it into \( k \) equal-length segments of \( n/k \) points each.

Compute an \( R(t) \) for each segment and then add all \( k \) of the \( R(t) \) together, or (2) form \( R(t) \) from the \( n \)-point sample. Of the \( n/2 \) independent amplitudes, replace each one by an average over its \( k \) nearest neighbors. Whichever method, (1) or (2), is used it will be found that \( \Delta R = 0.5k/n \) and \( (\Delta R)^2 = \text{inverse of number of degrees of freedom averaged over } 1/k \). Thus, we have

\[
\Delta R = \frac{(\Delta R)^2}{p} = 0.5p
\]

If some of the data are not used, or are not used effectively, we get the usual inequality

\[
\Delta R = \frac{(\Delta R)^2}{p} \geq \frac{0.5}{\pi}
\]

Thus we see that, if there are enough data available (\( n \) large enough), we can get as good resolution as we like. Otherwise, improved statistical resolution is at the cost of frequency resolution and vice versa.

We are right on the verge of recognizing a resolution tradeoff, not only between \( \Delta R \) and \( \Delta \) but also with \( \Delta t = \pi \), the time duration of the data sample. Recognizing now that the time duration of our data sample is given by \( \Delta t = \sqrt{n} \), we obtain the inequality

\[
\Delta R = \frac{(\Delta R)^2}{p} \geq \frac{1}{2}
\]

This inequality will be further interpreted and derived from a somewhat different point of view in the next section.

In time-series analysis we have the concept of coherency which is analogous to the concept of correlation defined in Sec. 4-2. There we had for two random variables \( x \) and \( y \) that

\[
C = \frac{E(xy)}{\sqrt{E(x^2)E(y^2)}}
\]

Now if \( x_i \) and \( y_i \) are time series, they may have a relationship between them which depends on time-delay, scaling, or even filtering. For example, perhaps \( y_i = f_i(x_i) + n_i \) where \( f_i(x_i) \) is a filter and \( n_i \) is unrelated noise. The generalization of the correlation concept is to define coherency by

\[
C = \frac{E(x_i^*)y_i}{\sqrt{E(x_i^*x_i)E(y_i^*y_i)}}
\]

Correlation is a real scalar. Coherency is complex and expresses the frequency dependence of correlation. In forming an estimate of coherency it is always essential to simulate some ensemble averaging. Note that if the ensemble averaging were to be omitted, the coherency (squared) calculation would give

\[
|C|^2 = \frac{E(x_i^*y_i)}{E(x_i^*x_i)E(y_i^*y_i)} = +1
\]
which states that the coherency squared is +1 independent of the data. Because correlation scatters away from zero we find that coherency squared is biased away from zero.

### 4.4 TIME-FREQUENCY-STATISTICAL RESOLUTION

Many time functions are not completely random from point to point, but become more random when viewed over a longer time scale. A popular mathematical model embodying this concept is to make a so-called *stationary time series* by putting random numbers into a filter as depicted in Fig. 4-4. The input $x_t$ may be *independent random numbers or white light*. [The two terms mean nearly the same thing in practice but the first term is the stronger; it means that $x_t$ is in no way related to $x_t$ if $t \neq s$, whereas white light means that $E(x_t, x_s) = 0$ if $t \neq s$.] The output random time series $y_t$ may vary rather slowly from point to point if $f_t$ is a low-pass filter. This is the usual case when we are modeling continuous time functions. The random time series may be called a *stationary* random time series if neither the filter nor any property of the random numbers (such as $m$ or $\sigma$) vary with time. Stationarity is often assumed even where it cannot be strictly true.

This model will be useful later when we consider the problem of predicting a future point on $y_t$ from knowledge of past values. Now we will use the model to examine the estimation of the spectrum of $y_t$ given a sample of $n$ points of $y_t$. To begin with, we have a very precise meaning for the spectrum of $y_t$. We have

$$Y(Z) = F(Z)X(Z) \quad (4.4-1)$$

and its conjugate

$$Y(1/Z) = F(1/Z)X(1/Z) \quad (4.4-2)$$

Multiplying (4.4-1) by (4.4-2) we get

$$Y(1/Z) = X(1/Z)X(Z)F(1/Z)F(Z) \quad (4.4-3)$$

but, from the previous section, we learned that $E(X_t X_s) = \sigma^2$. Considering $\sigma^2$ to be unity, we see that the expected power spectrum of the output $Y$ is the energy spectrum of the filter $F$. The overall situation is depicted in Fig. 4-5. The interesting question is how well can we estimate the spectrum when we start with an $n$-point sample of $y_t$. We will describe three computationally different methods, all having the same fundamental limitations.

The first method uses a bank of filters as shown in Fig. 4-6. When random numbers excite the narrowband filter, the output is somewhat like a sine wave. It differs in one important respect. A sine wave has constant amplitude, but the output of a narrowband filter has an amplitude which swings over a range. This is illustrated in Fig. 4-7. If the bandwidth is narrow, the amplitude changes slowly. If the impulse response of the filter has duration $\Delta t_{filter}$, then the output amplitude at time $t$ will be randomly related to the amplitude at time $t + \Delta t_{filter}$. Thus,
in statistical averaging, it is not the number of time points but the number of intervals $\Delta t_{\text{filter}}$ which enhance the reliability of the average. Consequently, the decay time of the integrator $\Delta t_{\text{integrator}}$ will generally be chosen to be greater than $\Delta t_{\text{filter}} = 1/\Delta f$. The variability $\Delta p$ of the output $p$ decreases as $\Delta t_{\text{integrator}}$ increases. Since $\eta$, has independent values over time spans of about $\Delta t_{\text{filter}} = 1/\Delta f$, then the "degrees of freedom" smoothed over can be written $\Delta f_{\text{integrator}} / \Delta t_{\text{filter}} = \Delta f / \Delta t_{\text{integrator}}$. The variability $\Delta p / p$ is proportional to the inverse square root of the number of degrees of freedom, and so we get

$$\left( \frac{\Delta p}{p} \right)^2 = \frac{1}{\Delta f \Delta t_{\text{integrator}}}$$

or, introducing the usual inequality,

$$\Delta t \Delta f \left( \frac{\Delta p}{p} \right)^2 > 1 \quad (4-4-4)$$

The inequality (4-4-4) indicates the three-parameter uncertainty which is fundamental to estimating power spectra of random functions. Two other methods of estimating the spectrum of $Y$ from a sample of length $n$ are exactly the same as the methods described in Sec. 4-3 as ways of estimating the spectrum of white light. In fact, (4-4-4) turns out to be the same as (4-3-9).

The usual interpretation is that to attain a frequency resolution of $\Delta f$ and a relative accuracy of $\Delta p / p$ a time sample of duration at least $\Delta t \geq 1 / [\Delta f (\Delta p / p)^2]$ will be required. Although this sort of interpretation is generally correct, it will break down for highly resonant series recorded for a short time. Then the data sample may be predictable an appreciable distance off its ends so that the effective $\Delta t$ is somewhat (perhaps appreciably) larger than the sample length.

**Exercises**

1. It is popular to taper the ends of a data sample so that the data go smoothly to zero at the ends of the sample. Choose a weighting function and discuss in a semiquantitative fashion its effect on $\Delta t$, $\Delta f$, and $\Delta p / p$.

2. Answer the question of Exercise 1, where the autocorrelation function is tapered rather than the data sample.

**4-5 The Central-Limit Theorem**

The central-limit theorem of probability and statistics is perhaps the most important theorem in the fields of probability and statistics. A derivation of the central limit theorem explains why the gaussian probability function is so frequently encountered in nature; not just in physics but also in the biological and social sciences. No experimental scientist should be unaware of the basic ideas behind this theorem. Although the result is very deep and is even today the topic of active research, we can get to the basic idea quite easily.

One way to obtain random integers from a known probability function is to write integers on slips of paper and place them in a hat. Draw one slip at a time. After each drawing replace the slip in the hat. The probability of drawing the integer $i$ is given by the ratio of $a_i$ of the number of slips containing the integer $i$ divided by the total number of slips. Obviously the sum over $i$ of $a_i$ must be unity. Another way to get random integers is to throw one of a pair of dice. Then all $a_i$ equal zero except $a_1 = a_2 = a_3 = a_4 = a_5 = a_6 = 1$. The probability that the integer $i$ will occur on the first drawing and the integer $j$ will occur on the second drawing is $a_1 a_2$. If you draw two slips or throw a pair of dice, then the probability that the sum of $i$ and $j$ equals $k$ is readily seen to be

$$c_k = \sum_i a_{i+k} \quad (4-5-1)$$

Since this equation is a convolution, we may look into the meaning of the Z transform

$$A(Z) = \cdots a_{-2} z^{-2} + a_0 + a_1 z + a_2 z^2 + \cdots \quad (4-5-2)$$
In terms of \( Z \) transforms the probability that \( i \) plus \( j \) equals \( k \) is simply the coefficient of \( Z^k \) in

\[
C(Z) = A(Z)^n
\]  

(4-5-3)

Obviously, if we add \( n \) of the random numbers, the probability that the sum of them equals \( k \) is given by the coefficient of \( Z^k \) in

\[
G(Z) = [A(Z)]^n
\]  

(4-5-4)

The central-limit theorem of probability says that as \( n \) goes to infinity the polynomial \( G(Z) \) goes to a special form, almost regardless of the specific polynomial \( A(Z) \). The specific form is such that a graph of the coefficients of \( G(Z) \) comes closer and closer to fitting under the envelope of the bell-shaped gaussian function. Let us see why this happens. Our development will lack a mathematical rigor because the theorem is not always true. There are pathological \( A \) functions which do not result in \( G \) tending to gaussian. Despite the fact that some of the pathological functions sometimes turn up in applications, we will not take the time here to look at such instances.

Consider the size of \( A(Z) \) for real \( \omega \). If \( \omega = 0 \), the sum of the terms of \( A(Z) \) may be visualized in the complex plane as a sum of vectors \( a_k e^{ik\omega} \) all pointing in the positive real direction. If \( \omega \neq 0 \) the vectors point in different directions. This is shown in Fig. 4-8.

In raising \( A(e^{iz}) \) to the \( n \)th power, the values of \( \omega \) of greatest concern are those near \( \omega = 0 \) where \( A \) is largest—because in any region where \( A \) is small \( A^n \) will be extremely small. Near \( \omega = 0 \) or \( Z = 1 \) we may expand \( A(Z) \) in a power series in \( \omega \)

\[
A(e^{iz}) = A(1) + \frac{\partial A}{\partial \omega} \bigg|_1 \omega + \frac{\partial^2 A}{\partial \omega^2} \bigg|_1 \frac{\omega^2}{2!} + \cdots
\]  

(4-5-5)

Note that the coefficients of this power series are proportional to the moments \( m_i \) of the probability function; that is

\[
A(e^{iz}) = \sum_k a_k e^{ik\omega}
\]  

(4-5-6)

\[
A(1) = \sum_k a_k = m_0 = 1
\]  

(4-5-7)

\[
\frac{\partial}{\partial \omega} A = \sum_k ik a_k e^{ik\omega}
\]  

(4-5-8)

\[
\frac{\partial^2}{\partial \omega^2} A \bigg|_0 = \sum_k ik^2 a_k = im_1
\]  

(4-5-9)

\[
\frac{\partial^2}{\partial \omega^2} A \bigg|_0 = -\sum_k k^2 a_k = -m_2
\]  

(4-5-10)

When we raise \( A(Z) \) to the \( n \)th power we will make the conjecture that only the given first three terms of the power series expansion will be important. (This assumption clearly fails if any of the moments of the probability function are infinite.) Thus, we are saying that as far as \( G \) is concerned the only important things about \( A \) are its mean value \( m = m_1 \) and its second moment \( m_2 \). If this is really so, we may calculate \( G \) by replacing \( A \) with any function \( B \) having the same mean and same second moment as \( A \). We may use the simplest function we can find. A good choice is the so called binomial probability function given by

\[
B = \frac{Z^n(Z^n + Z^{-n})}{2}
\]  

(4-5-11)

\[
e^{i(n+\sigma)s} + e^{i(n-\sigma)s}
\]  

(4-5-12)

Let us verify its first moment

\[
\frac{\partial B}{\partial \omega} = \frac{i(l(m+\sigma)e^{i(n+\sigma)s}+(m-\sigma)e^{i(n-\sigma)s})}{2}
\]  

(4-5-13)

\[
\frac{\partial^2 B}{\partial \omega^2} \bigg|_0 = im
\]  

(4-5-14)

Now let us verify its second moment

\[
\frac{\partial^2 B}{\partial \omega^2} \bigg|_0 = \frac{(m+\sigma)^2+(m-\sigma)^2}{2}
\]  

(4-5-15)

\[
\frac{\partial^2 B}{\partial \omega^2} \bigg|_0 = -m_2 + \sigma^2
\]  

(4-5-16)

Hence, \( \sigma \) should be chosen so that

\[
m_2 = m^2 + \sigma^2
\]  

(4-5-17)
Of course, we cannot expect that $m$ and $\sigma$ will necessarily turn out to be integers; therefore (4-5-11) will not necessarily be a $Z$ transform in the usual sense. It does not really matter; we simply interpret (4-5-11) as saying:

1. The probability of drawing the number $m + \sigma$ is one-half.
2. The probability of $m - \sigma$ is one-half.
3. The probability of any other number is zero.

Now, raising $(Z^N + Z^{-N})$ to the $n$th power gives a series in powers of $Z^0$ whose coefficients are symmetrically distributed about $Z$ to the zero power and whose magnitudes are given by the binomial coefficients. A sketch of the coefficients of $B(Z)^n$ is given in Fig. 4.9.

We will now see how, for large $n$, the binomial coefficients asymptotically approach a gaussian. Approaching this limit is a bit tricky. Obviously, the sum of $n$ random integers will diverge as $\sqrt{n}$. Likewise the coefficients of powers of $Z$ in $(1 + Z)^n$ individually get smaller while the number of coefficients gets larger. We recall that in time series analysis we used the substitution $Z = e^{\lambda \Delta t}$. We commonly chose $\Delta t = 1$, which had the meaning that data points were given at integral points on the time axis. In the present probability theory application of $Z$ transforms, the choice $\Delta t = 1$ arises from our original statement that the numbers chosen randomly from the slips of paper were integers. Now we wish to add $n$ of these random numbers together; and so, it makes sense to rescale the integers to be integers divided by $\sqrt{n}$. Then we can make the substitution $Z = e^{\lambda \Delta t} = e^{\lambda n/\sqrt{n}}$.

The coefficient of $Z^n$ now refers to the probability of drawing the number $n/\sqrt{n}$. Raising $(Z^N + Z^{-N})/2$ to the $n$th power to find the probability distribution for the sum of $n$ independently chosen numbers, we get

$$B(Z)^n = \left(\frac{Z^N + Z^{-N}}{2}\right)^n = \left(\cos \frac{\sigma \Delta t}{\sqrt{n}}\right)^n$$

Using the first term of the series expansion for cosine we have

$$B(Z)^n = \left(1 - \sigma \Delta t^2/2n\right)^n$$

Using the well-known fact that $(1 + x/n)^n \approx e^x$, we have for large $n$

$$B(Z)^n \approx e^{-\sigma \Delta t^2/2}$$  \hspace{1cm} (4-5-18)

The probability that the number $\tau$ will result from the sum is now found by inverse Fourier transformation of (4-5-18). The Fourier transform of the gaussian (4-5-18) may be looked up in a table of integrals. It is found to be the gaussian

$$p(\tau) = \frac{1}{\sigma \sqrt{2\pi}} e^{-\sigma^2 \tau^2/2\pi}$$

4-6 CONFIDENCE INTERVALS

It is always important to have some idea of the size and influence of random errors. It is often important to be able to communicate this idea to others in the form of a statement such as

$$m = \bar{m} \pm \sigma /\sqrt{n}$$

In a matter of any controversy you may be called upon to define a probability that the true mean lies in your stated interval; in other words, what is your confidence that $m$ lies in the interval

$$\bar{m} - \Delta m < m < \bar{m} + \Delta m$$

Before you can answer questions about probability, it is necessary to make some assumptions and assertions about the probability functions which control your random errors. The assertion that errors are independent of one another is your most immediate hazard. If they are not, as is often the case, you may be able to readjust the numerical value of $n$ to be an estimate of the number of independent errors. We did something like this in time series analysis when we took $n$ to be not the number of points on the time series but the number of intervals of length $\Delta t_{	ext{real}}$. The second big hazard in trying to state a confidence interval is the common assumption that, because of the central-limit theorem and for lack of better information, the errors follow a gaussian probability function. If in fact the data errors include blunders which arise from human errors or blunders from transient electronic equipment difficulties, then the gaussian assumption can be very wrong and can lead you into serious errors in geophysical interpretation. Some useful help is found in the field of nonparametric statistics (see, for example, Ref. [18]).

To begin with, it is helpful to rephrase the original question into one involving the median rather than the mean. The median $m_1$ is defined as that value which is expected to be less than half of the population and greater than the other half. In many—if not most—applications the median is a ready, practical substitute for
the arithmetic mean. The median is insensitive to a data point, which, by some blunder, is near infinity. In fact, median and mean are equal when the probability function is symmetrical. For a sample of N numbers \((x_i, i = 1, 2, \ldots, n)\), the median \(m_1\) may be estimated by reordering the numbers from smallest to largest and then selecting the number in the middle as the estimate of the median \(m_1\). Specifically, let the reordered \(x_i\) be denoted by \(x'_i\) where \(x'_i \leq x'_{i+1}\). Then we have \(m_1 = x'_{n/2}\). Now it turns out that without knowledge of the probability density function for the random variables \(x_i\) we will still be able to compute the probability that the true median \(m_1\) is contained in the interval

\[
x'_{n/2 - a\sqrt{n}} < m_1 < x'_{n/2 + a\sqrt{n}}
\]  \hspace{1cm} (4-6-1)

For example, set \(\alpha = 1\) and \(N = 100\), the assertion is that we can now calculate the probability that the true median \(m_1\) lies between the 40th and the 60th percentile of our data. The trick is this: Define a new random variable

\[
y = \text{step}(x - m_1)
\]  \hspace{1cm} (4-6-2)

The step function equals +1 if \(x > m_1\) and equals 0 if \(x < m_1\). The new random variable \(y\) takes on only values of zero and one with equal probability; thus we know its probability function even though we may not know the probability function for the random variable \(x\). Now define a third random variable \(s\) as

\[
s = \frac{1}{\alpha} \sum_{i=1}^{n} y_i
\]

Since each \(y_i\) is zero or one, then \(s\) must be an integer between zero and \(n\). Furthermore, the probability that \(s\) takes the value \(j\) is given by the coefficient of \(Z^j\) of \((1 + Z/\alpha)^n\). Now the probability that \(s\) lies in the interval \(n/2 - \alpha\sqrt{n} < s < n/2 + \alpha\sqrt{n}\) is readily determined by adding the required coefficients of \(Z^j\), and this probability is by definition equal to the probability that the median \(m_1\) lies in the interval (4-6-1). For \(\alpha = 1\) and large \(n\) this probability works out to about 95 percent.

Familiarity with matrices is essential to computer modeling in both physical and social sciences. As this is a big subject covered by many excellent texts at all levels, our review will be a quick one. We focus on those properties required in the succeeding chapters. We avoid proofs, and although constructions given should be useful in most situations, there will be occasional matrices (which we will dismiss as pathological cases) in which our constructions will fail. In practice, the user should always check computed results. Unfortunately, the so-called pathological cases arise in practice far more often than might be expected. When matrix difficulties arise, the first tendency of the scientist is to use a higher-precision arithmetic. In the author's experience, physically meaningful calculations rarely require high precision. When higher precision seems to be needed, it is often because something is happening physically which shows that the problem being solved is a poorly posed problem. If a slight change in the problem should not make a drastic change in the answer, then it may happen that a different organization of the calculations will obviate the need for high precision. Anyway, our discussion here will focus on the nonpathological cases, but the reader is warned that pathological cases will certainly be encountered in practice and when they are they will be a stern test of the reader's mathematical knowledge and physical insight.