(1980) for the same crystal and radiation were used to derive the relative magnitudes of the extinction corrections for the two sets of data. The reasonableness of using the same extinction tensor can be understood from the fact that the crystal had not been removed from its support between the two experiments; consequently the directions of the normals to the planes of diffraction for corresponding reflections remained unchanged. To obviate the need for making TDS corrections, the peak widths of the monochromator data were adjusted to be identical to their equivalents in the filtered data with peak maxima in identical positions relative to the peak widths. It was felt prudent not to include the $\{111\}$ reflections in the determination of $K$ since the absorption edge of the $\beta$ filter fell beneath their peaks. No reflections were rejected on the basis of the difference between extinction corrections [maximum value of $y p c$ was 0.0007 corresponding to $y_{\text {ext }}=0.8885$ (monochromator) and 0.8891 (filter) for the 220 reflectionl. The value of $K$ determined under these conditions was 0.805 (11).

In summary, our experimentally derived $K$ for Cu $K \alpha$ radiation falls between the expected values for an ideally mosaic and perfect graphite monochromator. In contrast, the values we have obtained for Mo $K \alpha$ and $\mathrm{Ag} K \alpha$ radiation fall outside this range. Moreover, and contrary to the results of Jennings (1968) and Le Page, Gabe \& Calvert (1979), they are smaller than those expected for an ideally mosaic crystal.

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# On the Polarization Factor for Crystal-Monochromated X-radiation. III. A Weighting Scheme for Products 

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#### Abstract

The minimum-variance weighting scheme for calculating the weighted mean of a quantity for the following unusual case is analysed. The quantity is derived from the product of two independent quantities for each of which a set of data is available. All the products between the two sets of data are taken into consideration. The application to measurements of the polarization state of X -radiation is outlined.

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## Introduction

In a previous paper (Vincent \& Flack, 1980), a method is described for determining the polarization ratio, $K$, of X-radiation reflected from a monochromator. This value is estimated by taking the weighted mean of $M \times$ $N$ individual values $K_{i j}$, i.e.

$$
\begin{equation*}
\hat{K}=\sum_{i=1}^{M} \cdot \sum_{j=1}^{N} a_{i j} K_{i j} / \sum_{i=1}^{M} \sum_{j=1}^{N} a_{i j} \tag{1}
\end{equation*}
$$

where $a_{i j}$ is the weight attributed to $K_{i j}$.

The $K_{i j}$ are obtained from the relationship

$$
\begin{equation*}
K_{i j}=c_{i j} A_{i} B_{j}+d_{i j} \tag{2}
\end{equation*}
$$

where $c_{i j}$ and $d_{i j}$ are known constants and $A_{i},(i=1,2$, $\ldots, M)$ and $B_{j},(j=1,2, \ldots, N)$ are independent series of data (with e.s.d.'s) derived from intensity measurements. The subject of this paper is the choice of the 'proper' or 'best' weighting scheme, $a_{i j}$, to apply in the above-defined problem.

A simple choice for the weighting scheme would be to use

$$
a_{i j}=1 / \sigma_{i j}^{2},
$$

where $\sigma_{i j}^{2}$ could be estimated from

$$
\sigma_{i j}^{2}=B_{j}^{2} \sigma_{i}^{2}+A_{i}^{2} \sigma_{j}^{2}
$$

where $\sigma^{2}$ represents the population variance of the subscripted quantity (Hamilton, 1964). From a theoretical point of view this weighting scheme is unsatisfactory as it treats the $M \times N$ values of $K_{i j}$, derived from $M+N$ values of $A_{i}$ and $B_{j}$, as independent observations whereas the error disturbances in the $K_{i j}$ must be correlated.

From a practical point of view this simple weighting scheme was tested in a simulation with model data of $A_{i}$ and $B_{j}$ derived from a chosen value of $K$ and with a Gaussian pseudo-random error added in. The value of $\hat{K}$ was frequently unreasonably different from the model value of $K$. The discrepancy between $\hat{K}$ and the model $K$ increased with decreasing wavelength.

We were thus driven to determine a more satisfactory weighting scheme. The criterion to be applied (Hamilton, 1964) is that $\hat{K}$ should be an unbiased estimator of $K$ and that its variance should be a minimum. The steps in the analysis, given below, to determine the $a_{i j}$ are as follows. The error disturbances on the $A_{i}$ and $B_{j}$, and the expected values of their first-second-, third- and fourth-order moments are defined. It is then shown that, for any weighting scheme, $\hat{K}$ is an unbiased estimator of $K$. Next, the variance of $\hat{K}$ is determined and this value is minimized with respect to the $a_{i j}$. An expression for the weighting scheme in terms of the population variances of the $A_{i}$ and $B_{j}$ is thus obtained. We next suppose that estimates of these population variances are available but that these are only known within a scale factor. It is shown that the variance of $\hat{K}$ is proportional to this (as yet unknown) scale factor. It may, however, be estimated by considering the expected value of the weighted squared deviations of the $K_{i j}$ about $\hat{K}$. Finally, a practical method for calculating the weighting scheme in a computer program is presented. As much of the algebra is straightforward, only the important intermediate equations are given in the text.

## Analysis

Let us write

$$
\begin{equation*}
\gamma_{i j}=A_{i} B_{j}(i=1,2, \ldots, M ; j=1,2, \ldots, N) \tag{3}
\end{equation*}
$$

and $K$ for the population mean of the desired quantity. In the following analysis it will be necessary to have multiple summations over the observations. As a matter of convention, $i, k$ and $m$ take values from 1 to $M$, and $j, l$ and $n$ take values from 1 to $N$. It is convenient also to choose weights $a_{i j}$ so that

$$
\begin{equation*}
\sum_{i} \sum_{j} a_{i j}=1 \tag{4}
\end{equation*}
$$

## Definition of errors

Let us assume that our $A_{i}, B_{j}$ of (2) are subject to random errors. We may write

$$
\begin{equation*}
A_{i}=E_{i}+e_{i} \quad \text { and } \quad B_{j}=F_{j}+f_{j} \tag{5}
\end{equation*}
$$

where $E_{i}$ and $F_{j}$ are the population means and $e_{i}$ and $f_{j}$ represent the random errors on these quantities. Taking expected values, $\varepsilon\left\{A_{i}\right\}=E_{i}$ and $\varepsilon\left\{B_{j}\right\}=F_{j}$, we assume that
$\varepsilon\left\{e_{i}\right\}=\varepsilon\left\{f_{j}\right\}=0 ; \quad i=1, \ldots, M, \quad j=1, \ldots, N$.
Further, we will assume that the individual error terms have a finite variance

$$
\begin{equation*}
\varepsilon\left\{e_{i}^{2}\right\}=\sigma^{2}\left(E_{i}\right) \quad \text { and } \quad \varepsilon\left\{f_{j}^{2}\right\}=\sigma^{2}\left(F_{j}\right) \tag{6b}
\end{equation*}
$$

and that the error disturbances are independent one from another, i.e.

$$
\begin{align*}
& \varepsilon\left\{e_{i} f_{j}\right\}=0, \quad \varepsilon\left\{e_{i} e_{k}\right\}=0 \text { for } i \neq k, \\
& \varepsilon\left\{f_{j} f_{l}\right\}=0 \text { for } j \neq l, \\
& \varepsilon\left\{e_{i}^{2} f_{j}^{2}\right\}=\sigma^{2}\left(E_{i}\right) \sigma^{2}\left(F_{j}\right), \\
& \varepsilon\left\{e_{i} f_{j} e_{k} f_{l}\right\}=0 \text { for } i \neq k \text { or } j \neq l, \\
& \varepsilon\left\{e_{i} f_{j} e_{k}\right\}=0 \text { and } \varepsilon\left\{e_{i} f_{j} f_{l}\right\}=0 . \tag{6c}
\end{align*}
$$

## The expected value of $\hat{K}$

Substituting (5) into (2) and the result into (1), with (4), we obtain

$$
\begin{equation*}
\left.\hat{K}=K+\sum_{i} \sum_{j} a_{i j} c_{i j} \mid E_{i} f_{j}+F_{j} e_{i}+e_{i} f_{j}\right\rfloor \tag{7}
\end{equation*}
$$

Taking expected values and applying (6) it follows that $\left.\varepsilon_{\}} \hat{K}\right\}=K$. i.e. $\hat{K}$ is an unbiased estimator of $K$.

## The variance of $\hat{K}$

Squaring (7) we obtain

$$
\begin{align*}
(\hat{K}-K)^{2}=\sum_{i=1}^{M} & \sum_{j=1}^{N} \sum_{k=1}^{M} \sum_{l=1}^{N} a_{i j} a_{k l} c_{i j} c_{k l}\left[E_{i} E_{k} f_{j} f_{l}\right. \\
& +F_{j} F_{l} e_{i} e_{k}+e_{i} e_{k} f_{j} f_{l}+2 E_{i} F_{l} f_{j} e_{k} \\
& \left.+2 F_{j} e_{i} e_{k} f_{l}+2 E_{i} e_{k} f_{j} f_{l}\right] \tag{8}
\end{align*}
$$

Taking expected values in (8) and applying (6), one finds

$$
\begin{align*}
\varepsilon\left\{(\hat{K}-K)^{2}\right\}=\sum_{i} & \sum_{j} a_{i j}^{2} c_{i j}^{2}\left[E_{i}^{2} \sigma^{2}\left(F_{j}\right)\right. \\
& \left.+F_{j}^{2} \sigma^{2}\left(E_{i}\right)+\sigma^{2}\left(E_{i}\right) \sigma^{2}\left(F_{j}\right)\right] \\
& +\sum_{i} \sum_{j} \sum_{k}^{k \neq i} a_{i j} a_{k j} c_{i j} c_{k j} E_{i} E_{k} \sigma^{2}\left(F_{j}\right) \\
& +\sum_{i} \sum_{j} \sum_{l}^{l \neq j} a_{i j} a_{i l} c_{i j} c_{i l} F_{j} F_{l} \sigma^{2}\left(E_{i}\right) \tag{9}
\end{align*}
$$

which is an expression for the variance of $\hat{K}$.

## Minimizing the variance of $\hat{K}$

Let us now minimize $\varepsilon\left\{(\hat{K}-K)^{2}\right\}$ subject to the restriction (4) that $\sum_{i} \sum_{j} a_{i j}=1$. This may be achieved by differentiation of (9) including a Lagrangian multiplier. Set

$$
\begin{equation*}
Q=\varepsilon\left\{(\hat{K}-K)^{2}\right\}+\lambda\left(\sum_{i} \sum_{j} a_{i j}-1\right) \tag{10}
\end{equation*}
$$

and the conditions $\partial Q / \partial a_{i j}=0$ produce

$$
\begin{equation*}
a_{i j}=-\eta_{i j}\left[\lambda+c_{i j} u_{j} E_{i} \sigma^{2}\left(F_{j}\right)+c_{i j} v_{i} F_{j} \sigma^{2}\left(E_{i}\right)\right] \tag{11}
\end{equation*}
$$

where

$$
\begin{equation*}
\eta_{i j}=1 / c_{i j}^{2}\left[E_{i}^{2} \sigma^{2}\left(F_{j}\right)+F_{j}^{2} \sigma^{2}\left(E_{i}\right)+2 \sigma^{2}\left(E_{i}\right) \sigma^{2}\left(F_{j}\right)\right],( \tag{12}
\end{equation*}
$$

$$
\begin{equation*}
u_{j}=\sum_{k} a_{k j} c_{k j} E_{k} \quad \text { and } \quad v_{i}=\sum_{l} a_{i l} c_{i l} F_{l} \tag{13}
\end{equation*}
$$

$\lambda$ in (11) may be eliminated by imposing (4). One obtains

$$
\begin{align*}
a_{i j}= & \eta_{i j}\left\{(1+\psi) / \sum_{k} \sum_{i} \eta_{k l}\right. \\
& -c_{i j}\left\lceil u_{j} E_{i} \sigma^{2}\left(F_{j}\right)+v_{i} F_{j} \sigma^{2}\left(E_{i}\right)\right\}_{\}} \tag{14}
\end{align*}
$$

where

$$
\begin{equation*}
\varphi=\sum_{i} \sum_{j} c_{i j} \eta_{i j}\left[u_{j} E_{i} \sigma^{2}\left(F_{j}\right)+v_{i} F_{j} \sigma^{2}\left(E_{i}\right)\right] . \tag{15}
\end{equation*}
$$

Thus, (14) defines the weighting scheme that we need although it is in an implicit form as the $a_{i j}$ also appear
on the right-hand side of (14) through the $u_{j}$ and $v_{i}$ in (13).

## Use of sample variances of $A_{i}$ and $B_{j}$

The weighting scheme defined by (14) requires values of $\sigma^{2}\left(E_{i}\right)$ and $\sigma^{2}\left(F_{j}\right)$ to be known. We will now suppose that we have experimental values of these quantities available $s^{2}\left(E_{i}\right), s^{2}\left(F_{j}\right)$ but that they are only known to within an unknown scale factor, i.e.

$$
\begin{equation*}
\sigma^{2}\left(E_{i}\right)=\beta s^{2}\left(E_{i}\right) \quad \text { and } \quad \sigma^{2}\left(F_{j}\right)=\beta s^{2}\left(F_{j}\right) \tag{16}
\end{equation*}
$$

In order to estimate $\beta$, it is useful to substitute (16) into (12) - (15), omit all terms in $s^{4}$ and of higher order and replace $a_{i j}$ by a proportional quantity $w_{i j}$. We obtain

$$
\begin{gather*}
x_{i j}=\beta \eta_{i j}=1 / c_{i j}^{2}\left[E_{i}^{2} s^{2}\left(F_{j}\right)+F_{j}^{2} s^{2}\left(E_{i}\right)\right]  \tag{17}\\
w_{i j}=x_{i j}\left\{1+\varphi-c_{i j}\left[u_{j} E_{i} s^{2}\left(F_{j}\right)\right.\right. \\
\left.\left.+v_{i} F_{j} s^{2}\left(E_{i}\right)\right]\left(\sum_{k} \sum_{l} x_{k l}\right)\right\}  \tag{18}\\
u_{j}=\frac{\sum_{k} w_{k j} c_{k j} E_{k}}{\sum_{i} \sum_{l} w_{i l}} \text { and } v_{i}=\frac{\sum_{l} w_{i l} c_{i l} F_{l}}{\sum_{k} \sum_{j} w_{k j}}  \tag{19}\\
\varphi=\sum_{i} \sum_{j} c_{i j} x_{i j}\left[u_{j} E_{i} s^{2}\left(F_{j}\right)+v_{i} F_{j} s^{2}\left(E_{i}\right)\right]  \tag{20}\\
\hat{K}=\sum_{i} \sum_{j} w_{i j} K_{i j} / \sum_{i} \sum_{j} w_{i j} \tag{21}
\end{gather*}
$$

and
variance $(\hat{K})=\varepsilon\left\{(\hat{K}-K)^{2}\right\}$

$$
\begin{equation*}
=\beta\left[1+\varphi-\sum_{i} \sum_{j} \frac{w_{i j}^{2}}{x_{i j}} / \sum \sum x_{i j}\right] / \sum_{i} \sum_{j} w_{i j} \tag{22}
\end{equation*}
$$

The above relationships are readily derived when it is realized that

$$
\begin{equation*}
\sum_{i} \sum_{j} x_{i j}=\sum_{i} \sum_{j} w_{i j} \tag{23}
\end{equation*}
$$

## Estimation of $\beta$

Consider $\sum_{i} \sum_{j} w_{i j}\left(K_{i j}-\hat{K}\right)^{2}$ and its expected value. From (7), (2) and (5) one finds that

$$
\begin{align*}
K_{k l}-\hat{K}=c_{k l} & E_{k} f_{l}+c_{k l} F_{l} e_{k} \\
& -\sum_{i} \sum_{j} w_{i j} c_{i j} E_{i} f_{j} / \sum_{i} \sum_{j} w_{i j} \\
& -\sum_{i} \sum_{j} w_{i j} c_{i j} F_{j} e_{i} \mid \sum_{i} \sum_{j} w_{i j} \tag{24}
\end{align*}
$$

Taking expected values in (24) one finds

$$
\begin{align*}
& \varepsilon\left\{\sum_{k} \sum_{l} w_{k l}\left(K_{k l}-\hat{K}\right)^{2}\right\}=\beta\left[\sum_{k} \sum_{l}\left(w_{k l} / x_{k l}\right)-1-\varphi\right. \\
& \left.\quad+\sum_{i} \sum_{j}\left(w_{i j}^{2} / x_{i j}\right) / \sum_{i} \sum_{j} x_{i j}\right] \tag{25}
\end{align*}
$$

which allows the value of $\beta$ to be estimated by assuming that $\sum_{k} \sum_{l} w_{k l}\left(K_{k l}-\hat{K}\right)^{2}$ has attained its expected value.

## Evaluation of $w_{i j}$

The equations (18) for $w_{i j}$ represent a $M N \times M N$ system of linear equations in $M N$ unknowns. Rather than solve directly such a system by matrix inversion, we have preferred an iterative procedure which progressively improves the weights in a series of cycles. Suppose we have an 'old' system of weights $w o_{l j}$ and that a 'new' better system is required by adding in a shift $\Delta w_{i j}$, i.e.

$$
\begin{equation*}
w n_{l j}=w o_{i j}+\Delta w_{i j} . \tag{26}
\end{equation*}
$$

$w n_{i j}$ will be made to obey (18) approximately, keeping in mind that (23) needs to hold, i.e.

$$
\sum_{i} \sum_{j} w n_{i j}=\sum \sum w o_{l j}=\sum \sum x_{i j}
$$

or

$$
\begin{equation*}
\sum \sum \Delta w_{i j}=0 . \tag{27}
\end{equation*}
$$

$u o_{i}, v o_{j}$ and $u n_{i}, v n_{j}$ are the values defined by (19) with the 'old' and 'new' system of weights respectively. Likewise, $\varphi_{o}$ and $\varphi_{n}$ are defined from (20). Substituting (26) into (18) to make $w n_{i j}$ obey (18) exactly one finds,

$$
\begin{align*}
\Delta w_{i j}=- & w o_{i j}+x_{i j}+x_{i j} \varphi_{o} \\
& +x_{i j} \sum_{k} \sum_{l} c_{k l} x_{k l}\left[s^{2}\left(F_{i}\right) E_{k}\left(\sum_{m} \Delta w_{m l} c_{m l} E_{m}\right)\right. \\
& \left.+s^{2}\left(E_{k}\right) F_{l}\left(\sum_{n} \Delta w_{k n} c_{k n} F_{n}\right)\right] / \sum_{k} \sum_{l} x_{k l} \\
& -c_{i j} x_{i j}\left[u o_{j} E_{i} s^{2}\left(F_{j}\right)+v o_{i} F_{j} s^{2}\left(E_{i}\right)\right] \sum_{k} \sum_{l} x_{k l} \\
& -c_{i j} x_{i j}\left[E_{i} s^{2}\left(F_{j}\right)\left(\sum_{m} \Delta w_{m j} c_{m j} E_{m}\right)\right. \\
& \left.+F_{j} s^{2}\left(E_{i}\right)\left(\sum_{n} \Delta w_{i n} c_{i n} F_{n}\right)\right] . \tag{28}
\end{align*}
$$

To simplify (28) we make a diagonal approximation of the $M N \times M N$ matrix defining $\Delta w_{i j}$ by taking only constants and terms in $\Delta w_{i j}$. This is tantamount to assuming that all weights other than one individual $w_{i j}$
are correct. This approximation leads to

$$
\begin{align*}
& \Delta^{\prime} w_{i j}\left\{2-c_{i j} x_{i j}\left[u_{j}^{\prime} s^{2}\left(F_{j}\right) E_{i}+v_{i}^{\prime} s^{2}\left(E_{i}\right) F_{j}\right]\right\} \\
& =-w o_{i}+x_{i j}+x_{i j} \varphi_{o}-c_{i j} x_{i j}\left[u o_{j} s^{2}\left(F_{j}\right) E_{i}\right. \\
& \left.\quad+v o_{i} s^{2}\left(E_{i}\right) F_{j}\right] \sum_{k} \sum_{l} x_{k l}, \tag{29}
\end{align*}
$$

where

$$
u_{j}^{\prime}=\frac{\sum_{k} c_{k j} x_{k j} E_{k}}{\sum_{k} \sum_{l} x_{k l}}, \quad v_{i}^{\prime}=\frac{\sum_{l} c_{i l} x_{i l} F_{l}}{\sum_{k} \sum_{l} x_{k l}} .
$$

The second term in the braces \{\} on the left-hand side of (29) is small in comparison to 2 and may be neglected. It is also desirable to damp the shift $\Delta^{\prime} w_{l j}$ by a factor $\delta$ in a practical application. Hence we obtain

$$
\begin{align*}
w n_{i j}= & (1-\delta / 2) w o_{i j}+\delta / 2\left(x _ { i j } \left\{1+\varphi_{0}\right.\right. \\
& \left.\left.-c_{i j}\left[u o_{j} s^{2}\left(F_{j}\right) E_{i}+v o_{i} s^{2}\left(E_{i}\right) F_{j}\right] \sum_{k} \sum_{l} x_{k l}\right)\right) \cdot(3 \tag{30}
\end{align*}
$$

A suitable starting set of weights are the $x_{i j}$. In principle, it is the (unknown) population means $E_{i}$ and $F_{j}$ of the observations which enter into (30). In practice it is adequate to use the observed values $A_{i}$ and $B_{j}$ in place of $E_{l}$ and $F_{j}$.

## Application to polarization measurements

Details of the measurements and the technique are given in Vincent \& Flack (1980). Intensity measurements are available from the same crystal for both monochromated and filtered radiation. These intensity measurements are divided into two sets, which are mutually exclusive, based on values of $\theta$. We may write

$$
\begin{align*}
A_{i}= & {\left[\left(1+\cos ^{2} 2 \theta_{i}\right) / 2\right] I_{i} \text { (monochromator) } / I_{i} \text { (filter) }, } \\
& i=1,2, \ldots, M,  \tag{3}\\
B_{j}= & {\left[2 /\left(1+\cos ^{2} 2 \theta_{j}\right)\right] I_{j} \text { (filter) } / I_{j} \text { (monochromator) }, } \\
& j=1,2, \ldots, N,
\end{align*}
$$

where $I$ is the intensity measurement for a particular reflection under the stated conditions. From (3), $\gamma_{i j}=$ $A_{i} B_{j}$, the polarization state of the radiation may be written as

$$
\begin{equation*}
K_{i j}=\left(\cos ^{2} 2 \theta_{i}-\gamma_{i j} \cos ^{2} 2 \theta_{j}\right) /\left(\gamma_{l j}-1\right) . \tag{32}
\end{equation*}
$$

To apply the weighting scheme described above it is necessary to linearize (32). This is achieved by
expanding (32) to first order in a Taylor series about some 'guessed' value of $K, K_{g}$. One obtains the following values to be used in (2).

$$
\begin{align*}
c_{i j} & =\frac{\left(\cos ^{2} 2 \theta_{j}+K_{g}\right)^{2}}{\cos ^{2} 2 \theta_{j}-\cos ^{2} 2 \theta_{i}} \\
d_{i j} & =K_{g}-\frac{\left(\cos ^{2} 2 \theta_{i}+K_{g}\right)\left(\cos ^{2} 2 \theta_{j}+K_{g}\right)}{\cos ^{2} 2 \theta_{j}-\cos ^{2} 2 \theta_{i}} \tag{33}
\end{align*}
$$

The weighting scheme was implemented in a Fortran program. The weights were refined in cycles according to (30) until the maximum relative shift was less than a specified quantity (a value of $0.5 \%$ was chosen). It was found that $\delta$ had to be small $(0.03)$ to avoid divergence. A considerable number of cycles ( $\sim 650$ ) were necessary in order for the weights to converge but the calculation nevertheless is extremely rapid on a modern computer. We have not attempted to improve the rate of convergence. The value of $K_{g}$ initially chosen was for
a kinematical monochromator and the whole calculation was rerun with $K_{g}$ equal to $\hat{K}$ obtained from the first run. The value of $\hat{K}$ determined is highly insensitive to the value of $K_{g}$.

Some simulations were carried out to test the method and the program. These were based on the data used in paper II of this series but with intensities calculated from a model value of $K$, and with a pseudo-random Gaussian error added in to match the counting statistics of the experiment. From these simulations it is clear that the weighting scheme produces values of $\hat{K}$ closer to $K$ than the simple weighting scheme of $x_{i j}$ and that the e.s.d. of $\hat{K}$ is also reduced.

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# On the Theory and Estimation of the Three-Phase Structure Seminvariant in $\boldsymbol{P} \overline{\mathbf{1}}$ 

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#### Abstract

By embedding the three-phase structure seminvariant $T$ and its three symmetry-related variants in suitable quintets $Q$ (five-phase structure invariants) one obtains the extensions $Q$ of $T$ to which $T$ is related via the space-group symmetries. The neighborhoods of $T$ are then defined in terms of the neighborhoods of its extensions. The conditional probability distribution of $T$, given the seven magnitudes $|E|$ in its first neighborhood, is derived. The distribution yields a reliable estimate ( 0 or $\pi$ ) for $T$ in the favorable case that the variance of the distribution happens to be small.


## 1. Introduction

In recent years the basic concepts and mathematical formalism needed for the development of the probabilistic theory of the structure seminvariants have been elucidated. For example, it has long been known that, for fixed enantiomorph, the collection of observed magnitudes $|E|$ determines, in general, the values of all the structure seminvariants. A major recent advance is
the neighborhood principle: For fixed enantiomorph, the value of any structure seminvariant $T$ is primarily determined, in favorable cases, by the values of one or more small sets of magnitudes $|E|$, the neighborhoods of $T$, and is relatively insensitive to the values of the great bulk of remaining magnitudes (Hauptman, 1975). The conditional probability distribution of $T$, given the magnitudes in any of its neighborhoods, yields an estimate for $T$ which is particularly good in the favorable case that the variance of the distribution happens to be small. With the identification of systems of neighborhoods for the structure invariants, the probabilistic theory of the structure invariants developed rapidly, especially in space groups $P 1$ and $P \overline{1}$, but a great deal of work still remains to be done in deriving accurate and readily computable probability distributions, particularly for the space-group-special structure invariants.

Again, with the formulation of the extension concept [Hauptman (1977b, 1978); but see Giacovazzo (1977) for the equivalent concept called representation theory), the probabilistic theory of the structure seminvariants was reduced to that of the structure invariants, which is more highly developed. In particular, the

