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# Extinction in the Framework of Transfer Equations for General-Type Crystals 

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#### Abstract

An improvement of the classical theory of extinction in mosaic crystals is made by starting from the energy transfer equations valid for a general-type crystal according to Zachariasen's [Acta Cryst. (1967), 23, 558-564] classification. Within the assumption that only the integrated intensity of the diffraction peak is needed, the equations are first simplified and then solved. The result obtained for the extinction factor is similar to that of Becker \& Coppens [Acta Cryst. (1974), A30, 129-147], but two new parameters appear if the crystal is not of type I. One of them, determining the peculiarity of the transfer equations, gives differences in the extinction factor not greater than $8 \%$. The other, representing the ratio of the kinematical cross-section strengths along the diffracted and incident beams, gives differences up to $50 \%$. For crystals of ellipsoidal shape, empirical formulae appropriate for structure refinement programs are proposed.


## 1. Introduction

In this paper we re-analyse the problem of secondary extinction in the framework of classical transfer theory. The transfer equations for secondary extinction in finite crystals were first written by Hamilton (1957) and were based on the mosaic model of Darwin (1922) for the ideal imperfect crystal. Zachariasen (1967) has used similar equations to describe extinction in real crystals, so henceforth we will call these equations Hamilton-Zachariasen (HZ) equations.

[^0]Zachariasen stated that any real crystal is situated between two limiting types, distinguished by the nature of the peak width: type I if the width is given exclusively by the mosaic and type II if the width is given by the crystallite size only. Correspondingly, the secondary extinction follows the same classification. So far as primary extinction in small mosaic blocks is concerned, a description by the same transfer equations has been considered good enough under the assumption that this extinction is weak. The unified theory of Zachariasen has been very much criticized both for some mathematical errors and for its physical basis. On the same basis, Becker \& Coppens (1974a) (BC) have re-analysed the HZ equations. The solution which they provided has become very popular both for its convenient parametrization for least-squares-refinement programs and for its resistance to numerous experimental tests (see e.g. Hutton, Nelmes \& Scheel, 1981).

The limitations on the classical theory of extinction in real crystals were clarified by the new dynamical statistical theory of Kato (1976a, b, 1979, 1980). Starting from the dynamical equations for a distorted crystal and assuming a homogeneous and isotropic distribution of the defects, Kato derived a system of energy transfer equations valid for extinction only if the coherence distance $t_{c}$ is smaller than the extinction distance $\Lambda=(n \lambda|F|)^{-1}$. Here $\lambda$ is the wavelength, $F$ the structure factor and $n$ the density of unit cells. The energy transfer equations of Kato are similar to but not identical with the HZ equations. The differences discussed in detail by Kato (1976b, 1979) are in the form and physical interpretation of the coupling constants. Analysing the equivalence between the two kinds of energy transfer equations, Becker (1977) concludes that the range of validity found by Kato
for his energy transfer equations is the same for the HZ equations. The Kato theory, giving the extinction for any degree of crystal perfection, is more extensive than the classical theory. For $t_{c}>\Lambda$ the Kato theory replaces the energy transfer equations by equations for the averaged wave functions. On the basis of Kato's energy transfer equations, Kawamura \& Kato (KK) (1983) have elaborated a practical formula for the secondary extinction in a cylinder and a sphere, valid for a Bragg angle smaller than $30^{\circ}$ and extinction parameter smaller than 2 . Comparing their formula numerically with the BC calculations, KK have found significant differences. But recently Harada, Miyatake \& Sakata (1984) have compared KK and BC formulae experimentally using neutron diffraction data. Although some reflections were severely affected by extinction, they have found only small differences in the refined structure parameters.

It seems that in spite of its limitations there are not enough arguments for abandoning the classical theory of secondary extinction. But as was emphasized for a long time (Werner, 1974) the HZ equations describe the secondary extinction rigorously only in crystals of type I. When the natural broadening becomes competitive with the mosaic broadening or dominant, the classical energy transfer equations take a different form, the HZ equations being only a limiting case. These more general equations, labelled below as (1), can be directly obtained (see e.g. Popa, 1976) from the neutron transport equation of Vineyard (1954), which is a classical equation (see e.g. Sears, 1975, 1978). Thus the equations (1) are also classical, and like their HZ limits they describe the extinction only for crystals and diffraction maxima which fulfil the condition given above: $t_{c} \leqslant \Lambda$. This condition is fulfilled by type II rather than by type I crystals. Indeed the type II extinction is associated with a small block size, while in type I crystals the block size and the primary extinction may be large. In this paper we solve (1) in order to find a formula for the secondary extinction (alternative to BC) valid (in the limit $t_{c} \leqslant \Lambda$ ) for any type of crystals: type I, type II and a mixed type. The correction to the BC result is small for isotropic crystals of type II, but becomes significant for anisotropic crystals. This may explain why in many cases the BC formula fits the experimental data well, but sometimes, although the mosaic is small, the type II model for extinction gives an odd result (see Hutton, Nelmes \& Scheel, 1981).

This paper contains six parts. In § 2 the arguments permitting simplification of the transfer equations which are then solved in the next two parts are discussed. Two new quantities are defined, which differentiate our result from the BC result. In § 5 an empirical formula is proposed for extinction in spherical and ellipsoidal crystals with anisotropy. A comparison with the BC result is made in § 6.

## 2. Transfer equations for intensities

For small primary extinction and weak absorption the general form of the energy transfer equations is (Werner, 1974)

$$
\begin{align*}
\partial I_{1}\left(\mathbf{r}, \mathbf{k}_{1}\right) / \partial x_{1}= & -\left[\mu+\int \mathrm{d} \mathbf{k}_{2} \bar{\sigma}\left(\mathbf{k}_{1} \rightarrow \mathbf{k}_{2}\right)\right] I_{1}\left(\mathbf{r}, \mathbf{k}_{1}\right) \\
& +\int \mathrm{d} \mathbf{k}_{2} \bar{\sigma}\left(\mathbf{k}_{2} \rightarrow \mathbf{k}_{1}\right) I_{2}\left(\mathbf{r}, \mathbf{k}_{2}\right)  \tag{1a}\\
\partial I_{2}\left(\mathbf{r}, \mathbf{k}_{2}\right) / \partial x_{2}= & -\left[\mu+\int \mathrm{d} \mathbf{k}_{1} \bar{\sigma}\left(\mathbf{k}_{2} \rightarrow \mathbf{k}_{1}\right)\right] I_{2}\left(\mathbf{r}, \mathbf{k}_{2}\right) \\
& +\int \mathrm{d} \mathbf{k}_{1} \bar{\sigma}\left(\mathbf{k}_{1} \rightarrow \mathbf{k}_{2}\right) I_{1}\left(\mathbf{r}, \mathbf{k}_{1}\right) \tag{1b}
\end{align*}
$$

The indices 1 and 2 refer to the incident and diffracted beams respectively, $I$ is the intensity, $\mathbf{k}$ is the wave vector, $\mu$ is the linear absorption coefficient (including all non-Bragg scattering). $x_{1}, x_{2}$ are the coordinates in the oblique system with axes $\mathbf{i}_{1}, \mathbf{i}_{2}$ along the mean incident and diffracted beam respectively and $\mathbf{i}_{3}$ normal to the ( $\mathbf{i}_{1}, \mathbf{i}_{2}$ ) plane, which in the following will be considered horizontal (see Fig. 1). The quantity $\bar{\sigma}\left(\mathbf{k}_{1} \rightarrow \mathbf{k}_{2}\right)$ is the average over the mosaic distribution of the kinematical Bragg cross section per unit volume for the process $\mathbf{k}_{1} \rightarrow \mathbf{k}_{2}$. It is independent of $r$ if the crystal is homogeneous. If the primary extinction factor $y_{p}$ is not unity and the average size of the perfect crystallites is smaller than the extinction distance, then the same equations (1) will be used with replacement of $\bar{\sigma}$ by $y_{p} \bar{\sigma}$.

The explicit form of the cross section $\bar{\sigma}$ depends on the model chosen for the crystal microstructure: the shape and the dimension of the perfect blocks and the distribution of their orientations. The model frequently used for the anisotropic mosaic crystals has at most twelve parameters. These are the lengths and the orientations of the principal axes of two ellipsoids; one ellipsoid approximates the average shape of the perfect mosaic block (Coppens \& Hamilton, 1970) and the second is the surface of constant probability associated with the three-dimensional distribution function $W(\mathbf{\Delta})$ describing the block misorientations (Nelmes, 1980). Here $\Delta$ is a vector of which components represent small rotations around


Fig. 1. Diagram of the diffraction process and the coordinate systems (i), ( $\boldsymbol{\tau}$ ), ( $\mathbf{n}$ ). The wave vectors $\mathbf{k}_{1}, \mathbf{k}_{2}$ are drawn for convenience in the plane of their averages, $\mathbf{k}_{10}$ and $\mathbf{k}_{20}$, respectively. The unit vectors $i_{3}=\tau_{3}=n_{3}$ are perpendicular to the plane of the figure.
the axes of an orthogonal coordinate system. The number of model parameters reduces to two if no anisotropy exists, and to one if, moreover, the crystal is of type I or type II.

The Bragg cross section is elastic. It contains the factor $\delta\left(\mathbf{k}_{1}-\mathbf{k}_{2}\right)$ (which in the following will be omitted); thus all the integrals in (1) become double integrals.
If the profile of the diffraction peak is desired, (1) must be solved under the boundary conditions imposed by the instrument. But both the boundary conditions and (1) become simpler if one wants only the integrated intensity of the diffraction peak, a quantity independent of resolution. There are many possible ways to record the integrated intensity. If, for example, the diffractometer is set to have a very good resolution, a three-dimensional scan must be performed. The dimension of the scan can be decreased if the resolution is reduced. For example, if the collimator in front of the detector is removed and the detector's window is large enough, the integrated intensity is provided by a one-dimensional scan. For the angular dispersive diffraction method this may be a scan with the crystal (detector fixed) or a scan with the crystal and detector in the ratio $1: 2$. If a strongly divergent monochromatic beam with uniform angular distribution of intensity is available, the integrated intensity can be measured with both crystal and detector fixed (no scan). For the integrated intensity all these procedures are equivalent, the practical choice being dictated by other considerations. The last procedure entailing the simplest boundary conditions for (1) is the most convenient for our aim. In the energy dispersive method the integrated intensity measurement without scan may be realized if the channel width of the energy analyser is chosen large enough to cover the entire diffraction peak. In this method the incident beam is well collimated, but it has uniform wavelength distribution of intensity. In the following the angular dispersive method and the energy dispersive method will be treated together, as they are characterized by the same extinction factor (Tomiyoshi, Yamada \& Watanabe, 1980).
The Bragg cross section for a given block does not depend on $\mathbf{k}_{1}$ and $\mathbf{k}_{2}$ separately; it depends on their difference-more exactly, on the vector $h=$ $\mathbf{k}_{2}-\mathbf{k}_{1}-\mathbf{H}$, where $\mathbf{H}$ is a reciprocal vector for this block. This cross section is given by

$$
\begin{equation*}
\sigma(\mathbf{h})=n^{2}|F|^{2}\left|\int_{v} \exp (i \mathbf{h r}) \mathrm{d} v\right|^{2}, \tag{2}
\end{equation*}
$$

$v$ being the volume of the perfect block. Let us denote by $\mathbf{H}_{0}$ the vector $\mathbf{H}$ for the most probable mosaic block. This vector together with the line connecting the centres of the sample and detector define the diffraction plane, considered horizontal. The wave vectors $\mathbf{k}_{10}$ and $\mathbf{k}_{20}$ lying in the horizontal plane and
fulfilling exactly the Bragg condition $\mathbf{k}_{20}-\mathbf{k}_{10}=\mathbf{H}_{0}$ are the averages of the vectors $\mathbf{k}_{1}$ and $\mathbf{k}_{2}$ respectively. The vector $h$ is determined only by the deviations from the averages of the vectors $\mathbf{k}_{1}, \mathbf{k}_{2}, \mathbf{H}$. These deviations can be easily written, if two new coordinate systems (see Fig. 1) are introduced: the oblique system $\left(\tau_{i}\right)$ with $\tau_{1}, \tau_{2}$ lying in the horizontal plane and perpendicular to $i_{1}, i_{2}$ respectively, $\tau_{3}=i_{3}$; and the orthogonal system ( $\mathbf{n}_{i}$ ) with $\mathbf{n}_{2}$ along $\mathbf{H}_{0}, \mathbf{n}_{1}$ perpendicular to $\mathbf{H}_{0}$ in the horizontal plane and $\mathbf{n}_{3}=\mathbf{i}_{3}$. Then

$$
\begin{align*}
& \Delta \mathbf{k}_{l}=\Delta k \mathbf{i}_{l}+k_{0} \gamma_{l h} \tau_{l}+k_{0} \gamma_{l v} \tau_{3}, \quad(l=1,2),  \tag{3a}\\
& \Delta \mathbf{H}=\Delta \times \mathbf{H}_{0}=2 k_{0} \sin \theta\left(-\varepsilon_{3} \mathbf{n}_{1}+\varepsilon_{1} \mathbf{n}_{3}\right) . \tag{3b}
\end{align*}
$$

Here $\theta$ is the Bragg angle, $\gamma_{l h}, \gamma_{l v}(l=1,2)$ are the horizontal ( $h$ ) and vertical (v) divergence angles, $k_{0}=\left|\mathbf{k}_{10}\right|=\left|\mathbf{k}_{20}\right|=H /(2 \sin \theta), \quad \Delta k=k_{1}-k_{0}=k_{2}-k_{0}$ and $\varepsilon_{i}$ are the components of $\Delta$ in the $\left(\mathbf{n}_{i}\right)$ system. Writing all the vectors in the ( $\boldsymbol{\tau}_{i}$ ) system, one obtains the vector $h$ as follows:
$\mathbf{h}=k_{0}\left[-\left(\Gamma_{1}-\varepsilon_{3}\right) \boldsymbol{\tau}_{1}+\left(\Gamma_{2}-\varepsilon_{3}\right) \boldsymbol{\tau}_{2}+\left(\Gamma_{3}-2 \sin \theta \varepsilon_{1}\right) \boldsymbol{\tau}_{3}\right]$
where $\Gamma_{i}(i=1,2,3)$ are defined as follows:

$$
\begin{align*}
& \Gamma_{1}=\gamma_{1 h}-\Delta k \tan \theta / k_{0}, \\
& \Gamma_{2}=\gamma_{2 h}+\Delta k \tan \theta / k_{0},  \tag{5}\\
& \Gamma_{=}=\gamma_{2 v}-\gamma_{1 v} .
\end{align*}
$$

The quantities $\Gamma_{i}$, chosen to give a unique treatment for both diffraction methods, are equivalent to divergence angles, though their nature is not purely angular (except for $\Gamma_{3}$ ). By convention we call them equivalent divergence angles. The average cross sections $\bar{\sigma}\left(\mathbf{k}_{1} \rightarrow \mathbf{k}_{2}\right)=\bar{\sigma}\left(\mathbf{k}_{2} \rightarrow \mathbf{k}_{1}\right)$ are then

$$
\begin{align*}
\bar{\sigma}\left(\Gamma_{1}, \Gamma_{2}, \Gamma_{3}\right)= & \iiint_{-\infty}^{\infty} \sigma\left(\Gamma_{1}-\varepsilon_{3}, \Gamma_{2}-\varepsilon_{3}, \Gamma_{3}-2 \varepsilon_{1} \sin \theta\right) \\
& \times W\left(\varepsilon_{1}, \varepsilon_{2}, \varepsilon_{3}\right) \mathrm{d} \varepsilon_{1} \mathrm{~d} \varepsilon_{2} \mathrm{~d} \varepsilon_{3} . \tag{6}
\end{align*}
$$

The integrated cross sections $\int \mathrm{d} \mathbf{k}_{2} \bar{\sigma}\left(\mathbf{k}_{1} \rightarrow \mathbf{k}_{2}\right)$ and $\int \mathrm{d} \mathbf{k}_{1} \bar{\sigma}\left(\mathbf{k}_{2} \rightarrow \mathbf{k}_{1}\right)$ are in general different from each other and depend only on one equivalent divergence angle, $\Gamma_{1}$ and $\Gamma_{2}$ respectively. We denote them by $\bar{\sigma}_{1}\left(\Gamma_{1}\right)$ and $\bar{\sigma}_{2}\left(\Gamma_{2}\right)$.

Let us denote by $P_{i}$ the integral of $I_{i}$ over the vertical divergence angle $\gamma_{i v}(i=1,2)$. If $(1 a)$ is integrated over $\gamma_{1 v}$ in the range $(-\infty, \infty)$, then the left-hand side and the first term of the right-hand side become $\partial P_{1} / \partial x_{1}$ and $-\left[\mu+\bar{\sigma}_{1}\left(\Gamma_{1}\right)\right] P_{1}$. In the second term of the right-hand side this integration acts only on the cross section $\overline{\boldsymbol{\sigma}}\left(\Gamma_{1}, \Gamma_{2}, \Gamma_{3}\right)$ resulting in $\overline{\boldsymbol{\sigma}}\left(\Gamma_{1}, \Gamma_{2}\right)$, a function independent of $\gamma_{2 v}$; consequently, the existing integral over $\gamma_{2 v}$ acts only on $I_{2}$ and the second term becomes $\int_{-\infty}^{\infty} \mathrm{d} \gamma_{2 h} \bar{\sigma}\left(\Gamma_{1}, \Gamma_{2}\right) P_{2}$. Evidently, here $\mathrm{d} \gamma_{2 h}$ can be replaced by $\mathrm{d} \Gamma_{2}$. In the same manner we proceed with ( $1 b$ ), inverting the roles of $\gamma_{1 v}$ and $\gamma_{2 v}$. Thus we obtain two equations for the functions $P_{1}$,
$P_{2}$ similar to (1) except for the simple angular integral for the second term on the right-hand side. The cross section under this integral has the form

$$
\begin{equation*}
\bar{\sigma}\left(\Gamma_{1}, \Gamma_{2}\right)=\int_{-\infty}^{\infty} \mathrm{d} \varepsilon_{3} \sigma\left(\Gamma_{1}-\varepsilon_{3}, \Gamma_{2}-\varepsilon_{3}\right) W\left(\varepsilon_{3}\right) \tag{7a}
\end{equation*}
$$

where

$$
\begin{align*}
\sigma\left(\Gamma_{1}, \Gamma_{2}\right) & =\int_{-\infty}^{\infty} \mathrm{d} \Gamma_{3} \sigma\left(\Gamma_{1}, \Gamma_{2}, \Gamma_{3}\right)  \tag{7b}\\
W\left(\varepsilon_{3}\right) & =\int_{-\infty}^{\infty} \mathrm{d} \varepsilon_{1} \mathrm{~d} \varepsilon_{2} W\left(\varepsilon_{1}, \varepsilon_{2}, \varepsilon_{3}\right) . \tag{7c}
\end{align*}
$$

The formulae (7) are obtained by integrating (6) over $\Gamma_{3}$, an operation imposed by the integration of (1a) and ( $1 b$ ) over $\gamma_{1 v}$ and $\gamma_{2 v}$ respectively. The function $W\left(\varepsilon_{3}\right)$ is the one-dimensional mosaic distribution seen in a particular diffraction process and it is the projection on the axis $n_{3}$ of the three-dimensional mosaic distribution $W(\Delta)$ (Nelmes, 1980). The integral cross sections $\bar{\sigma}_{1}\left(\Gamma_{1}\right)$ and $\bar{\sigma}_{2}\left(\Gamma_{2}\right)$ are found by integrating ( $7 a$ ) over $\Gamma_{2}$ and $\Gamma_{1}$ respectively. Formally, $\bar{\sigma}\left(\Gamma_{1}, \Gamma_{2}\right)$ can be factorized as follows:

$$
\begin{equation*}
\bar{\sigma}\left(\Gamma_{1}, \Gamma_{2}\right)=\bar{\sigma}_{1}\left(\Gamma_{1}\right) Z_{1}\left(\Gamma_{1}, \Gamma_{2}\right)=\bar{\sigma}_{2}\left(\Gamma_{2}\right) Z_{2}\left(\Gamma_{2}, \Gamma_{1}\right) . \tag{8}
\end{equation*}
$$

Hence, the classical energy transfer equations allowing one to obtain the secondary extinction factor both for angular dispersive and energy dispersive diffraction methods are the following:

$$
\begin{align*}
\partial P_{1}\left(\mathbf{r}, \Gamma_{1}\right) / \partial x_{1}= & -\left[\mu+\bar{\sigma}_{1}\left(\Gamma_{1}\right)\right] P_{1}\left(\mathbf{r}, \Gamma_{1}\right) \\
& +\bar{\sigma}_{1}\left(\Gamma_{1}\right) \int_{-\infty}^{\infty} Z_{1}\left(\Gamma_{1} ; \Gamma_{2}\right) P_{2}\left(\mathbf{r}, \Gamma_{2}\right) \mathrm{d} \Gamma_{2} \tag{9a}
\end{align*}
$$

$$
\begin{align*}
\partial P_{2}\left(\mathrm{r}, \Gamma_{2}\right) / \partial x_{2}= & -\left[\mu+\bar{\sigma}_{2}\left(\Gamma_{2}\right)\right] P_{2}\left(\mathrm{r}, \Gamma_{2}\right) \\
& +\bar{\sigma}_{2}\left(\Gamma_{2}\right) \int_{-\infty}^{\infty} Z_{2}\left(\Gamma_{2}, \Gamma_{1}\right) P_{1}\left(\mathbf{r}, \Gamma_{1}\right) \mathrm{d} \Gamma_{1} . \tag{9b}
\end{align*}
$$

The boundary conditions for the functions $P_{1}$ and $P_{2}$ are as simple as for the functions $I_{1}$ and $I_{2}$. Thus, if the crystal has a convex shape and is bathed by an incident beam with uniform spatial distribution of intensity, the boundary conditions are (see Fig. 2)

$$
\left.P_{1}\left(\mathrm{r}, \Gamma_{1}\right)\right|_{A B C}=1 ;\left.\quad P_{2}\left(\mathbf{r}, \Gamma_{2}\right)\right|_{B C D}=0 . \quad(10 a, b)
$$

The equations (9) are different from HZ equations. Their right-hand sides contain angular integrals and different integrated cross sections $\bar{\sigma}_{i}\left(\Gamma_{i}\right)$. The existence of the integrals on the right-hand side of (9) is a direct consequence of the smallness of the
mean mosaic block. The problem was first discussed by Darwin (1922) and later by Werner (1974) and others (see e.g. Suorti, 1982). If the mosaic block is small, the size of the corresponding reflecting domain in reciprocal space is large. The effect of the mosaic misorientation is to extend the reflecting domain perpendicularly to the vector $\mathbf{H}_{0}$. Thus if a well collimated monochromatic beam is diffracted, the reflected beam becomes divergent. This divergence is determined by the section cut from the reflecting domain by the Ewald sphere. The successive reflexions make both incident and diffracted beams divergent; thus to write correctly the feedback term in the transfer equations we have to perform an integration over the divergence angle. If, moreover, the mean mosaic block has an anisotropic shape, the corresponding reflecting domain is ellipsoidal rather than spherical. Consequently the divergences of the incident and diffracted beams as well as the integrated cross sections $\bar{\sigma}_{1}$ and $\bar{\sigma}_{2}$ are different from each other (because the sections cut by the Ewald sphere are different). These effects become negligibly small for type I crystals. Indeed, if the mosaic width is much greater than the natural width, the reflecting domain is very flattened and the section cut by the Ewald sphere is small. Consequently (see $\S 4$ below) the functions $Z_{1}$ and $Z_{2}$ become $\delta\left(\Gamma_{1}-\Gamma_{2}\right), \bar{\sigma}_{1}$ and $\bar{\sigma}_{2}$ become identical and (9) are reduced to HZ equations.
3. The solution of the transfer equations and the general expression for the extinction factor
By definition, the extinction factor $y$ is the ratio between the diffracted integral intensity and its kinematical approximation. If the sample is totally immersed in the incident beam, then from Fig. 2 we have

$$
\begin{align*}
y= & {[\sin 2 \theta / Q V A(\mu)] } \\
& \times \int_{-\infty}^{\infty} \mathrm{d} \Gamma_{2} \int \mathrm{~d} x_{3} \int_{x_{1}(B)}^{x_{1}(D)} \mathrm{d} x_{1} P_{2}\left[x_{1}, x_{2}^{1}\left(x_{1}\right), x_{3}, \Gamma_{2}\right] . \tag{11}
\end{align*}
$$



Fig. 2. The cross section of the crystal in a plane parallel to the diffraction plane. The hatched area $S\left(x_{1}, x_{2}\right)$ is the integration area for equations (20), (28).

Here $V$ is the sample's volume and the quantities $Q$, $A(\mu)$ are

$$
\begin{gather*}
Q=n^{2}|F|^{2} \lambda^{3} / \sin 2 \theta  \tag{12}\\
A(\mu)=V^{-1} \int_{V} \mathrm{~d} V \exp \left[-\mu\left(t_{1}+t_{2}\right)\right] \tag{13}
\end{gather*}
$$

where $t_{1}, t_{2}$ represent the path in the sample of the incident and diffracted beams, respectively:

$$
\begin{align*}
& t_{1}\left(x_{1}, x_{2}\right)=x_{1}-x_{1}^{0}\left(x_{2}\right),  \tag{14a}\\
& t_{2}\left(x_{1}, x_{2}\right)=x_{2}^{1}\left(x_{1}\right)-x_{2} . \tag{14b}
\end{align*}
$$

As the problem is a plane one, in the following we will omit $x_{3}$. The diffracted intensity $P_{2}$ in (11) is found by solving (9) under the boundary conditions (10). This can be done for any form of the functions $\bar{\sigma}_{i}\left(\Gamma_{i}\right)$ and $Z_{i}\left(\Gamma_{i}, \Gamma_{j}\right)$. The strategy is the same as in Becker \& Coppens (1974a), but the calculation here is longer and for the sake of brevity some elementary but tiresome details are skipped.

In the first stage (9) and (20) are transformed into an integral equation. To do that let us define the functions $\psi_{i}$ and $\xi^{(i)}$ as follows:

$$
\begin{array}{r}
P_{i}\left(x_{1}, x_{2}, \Gamma_{i}\right)=\psi_{i}\left(x_{1}, x_{2}, \Gamma_{i}\right) \exp \left[-\mu\left(x_{1}+x_{2}\right)\right] \\
(i=1,2) ; \\
\xi^{(i)}\left(x_{1}, x_{2}, \Gamma_{i}\right)=\int_{-\infty}^{\infty} Z_{i}\left(\Gamma_{i}, \Gamma_{j}\right) \psi_{j}\left(x_{1}, x_{2}, \Gamma_{j}\right) \mathrm{d} \Gamma_{j} \\
(i \neq j=1,2) \tag{16}
\end{array}
$$

Then (9) and the boundary conditions (10) become

$$
\begin{align*}
\partial \psi_{i}\left(x_{1}, x_{2}, \Gamma_{i}\right) / \partial x_{i}= & -\bar{\sigma}_{i}\left(\Gamma_{i}\right) \psi_{i}\left(x_{1}, x_{2}, \Gamma_{i}\right) \\
& +\bar{\sigma}_{i}\left(\Gamma_{i}\right) \xi^{(i)}\left(x_{1}, x_{2}, \Gamma_{i}\right) \\
& \quad(i=1,2) ;  \tag{17}\\
\psi_{1}\left(x_{1}^{0}, x_{2}\right)= & \exp \left[\mu\left(x_{1}^{0}+x_{2}\right)\right]=g\left(x_{2}\right)  \tag{18a}\\
\psi_{2}\left(x_{1}, x_{2}^{0}\right)= & 0 \tag{18b}
\end{align*}
$$

The equations (17) under the boundary conditions (18) are solved as inhomogeneous differential equations and give the result

$$
\begin{align*}
\psi_{1}\left(x_{1}, x_{2}, \Gamma_{1}\right)= & g\left(x_{2}\right) \exp \left[-\bar{\sigma}_{1}\left(\Gamma_{1}\right) t_{1}\left(x_{1}, x_{2}\right)\right] \\
& +\bar{\sigma}_{1}\left(\Gamma_{1}\right) \int_{x_{1}^{0}}^{x_{1}} \exp \left[-\bar{\sigma}_{1}\left(\Gamma_{1}\right)\left(x_{1}-u_{1}\right)\right] \\
& \times \xi^{(1)}\left(u_{1}, x_{2}, \Gamma_{1}\right) \mathrm{d} u_{1},  \tag{19a}\\
\psi_{2}\left(x_{1}, x_{2}, \Gamma_{2}\right)= & \bar{\sigma}_{2}\left(\Gamma_{2}\right) \int_{x_{2}^{0}}^{x_{2}} \exp \left[-\bar{\sigma}_{2}\left(\Gamma_{2}\right)\left(x_{2}-u_{2}\right)\right] \\
& \times \xi^{(2)}\left(x_{1}, u_{2}, \Gamma_{2}\right) \mathrm{d} u_{2} . \tag{19b}
\end{align*}
$$

Now, multiplying (19a) by $Z_{2}\left(\Gamma_{2}, \Gamma_{1}\right)$, (19b) by $Z_{1}\left(\Gamma_{1}, \Gamma_{2}\right)$ and integrating over $\Gamma_{1}$ and $\Gamma_{2}$, respectively, one obtains two coupled integral equations for
the functions $\xi^{(i)}$. Eliminating $\xi^{(1)}$ one obtains

$$
\begin{align*}
\xi^{(2)}\left(x_{1}, x_{2}, \Gamma_{2}\right)= & g\left(x_{2}\right) \int_{-\infty}^{\infty} \mathrm{d} \Gamma_{1} Z_{2}\left(\Gamma_{1}, \Gamma_{1}\right) \\
& \times \exp \left[-\bar{\sigma}_{1}\left(\Gamma_{1}\right) t_{1}\left(x_{1}, x_{2}\right)\right] \\
& +\int_{-\infty}^{\infty} \mathrm{d} \Gamma_{1} Z_{2}\left(\Gamma_{2}, \Gamma_{1}\right) \bar{\sigma}_{1}\left(\Gamma_{1}\right) \\
& \times \int_{-\infty}^{\infty} \mathrm{d} \Lambda_{2} Z_{1}\left(\Gamma_{1}, \Lambda_{2}\right) \bar{\sigma}_{2}\left(\Lambda_{2}\right) \\
& \times \int_{x_{1}^{0}}^{x_{1}} \mathrm{~d} u_{1} \int_{x_{2}^{0}}^{x_{2}} \mathrm{~d} u_{2} \exp \left[-\bar{\sigma}_{1}\left(\Gamma_{1}\right)\left(x_{1}-u_{1}\right)\right. \\
& \left.-\bar{\sigma}_{2}\left(\Lambda_{2}\right)\left(x_{2}-u_{2}\right)\right] \xi^{(2)}\left(u_{1}, u_{2}, \Lambda_{2}\right) . \tag{20}
\end{align*}
$$

The extinction factor can be expressed by the function $\xi^{(2)}$; indeed, from (15) and (19b), (11) becomes

$$
\begin{align*}
y= & {[Q V A(\mu)]^{-1} \int_{-\infty}^{\infty} \mathrm{d} \Gamma_{2} \bar{\sigma}_{2}\left(\Gamma_{2}\right) \int \mathrm{d} x_{3} \int_{S_{0}\left(x_{3}\right)} \mathrm{d} S f\left(x_{1}\right) } \\
& \times \exp \left[-\bar{\sigma}_{2}\left(\Gamma_{2}\right) t_{2}\left(x_{1}, x_{2}\right)\right] \xi^{(2)}\left(x_{1}, x_{2}, \Gamma_{2}\right) \tag{21}
\end{align*}
$$

where we have denoted by $S_{0}\left(x_{3}\right)$ the crystal crosssectional area for a given $x_{3}$ and by $f\left(x_{1}\right)$ the function

$$
\begin{equation*}
f\left(x_{1}\right)=\exp \left[-\mu\left(x_{1}+x_{2}^{1}\right)\right] . \tag{22}
\end{equation*}
$$

In the second stage, the integral equation (20) is solved by successive approximations. For that, the following factorization for $\bar{\sigma}_{i}$ and series development of $\xi^{(2)}$ are necessary:

$$
\begin{gather*}
\bar{\sigma}_{i}\left(\Gamma_{i}\right)=Q G_{i}\left(\Gamma_{i}\right)  \tag{23}\\
\xi^{(2)}\left(x_{1}, x_{2}, \Gamma_{2}\right)=\sum_{n=0}^{\infty}\left[(-1)^{n} / n!\right] \xi_{n}\left(x_{1}, x_{2}, \Gamma_{2}\right) Q^{n} . \tag{24}
\end{gather*}
$$

Here the functions $G_{i}$ and $\xi_{n}$ do not depend on $Q$. Introducing (23), (24) into (20) and identifying the coefficients of $Q^{n}$, one obtains

$$
\begin{align*}
\xi_{n}\left(x_{1}, x_{2}, \Gamma_{2}\right)= & F_{n}\left(\Gamma_{2}\right) g\left(x_{2}\right) t_{1}^{n}\left(x_{1}, x_{2}\right) \\
& +n(n-1) \sum_{l=0}^{n-2} \sum_{m=0}^{l}\binom{n-2}{l}\binom{l}{m} \\
& \times \int_{-\infty}^{\infty} \mathrm{d} \Lambda_{2} H_{n-1-l}\left(\Gamma_{2}, \Lambda_{2}\right) \\
& \times G_{2}^{l-m+1}\left(\Lambda_{2}\right) \\
& \times \int_{S\left(x_{1}, x_{2}\right)} \mathrm{d} u_{1} \mathrm{~d} u_{2}\left(x_{1}-u_{1}\right)^{n-2-l} \\
& \times\left(x_{2}-u_{x}\right)^{l-m} \xi_{m}\left(u_{1}, u_{2}, \Lambda_{2}\right), \tag{25}
\end{align*}
$$

where we have denoted by $S\left(x_{1}, x_{2}\right)$ the hatched area in Fig. 2 and by $F_{n}$ and $H_{n+1}$ the following functions:

$$
\begin{align*}
F_{n}\left(\Gamma_{2}\right)= & \int_{-\infty}^{\infty} \mathrm{d} \Gamma_{1} Z_{2}\left(\Gamma_{2}, \Gamma_{1}\right) G_{1}^{n}\left(\Gamma_{1}\right) ; \quad(n \geq 0), \\
H_{n+1}\left(\Gamma_{2}, \Lambda_{2}\right)= & \int_{-\infty}^{\infty} \mathrm{d} \Gamma_{1} Z_{2}\left(\Gamma_{2}, \Gamma_{1}\right)  \tag{26}\\
& \times G_{1}^{n+1}\left(\Gamma_{1}\right) Z_{1}\left(\Gamma_{1}, \Lambda_{2}\right) ; \quad(n \geq 0) . \tag{27}
\end{align*}
$$

In (25) and below, the sum exists only if the upper limit is greater than or equal to the lower one.

The recurrence relation (25) can be solved using the BC approximation presented in Appendix $A$. Indeed, let us iterate (25) beginning with $n=0 . \xi_{0}$ and $\xi_{1}$ are simple functions [first term of (25)]; $\xi_{2}$ and $\xi_{3}$ contain an angular and a surface integral. But $\xi_{4}$ contains a double angular and a double surface integral ( $A 1$ ). The latter is reduced to a simple surface integral using the BC approximation (A4). If one continues with $\xi_{5}$ and $\xi_{6}$ in the same manner, enough information is obtained to suggest the following form for the general term $\xi_{n}$ :

$$
\begin{align*}
\xi_{n}\left(x_{1}, x_{2}, \Gamma_{2}\right)= & F_{n}\left(\Gamma_{2}\right) g\left(x_{2}\right) t_{1}^{n}\left(x_{1}, x_{2}\right) \\
& +n(n-1) \int_{S\left(x_{1}, x_{2}\right)} \mathrm{d} u_{1} \mathrm{~d} u_{2} g\left(u_{2}\right) \\
& \times \sum_{m=0}^{n-2} t_{1}^{m}\left(u_{1}, u_{2}\right) \\
& \times \sum_{l=m}^{n-2}\binom{n-2}{l}\binom{l}{m}\left(x_{1}-u_{1}\right)^{n-2-l} \\
& \times\left(x_{2}-u_{2}\right)^{l-m} K_{n-2, m, l}\left(\Gamma_{2}\right), \tag{28}
\end{align*}
$$

where the functions $K\left(\Gamma_{2}\right)$ fulfil the following recurrence relations ( $n \geq 2$ ):

$$
\begin{align*}
K_{n-2, m, l}\left(\Gamma_{2}\right)= & \int_{-\infty}^{\infty} \mathrm{d} \Lambda_{2} H_{n-1-l}\left(\Gamma_{2}, \Lambda_{2}\right) V_{l m}\left(\Lambda_{2}\right)  \tag{29a}\\
V_{l m}\left(\Gamma_{2}\right)= & G_{2}^{l-m+1}\left(\Gamma_{2}\right) F_{m}\left(\Gamma_{2}\right) \\
& +\sum_{j=m+1}^{l} G_{2}^{l-j+1}\left(\Gamma_{2}\right) \\
& \times \sum_{k=0}^{m-1} K_{j-2, k, j-m+k-1}\left(\Gamma_{2}\right) \tag{29b}
\end{align*}
$$

The correctness of (28) is verified in the last step of this complete induction procedure by replacing it in (25) and using once again the BC approximation (A4). Now, if (21) is developed in a power series of $Q$ and $\xi_{n}$ is replaced by (28), another double surface integral (A5) is obtained, which is reduced to a simple surface integral using the BC approximation (A6).

Finally one obtains

$$
\begin{align*}
y= & {[V A(\mu)]^{-1} \int_{-\infty}^{\infty} \mathrm{d} \Gamma_{2} \int_{V} \mathrm{~d} V \exp \left[-\mu\left(t_{1}+t_{2}\right)\right] } \\
& \times \sum_{n=0}^{\infty} \frac{(-1)^{n}}{n!} Q^{n} \sum_{m=0}^{n}\binom{n}{m} V_{n m}\left(\Gamma_{2}\right) t_{1}^{m} t_{2}^{n-m} \tag{30}
\end{align*}
$$

where the functions $V_{n m}\left(\Gamma_{2}\right)$ satisfy a recurrence relation obtained from (29):

$$
\begin{align*}
V_{n m}\left(\Gamma_{2}\right)= & G_{2}^{n-m+1}\left(\Gamma_{2}\right) F_{m}\left(\Gamma_{2}\right) \\
& +\sum_{j=m+1}^{n} G_{2}^{n-j+1}\left(\Gamma_{2}\right) \sum_{l=0}^{m-1} \int_{-\infty}^{\infty} \mathrm{d} \Lambda_{2} \\
& \times H_{l+1}\left(\Gamma_{2}, \Lambda_{2}\right) V_{j-2-l, m-l-1} . \tag{31}
\end{align*}
$$

The formula (30) is the most general expression for the extinction factor in the limits of the transfer theory and BC approximation, being valid for any model of the crystal microstructure. For practical purposes a closed formula must be obtained. This is a simple task if the crystal is of type $I$, but for the general type we need one more approximation. In the following the calculations will be much simplified if the mosaic model is introduced explicitly.

## 4. Extinction for definite mosaic distributions

Traditionally, the misorientation of the mosaic is described by the Gaussian distribution which in three dimensions is

$$
\begin{equation*}
W(\Delta)=g_{1} g_{2} g_{3} \exp \left(-\pi \sum_{i} g_{i}^{2} \Delta_{i}^{2}\right) . \tag{32}
\end{equation*}
$$

Here the variables $\Delta_{i}$ are small rotations around the principal axes of the constant probability ellipsoid and $g_{i}$ give the widths of the distribution at half height:

$$
w_{i}=2(\ln 2 / \pi)^{1 / 2} / g_{i}=0.939 / g_{i} .
$$

The principal axes are oriented along the unit vectors $\mathbf{m}_{i}$ and in order to write down the distribution in the $\left(\mathbf{n}_{i}\right)$ system an orthogonal matrix $\mathrm{E}^{(\boldsymbol{m})}$ connecting $\mathbf{m}_{i}$ and $\mathbf{n}_{i}$ must be introduced:

$$
\begin{equation*}
\mathrm{m}_{i}=\sum_{i} E_{i j}^{(m)} n_{j} \quad(i=1,2,3) . \tag{33}
\end{equation*}
$$

With

$$
\Delta=\sum_{i} \Delta_{i} \mathbf{m}_{i}=\sum_{i} \varepsilon_{i} \mathbf{n}_{i}
$$

one obtains

$$
\begin{gather*}
W\left(\varepsilon_{1}, \varepsilon_{2}, \varepsilon_{3}\right)=|\mathrm{C}|^{1 / 2} \exp \left(-\pi \sum_{j} \sum_{k} C_{j k} \varepsilon_{j} \varepsilon_{k}\right),  \tag{34}\\
C=\mathrm{E}^{(m)^{\prime}} \mathrm{G}^{2} \mathrm{E}^{(m)} ; \quad G_{i j}=g_{i} \delta_{i j}, \tag{35a,b}
\end{gather*}
$$

where $\delta_{i j}$ are Kronecker symbols. By integrating over $\varepsilon_{1}$ and $\varepsilon_{2}$ one obtains the one-dimensional mosaic distribution seen in the diffraction process:

$$
\begin{gather*}
W\left(\varepsilon_{3}\right)=g_{G} \exp \left(-\pi g_{G}^{2} \varepsilon_{3}^{2}\right)  \tag{36}\\
1 / g_{G}^{2}=\sum_{k} E_{k 3}^{(m) 2} / g_{k}^{2} \tag{37}
\end{gather*}
$$

which must be convoluted with the unit-volume cross section $\sigma\left(\Gamma_{1}, \Gamma_{2}\right)$ of the average perfect block.

The latter is calculated in Appendix $B(B 12, B 13)$ for an ellipsoid with principal axes $r_{i}$ oriented along the unit vectors $\mathbf{c}_{i}(i=1,2,3)$ and depends on the ellipsoid radii $\rho_{1}, \rho_{2}$ along the incident and diffracted directions respectively, and on the transformed Bragg angle $\theta^{\prime}$ (when the ellipsoid is transformed into a sphere). These quantities can be calculated if the orthogonal matrix $\mathrm{E}^{(c)}$ relating the systems ( $\mathbf{c}_{i}$ ) and $\left(\mathbf{n}_{i}\right)$ is introduced:

$$
\begin{equation*}
\mathbf{c}_{i}=\sum_{j} E_{i j}^{(c)} \mathbf{n}_{j} \quad(i=1,2,3) . \tag{38}
\end{equation*}
$$

Indeed, using the transformation (see Fig. 1)

$$
\begin{equation*}
\mathbf{i}_{l}=\cos \theta \mathbf{n}_{1}+(-1)^{l} \sin \theta \mathbf{n}_{2} \quad(l=1,2) \tag{39}
\end{equation*}
$$

and ( $C 1$ ), ( $C 2$ ) from Appendix $C$ we have

$$
\begin{array}{r}
1 / \rho_{l}^{2}=\sum_{k}\left[\cos \theta E_{k 1}^{(c)}+(-1)^{l} \sin \theta E_{k 2}^{(c)}\right]^{2} / r_{k}^{2} \\
(l=1,2) ; \tag{40}
\end{array}
$$

$\cos 2 \theta^{\prime}=\rho_{1} \rho_{2} \sum_{k}\left(\cos ^{2} \theta E_{k 1}^{(c) 2}-\sin ^{2} \theta E_{k 2}^{(c)}\right) / r_{k}^{2}$.
Now the convolution (7) can be performed; with the factorizations (8) and (23) one obtains

$$
\begin{align*}
G_{i}\left(\Gamma_{i}\right)= & \alpha_{i G} \exp \left(-\pi \alpha_{i G}^{2} \Gamma_{i}^{2}\right) \quad(i=1,2),  \tag{42}\\
Z_{i}\left(\Gamma_{i}, \Gamma_{j}\right)= & \left(\alpha_{i G} / \delta_{i G}\right) \exp \left[-\pi\left(\alpha_{i G}^{2} / \delta_{i G}^{2}\right)\right. \\
& \left.\times\left(\nu_{i G} \Gamma_{i}-\Gamma_{j}\right)^{2}\right] \quad(i \neq j=1,2), \tag{43}
\end{align*}
$$

where

$$
\begin{align*}
& \alpha_{i G}^{2}=\left(g_{G}^{2} c_{0}^{2} \rho_{j}^{2} \sin ^{2} 2 \theta / \lambda^{2}\right) /\left(g_{G}^{2}+c_{0}^{2} \rho_{j}^{2} \sin ^{2} 2 \theta / \lambda^{2}\right) \\
&(i \neq j=1,2),  \tag{44}\\
& \nu_{i G}=\left(g_{G}^{2} \cos 2 \theta^{\prime} \rho_{j} / \rho_{i}+c_{0}^{2} \rho_{j}^{2} \sin ^{2} 2 \theta / \lambda^{2}\right) \\
& \times\left(g_{G}^{2}+c_{0}^{2} \rho_{j}^{2} \sin ^{2} 2 \theta / \lambda^{2}\right)^{-1} \quad(i \neq j=1,2),  \tag{45}\\
& \delta_{1 G}^{2}=1 / a_{G}^{2}-\nu_{1 G}^{2}, \quad \delta_{2 G}^{2}=a_{G}^{2}-\nu_{2 G}^{2} .
\end{align*}
$$

Here $c_{0}=1.612$ and $a_{G}$ is the block-shape-anisotropy parameter defined as

$$
\begin{align*}
a_{G}= & \alpha_{2 G} / \alpha_{1 G} \\
= & \left(\rho_{1} / \rho_{2}\right)\left[\left(g_{G}^{2}+c_{0}^{2} \rho_{2}^{2} \sin ^{2} 2 \theta / \lambda^{2}\right)\right. \\
& \left.\times\left(g_{G}^{2}+c_{0}^{2} \rho_{1}^{2} \sin ^{2} 2 \theta / \lambda^{2}\right)^{-1}\right]^{1 / 2} . \tag{47}
\end{align*}
$$

The following supplementary relations hold:

$$
\begin{align*}
\nu_{2 G}=a_{G}^{2} \nu_{1 G} ; \quad \delta_{2 G} & =a_{G}^{2} \delta_{1 G}  \tag{48a,b}\\
\nu_{1 G} \nu_{2 G}+\delta_{1 G} \delta_{2 G} & =1 . \tag{48c}
\end{align*}
$$

Always, $0 \leq \nu_{1 G} \nu_{2 G} \leq 1$ and consequently $0 \leq \delta_{1 G} \delta_{2 G}$ $\leq 1$.

Let us now discuss the limiting situations. If $g_{G} \ll$ $c_{0} \sin 2 \theta \min \left(\rho_{1}, \rho_{2}\right) / \lambda$, then $a_{G}=1, \alpha_{1 G}=\alpha_{2 G}=g_{G}$, $\nu_{1 G}=\nu_{2 G}=1, \delta_{1 G}=\delta_{2 G}=0, \alpha_{i G} / \delta_{i G} \rightarrow \infty$ and as a consequence the functions $Z_{i}$ are $\delta$ functions. In this situation the crystal is of type $I$, the extinction is governed by the mosaic only and the transfer equations are HZ . If, on the contrary, $g_{G} \gg$ $c_{0} \sin 2 \theta \max \left(\rho_{1}, \rho_{2}\right) / \lambda$, then $a_{G}=\rho_{1} / \rho_{2} \neq 1, \alpha_{i G}=$ $c_{0} \rho_{j} \sin 2 \theta / \lambda, \quad \nu_{i G}=\rho_{j} \cos 2 \theta^{\prime} / \rho_{i}, \delta_{i G}=\rho_{j} \sin 2 \theta^{\prime} / \rho_{i}$, ( $i \neq j=1,2$ ); $\alpha_{i G} / \delta_{i G}$ are finite and $Z_{i}$ are not $\delta$ functions. In this case the crystal is of type II, the extinction is governed only by the block size and the transfer equations are no longer HZ , as in the intermediate situation where the crystal is of mixed type. If $\sin 2 \theta$ is very small, the crystal cannot be of pure type I; but in this case $\rho_{1} \simeq \rho_{2}$, hence $a_{G} \simeq 1, \nu_{1 G}=\nu_{2 G} \simeq 1$, $\delta_{1 G}=\delta_{2 G} \simeq 0$ and although the crystal is of the mixed type the transfer equations are HZ. It is seen that when the quantity $\left(\delta_{1} \delta_{2}\right)_{G}$ is zero the transfer equations are always $H Z$ and the extinction is correctly given by the BC theory. This parameter is directly obtained from (45) and (48c):

$$
\begin{align*}
\left(\delta_{1} \delta_{2}\right)_{G}= & g_{G}^{2}\left[g_{G}^{2} \sin ^{2} 2 \theta^{\prime}+c_{0}^{2} \sin ^{2} 2 \theta\right. \\
& \left.\times\left(\rho_{1}^{2}+\rho_{2}^{2}-2 \rho_{1} \rho_{2} \cos 2 \theta^{\prime}\right) / \lambda^{2}\right] \\
& \times\left[\left(g_{G}^{2}+c_{0}^{2} \sin ^{2} 2 \theta \rho_{1}^{2} / \lambda^{2}\right)\right. \\
& \left.\times\left(g_{G}^{2}+c_{0}^{2} \sin ^{2} 2 \theta \rho_{2}^{2} / \lambda^{2}\right)\right]^{-1} . \tag{49}
\end{align*}
$$

Once given $G_{i}$ and $Z_{i}$, it is possible to calculate the functions $V_{n m}\left(\Gamma_{2}\right)$. If (31) is iterated beginning with $n=0$, after a few steps it will be observed that every $V_{n m}\left(\Gamma_{2}\right)$ is a sum of $\binom{n}{m}$ Gaussians, all having the common factor $\alpha_{1 G}^{m} \alpha_{2 G}^{n-m+1}$. For type I crystals all the terms are identical and

$$
\begin{align*}
V_{n m}\left(\Gamma_{2}\right) & =\binom{n}{m} W^{n+1}\left(\Gamma_{2}\right) \\
& =\binom{n}{m} g_{G}^{n+1} \exp \left[-\pi(n+1) g_{G}^{2} \Gamma_{2}^{2}\right] . \tag{50}
\end{align*}
$$

For the general-type crystal the expected difficulties in finding exact $V_{n m}\left(\Gamma_{2}\right)$ are discouraging. On the other hand, we do not need the profile but the integral of $V_{n m}\left(\Gamma_{2}\right)$, and so we accept the following approximation suggested by (50):

$$
\begin{align*}
V_{n m}\left(\Gamma_{2}\right)= & \binom{n}{m} \alpha_{1 G}^{m} \alpha_{2 G}^{n-m+1} p_{n m} \\
& \times \exp \left[-\pi(n+1) \alpha_{2 G}^{2}\left(p_{n m} / v_{n m}\right)^{2} \Gamma_{2}^{2}\right] . \tag{51}
\end{align*}
$$

Introducing this into (31) and retaining only (a) the height and area or (b) the height and the second moment, we obtain two coupled recurrence relations for the coefficients $p_{n m}$ and $v_{n m}$. These relations depend on one parameter only, $\left(\delta_{1} \delta_{2}\right)_{G}$, and can be solved numerically. The quantities of interest [which determine the areas of $\left.V_{n m}\left(\Gamma_{2}\right)\right]$ are $v_{n m}$. They fulfil the conditions $v_{n m}=v_{n, n-m}$ and $v_{n m} \leq 1$, the equality taking place for $m=0$ and $m=n$ (and of course always for $\delta_{1} \delta_{2}=0$ ). The recurrence relations (a) overestimate $v_{n m}$ and, conversely, (b) underestimate them. Some coefficients $v_{n m}$ calculated for $\left(\delta_{1} \delta_{2}\right)_{G}=1$ are represented in Fig. 3.

With the approximation (51) the extinction factor (30) becomes

$$
\begin{align*}
& y=\sum_{n=0}^{\infty}\left[(-1)^{n} / n!\right]\left[\left(Q \alpha_{1 G}\right)^{n} /(n+1)^{1 / 2}\right] \overline{t^{(n)}},  \tag{52a}\\
& \overline{t^{(n)}}= {[V A(\mu)]^{-1} \int_{V} \mathrm{~d} V \exp \left[-\mu\left(t_{1}+t_{2}\right)\right] } \\
& \times \sum_{m=0}^{n}\binom{n}{m}^{2} v_{n m} t_{1}^{m}\left(a_{G} t_{2}\right)^{n-m} . \tag{52b}
\end{align*}
$$

This differs from the corresponding Becker \& Coppens (1974a) result [their equation (51)] by the coefficients $v_{n m}$ and $a_{G}$ in $t^{(n)}$. The coefficient $a_{G}$ occurs in (52) for $n \geq 1$, then it occurs also in the main extinction parameter. As in BC this parameter is defined by (only the factor $2 / 3$ is ignored)

$$
\begin{equation*}
x=Q \alpha_{1 G} \overline{t^{(1)}} \tag{53}
\end{equation*}
$$

The quantity $\overline{t^{(1)}}$ here is the mean path in the crystal [or in the mosaic block if $V$ in ( $52 b$ ) is replaced by $v$ ] only for $a_{G}=1$.

The coefficients $v_{n m}$ occur in (52) only for $n \geq 2$ and they give a relatively small correction for the BC formula compared with $a_{G}$. To evaluate this correction we have calculated $y$ by direct summation of (52) for a crystal of parallelepipedic shape with the edges along $x_{1}, x_{2}, x_{3}, \mu=0$ and $a_{G}=1$. In this case the volume integral in ( $52 b$ ) is readily performed. $v_{n m}$ were calculated numerically for $\delta_{1} \delta_{2}=0,0 \cdot 4,0 \cdot 7,1 \cdot 0$


Fig. 3. Some of the coefficients $v_{n m}$ calculated for $\delta_{1} \delta_{2}=1 ;+v_{n m}$ overestimated; $\bigcirc v_{n m}$ underestimated.
and $n \leq 60$. For these $n$ it is possible to sum (52a) with an accuracy of three decimal digits in the range $0 \leq x \leq 5.7(y \geq 0 \cdot 2)$. The results of the summations made with underestimated and overestimated $v_{n m}$ differ from each other by at most $1 \cdot 1 \%$, therefore their mean value was considered satisfactory. The values of $y$ calculated in this way for values of $x$ varied with a step of $0 \cdot 1$ were used to fit the following empirical function:
$y\left(x, \delta_{1} \delta_{2}\right)=y_{0}(x)\left[1-\delta_{1} \delta_{2} A_{0} x^{C_{0}} /\left(1+B_{0} x^{C_{0}}\right)\right]$,
where $y_{0}(x)=y(x, 0)$. We have obtained $A_{0}=0 \cdot 02$, $B_{0}=0.26, C_{0}=1.5$ with

$$
R=\sum \mid y(\text { calc. })-y\left(x, \delta_{1} \delta_{2}\right) \mid / \sum y(\text { calc. })=0.008 .
$$

For $x \leq 30$ the maximum deviation of $y(x, 1)$ from $y_{0}(x)$ is only $8 \%$. This explains why the BC formula gives good results in many cases.

For the crystal with $\mu \neq 0, a \neq 1$ and of other shape or mosaic distribution, small variations of the parameters $A_{0}, B_{0}$ and $C_{0}$ are expected. But if the smallness of the correction (54) itself is taken into account, these variations can be neglected. If (54) is in general accepted, it only remains to express $y_{0}$. For that we must take $v_{n m}=1$ in (52b). But in this case, by using the identity

$$
\begin{align*}
& \sum_{n=0}^{\infty} \frac{(-1)^{n}}{n!} c^{n} \sum_{m=0}^{n}\binom{n}{m}^{2} u^{m} v^{n-m} \\
& \quad=\exp [-c(u+v)] I_{0}\left[2 c(u v)^{1 / 2}\right] \tag{55}
\end{align*}
$$

where $I_{0}$ is the modified Bessel function, one obtains $y_{0}$ as follows:

$$
\begin{align*}
y_{0}= & {[Q V A(\mu)]^{-1} \int_{V} \mathrm{~d} V \exp \left[-\mu\left(t_{1}+t_{2}\right)\right] } \\
& \times \int_{-\infty}^{\infty} \mathrm{d} \gamma \bar{\sigma}_{1}(\gamma) \exp \left[-\bar{\sigma}_{1}(\gamma)\left(t_{1}+a t_{2}\right)\right] \\
& \times I_{0}\left[2 \bar{\sigma}_{1}(\gamma)\left(a t_{1} t_{2}\right)^{1 / 2}\right] . \tag{56}
\end{align*}
$$

This expression is valid for all mosaic distributions leading to the functions $V_{n m}\left(\Gamma_{2}\right)$, in which the quantity $\alpha_{1}^{m} \alpha_{2}^{n-m+1}$ can be factorized.

It seems (e.g. Becker \& Coppens, 1974b) that the Lorentz mosaic distribution is more adequate for the secondary extinction than the Gauss distribution. But it must be noted that a three-dimensional Lorentz function has an infinite norm and cannot be a physical distribution. A function with tails longer than Gaussian which may be a mosaic distribution is

$$
\begin{equation*}
W(\Delta)=2 \pi g_{1} g_{2} g_{3} /\left[1+\beta_{0}\left(\sum_{i} g_{i}^{2} \Delta_{i}^{2}\right)^{m}\right], \quad(m \geq 2) \tag{57}
\end{equation*}
$$

where $\beta_{0}(m)$ results from the normalization condition. Taking $m=2$ we obtain $\beta_{0}=4 \pi^{4}$ and $w_{i}=$ $2^{1 / 2} /\left(\pi g_{i}\right)=0 \cdot 45 / g_{i}$. The one-dimensional mosaic
distribution becomes
$W\left(\varepsilon_{3}\right)=\left(\pi g_{L}^{\prime} / 2\right)\left[1-(2 / \pi) \arctan 2 \pi^{2} g_{L}^{\prime 2} \varepsilon_{3}^{2}\right]$,
where for $g_{L}^{\prime}$ one has exactly the same expression as for $g_{G}$. The asymptotic behaviour is $\varepsilon_{3}^{-2}$, when it can be well approximated by the Lorentz function:

$$
\begin{equation*}
W\left(\varepsilon_{3}\right)=g_{L} /\left(1+\pi^{2} g_{L}^{2} \varepsilon_{3}^{2}\right) \quad g_{L}=\pi g_{L}^{\prime} / 2 \tag{59}
\end{equation*}
$$

This is the sole argument for preserving the name Lorentz for this distribution. Further, to calculate the functions $G_{i}\left(\Gamma_{i}\right)$ it is enough [because (54) is accepted] to convolute (59) with the perfect block cross sections $\sigma_{i}\left(\Gamma_{i}\right)$ given by (B14). Their Lorentzian approximations ( $B 16$ ) are preferable for reasons of simplicity and the following results:

$$
\begin{gather*}
G_{i}\left(\Gamma_{i}\right)=\alpha_{i L} /\left(1+\pi^{2} \alpha_{i L}^{2} \Gamma_{i}^{2}\right) \quad(i=1,2) ;  \tag{60}\\
\alpha_{i L}=\left(1 \cdot 5 g_{L} \rho_{j} \sin 2 \theta / \lambda\right) /\left(g_{L}+1 \cdot 5 \rho_{j} \sin 2 \theta / \lambda\right) \\
(i \neq j=1,2) . \tag{61}
\end{gather*}
$$

The block anisotropy parameter $a_{L}$ becomes

$$
\begin{align*}
a_{L}= & \alpha_{2 L} / \alpha_{1 L} \\
= & \left(\rho_{1} / \rho_{2}\right)\left(g_{L}+1 \cdot 5 \rho_{2} \sin 2 \theta / \lambda\right) \\
& \times\left(g_{L}+1 \cdot 5 \rho_{1} \sin 2 \theta / \lambda\right)^{-1} \tag{62}
\end{align*}
$$

and, in place of (49),

$$
\begin{align*}
\left(\delta_{1} \delta_{2}\right)_{L}^{1 / 2}= & g_{L}\left[g_{L} \sin 2 \theta^{\prime}+1 \cdot 5 \sin 2 \theta\left(\rho_{1}^{2}+\rho_{2}^{2}\right.\right. \\
& \left.\left.-2 \rho_{1} \rho_{2}-2 \rho_{1} \rho_{2} \cos 2 \theta^{\prime}\right)^{1 / 2} / \lambda\right] \\
& \times\left[\left(g_{L}+1 \cdot 5 \rho_{1} \sin 2 \theta / \lambda\right)\right. \\
& \left.\times\left(g_{L}+1 \cdot 5 \rho_{2} \sin 2 \theta / \lambda\right)\right]^{-1} \tag{63}
\end{align*}
$$

is used.
It can be anticipated that if the crystal is nearly type II the Gaussian variant will work badly, because here $\bar{\sigma}_{i}\left(\Gamma_{i}\right)$ is practically given by $\sigma_{i}\left(\Gamma_{i}\right)$, roughly approximated by a Gaussian.

## 5. Application to spherical and ellipsoidal crystals

The diffraction data for precise structure determination are collected on spherical crystals, because this shape makes it relatively easy to make absorption and extinction corrections. For extinction, polyhedral crystals may be approximated by ellipsoids. If one denotes by $\rho_{10}$ and $\rho_{20}$ the ellipsoid radii along $\mathbf{i}_{1}$ and $\mathbf{i}_{2}$ and transforms the ellipsoid into a sphere of unit radius (when $\theta$ is transformed into $\theta_{0}^{\prime}$ ), (53), (56) become

$$
\begin{align*}
x= & Q\left(\alpha_{1} \alpha_{2} \rho_{10} \rho_{20}\right)^{1 / 2} \overline{t^{\prime}\left(a_{0}, a_{e}\right)}  \tag{64}\\
y_{0}= & {\left[3 / 4 \pi A\left(\zeta, a_{0}\right)\right] \int_{-\infty}^{\infty} \mathrm{d} \gamma \varphi(\gamma) } \\
& \times \int_{S_{1}} \mathrm{~d} V^{\prime} \exp \left[-\zeta\left(t_{1}^{\prime} a_{0}^{1 / 2}+t_{2}^{\prime} / a_{0}^{1 / 2}\right)\right] \tag{69}
\end{align*}
$$

$$
\begin{aligned}
y_{0}(x, \eta, \zeta, \chi)= & y_{00}(x, \eta)\left[1+\zeta A_{2} x^{1.5} /\left(1+B_{2} x^{1.5}\right)\right] \\
& \times\left[1-\frac{\chi^{2} A_{3} x^{1.5} /\left(1+B_{3} x^{1.5}\right)}{1+\chi^{2} C_{3} x^{1.5} /\left(1+D_{3} x^{1.5}\right)}\right]
\end{aligned}
$$

the coefficients $A_{2}, \ldots, D_{3}$ being polynomials in $\eta$ and $\zeta$. For the function $y_{00}(x, \eta)$ we have tried the formula proposed by Becker \& Coppens (1974a) and found that some values are reproduced with an error up to $10 \%$, though the $R$ factor (for the fit of $y_{00}$ only) is good. Then we modified the BC empirical formula as follows (for Gauss and Lorentz distributions):

$$
\begin{align*}
y_{00}^{(G)}(x, \eta)= & {\left[1+x\left(D_{1}+\frac{A_{1} x^{c_{1}}}{1+B_{1} x_{1} C_{1}}\right)\right]^{-1 / 2} ; \quad D_{1}=2^{1 / 2} }  \tag{70}\\
y_{00}^{(L)}(x, \eta)= & {\left[1+x\left(D_{1}+\frac{A_{1} x}{1+B_{1} x}\right.\right.} \\
& \left.\left.+\frac{C_{1} x^{2}}{\left(1+B_{1} x\right)^{2}}\right)\right]^{-1 / 2} ; \quad D_{1}=1, \tag{71}
\end{align*}
$$

the coefficients $A_{1}, B_{1}, C_{1}$ being polynomials in $\eta$. The expressions for the coefficients $A_{1}, \ldots, D_{3}$ obtained from the fit are
(a) for the Gauss mosaic:

$$
\begin{align*}
A_{1}= & 0.721+0.12 \eta-0.736 \eta^{2} \\
10 B_{1}= & 0.11+0.603 \eta-0.709 \eta^{2} \\
C_{1}= & 1-0.505 \eta^{2} \\
10 A_{2}= & 0.08+0.877 \eta-0.577 \eta^{2}-0.068 \eta \zeta \\
B_{2}= & 0.502-0.912 \eta-0.15 \zeta+0.734 \eta^{2} \\
& +0.138 \eta \zeta+0.021 \zeta^{2}  \tag{72}\\
10 A_{3}= & (0.124+0.082 \eta-0.031 \zeta) \eta^{2} \\
10 B_{3}= & 0.338+0.094 \eta-0.062 \zeta \\
10 C_{3}= & 0.585+0.233 \zeta+0.131 \eta \zeta-0.118 \zeta^{2} \\
D_{3}= & 0.66-1.414 \eta+0.9 \eta^{2} ;
\end{align*}
$$

(b) for the Lorentz mosaic:

$$
\begin{align*}
A_{1}= & 0.227+0.248 \eta-0.54 \eta^{2} \\
B_{1}= & 0.784-0.351 \eta \\
C_{1}= & 0.132-0.577 \eta+0.379 \eta^{2} \\
10 A_{2}= & 0.046+0.527 \eta-0.336 \eta^{2}-0.042 \eta \zeta \\
B_{2}= & 0.616-0.928 \eta-0.166 \zeta+0.713 \eta^{2} \\
& +0.147 \eta \zeta+0.023 \zeta^{2}  \tag{73}\\
10 A_{3}= & (0.114-0.017 \zeta) \eta^{2} \\
10 B_{3}= & 1.253-0.803 \eta-0.06 \zeta \\
10 C_{3}= & 0.668+0.221 \zeta+0.119 \eta \zeta-0.117 \zeta^{2} \\
D_{3}= & 0.819-1.629 \eta+0.975 \eta^{2} .
\end{align*}
$$

The factor $R$ was 0.0054 for the Gauss and 0.0027 for the Lorentz distribution. The empirical formula
reproduces (65) with an accuracy better than $1 \%$ if $y_{0}>0 \cdot 2$.

For $\zeta=0$ this empirical formula holds also for ellipsoidal crystals if $\eta$ is replaced by $\eta_{0}^{\prime}=\sin \theta_{0}^{\prime}$ and $\chi$ by $\chi_{e}=\ln a_{e}$. For $\zeta \neq 0$, in spite of its decreased accuracy, this formula still conserves two decimal digits in $y_{0}$ if $a_{0} \leq 5$.

The present theory (as Zachariasen's and BC) gives the primary extinction too, but this must be used with caution. To find it we must set $\zeta=0, \rho_{i 0}=\rho_{i}$ and $\varphi(\gamma)=\Phi_{1}(4 \pi \gamma / 3)$ in (64)-(67), where $\Phi_{1}$ is given by ( $B 15 b$ ). Since $a_{e}=1, y_{0}$ will be just $y_{00}$, depending on $\eta^{\prime}=\sin \theta^{\prime}$ and $x_{p}$ :

$$
\begin{equation*}
x_{p}=(9 / 4) Q \rho_{1} \rho_{2} \sin 2 \theta / \lambda . \tag{74}
\end{equation*}
$$

An excellent least-squares fit gives the formula (71) where one must take $D_{1}=1.25714$ (as above, $D_{1}$ is chosen to have the identity at $x=0$ between the first derivative of the exact and approximate $y_{0}$ ). As a result,

$$
\begin{align*}
& A_{1}=0.509+0.255 \eta^{\prime}-0.718 \eta^{\prime 2} \\
& B_{1}=0.139+0.436 \eta^{\prime}-0.166 \eta^{\prime 2}  \tag{75}\\
& C_{1}=-0.512-0.53 \eta^{\prime 2} .
\end{align*}
$$

The correction (54) must also be accounted for, with $\left(\delta_{1} \delta_{2}\right)_{p}=\sin ^{2} 2 \theta^{\prime}$.

## 6. Comparison with the Becker \& Coppens theory

As we have proved above, differences between the present and the BC theory appear only if the crystal is not of type $I$. There are two distinct situations: either the average mosaic block can or cannot be approximated by a sphere. In the first case only the second decimal digit in the extinction factor is changed, but significant qualitative and quantitative differences exist in the presence of the anisotropic block shape. In the latter case our extinction factor remains invariant to the permutation of the incident and diffracted directions, a property lost in the BC theory. This can readily be seen from (52). There the general term of the sum is proportional to $\alpha_{1}^{n} t^{(n)}$ which after inversion of the particle directions becomes $\alpha_{2}^{n} t^{(n)} / a^{n} \equiv \alpha_{1}^{n} t^{(n)}$ [where we have used $\binom{n}{m}=\binom{n}{n-m}$ and $v_{n m}=v_{n, n-m}$. In the BC theory $a$ is always unity and the last identity in general does not hold. Then the extinction factor is physically incorrect in BC theory if the crystal is of general type with an anisotropic block shape. This would have been of only academic importance if the quantitative differences were small. Fig. 4 speaks for itself. There $1 / y_{0}$ for a spherical crystal with $\mu=0, \theta=45^{\circ}$ and Lorentz mosaic distribution is shown versus the parameter $x$ from the BC theory. For this particular case the latter is related to the actual $x$ by $x($ actual $)=$ $x(\mathrm{BC})(1+a) / 2$. The curve for $a=1$ is just the BC
result. It is seen that even for moderate $a \neq 1$ the relative differences reach $40-50 \%$. For the Gaussian distribution the differences are even greater.

Therefore we can conclude that the results obtained in this paper are of practical importance and may explain some failure in the description of type II or mixed-type secondary extinction.

I am indebted to Dr A. M. Balagurov for encouragement and helpful discussions.

## APPENDIX A

## Becker-Coppens approximation

Here we want to clarify what the approximation used by Becker \& Coppens (1974a) consists of. This is necessary because there is an error in their paper: the second unlabelled formula on p. 142, as well as the subsequent arguments, are wrong. To arrive at the correct formula [their equation (B3)], the oblique hatched area $S\left(x_{1}, x_{2}\right)$ in Fig. 5(a) (partly reproducing Fig. 12 of BC ) must be replaced by the horizontally hatched area $S_{a}\left(x_{1}, x_{2}\right)$. As a consequence, there are four (not two) regions (Fig. $5 c$ ) where the approximation acts differently: if $\left(x_{1}, x_{2}\right)$ is in the region $A$, then $S=S_{a}$ and there is no approximation; in region $B, S<S_{a} ; S>S_{a}$ in $C$; in $D$ all three possibilities exist. Hence there are some compensation effects. On the other hand, it is not a priori evident that $y$ is always overestimated.

Let us use this approximation for our cases. In the first case we must reduce the integral

$$
\begin{align*}
J_{1}\left(x_{1}, x_{2}\right)= & \int_{S\left(x_{1}, x_{2}\right)} \mathrm{d} u_{1} \mathrm{~d} u_{2}\left(x_{1}-u_{1}\right)^{k}\left(x_{2}-u_{2}\right)^{l} \\
& \times \int_{S\left(u_{1}, u_{2}\right)} \mathrm{d} v_{1} \mathrm{~d} v_{2}\left(u_{1}-v_{1}\right)^{m}\left(u_{2}-v_{2}\right)^{n} \\
& \times t_{1}^{i}\left(v_{1}, v_{2}\right) g\left(v_{2}\right) . \tag{A1}
\end{align*}
$$

Replacing $S\left(u_{1}, u_{2}\right)$ by $S_{a}\left(u_{1}, u_{2}\right)$ (Fig. $5 b$ ) we can


Fig. 4. The inverse of $y_{0}$ versus $x(\mathrm{BC})$ for different values of the parameter $a$. The crystal is of spherical shape, $\theta=45^{\circ}, \mu=0$, Lorentz mosaic distribution.
write

$$
\begin{align*}
J_{1}\left(x_{1}, x_{2}\right)= & \int_{x_{1}^{0}}^{x_{1}} \mathrm{~d} u_{1} \int_{u_{2}^{0}}^{x_{2}} \mathrm{~d} u_{2}\left(x_{1}-u_{1}\right)^{k}\left(x_{2}-u_{2}\right)^{l} \\
& \times \int_{u_{2}^{0}}^{u_{2}} \mathrm{~d} v_{2} g\left(v_{2}\right)\left(u_{2}-v_{2}\right)^{n} \\
& \times \int_{v_{1}^{0}}^{u_{1}} \mathrm{~d} v_{1}\left(u_{1}-v_{1}\right)^{m}\left(v_{1}-v_{1}^{0}\right)^{i} . \tag{A2}
\end{align*}
$$

Performing the integral over $v_{1}$ and inverting the integrals over $u_{2}$ and $v_{2}$ one has

$$
\begin{align*}
J_{1}\left(x_{1}, x_{2}\right)= & \frac{i!m!}{(i+m+1)!} \int_{x_{1}^{1}}^{x_{1}} \mathrm{~d} u_{1}\left(x_{1}-u_{1}\right)^{k} \\
& \times \int_{u_{2}^{0}}^{x_{2}} \mathrm{~d} v_{2} g\left(v_{2}\right)\left(u_{1}-v_{1}^{0}\right)^{i+m+1} \\
& \times \int_{v_{2}}^{x_{2}} \mathrm{~d} u_{2}\left(x_{2}-u_{2}\right)^{l}\left(u_{2}-v_{2}\right)^{n} . \tag{A3}
\end{align*}
$$

Now the integral over $u_{2}$ is performed and after that by changing $u_{2}$ for $v_{2}$ (as an integration variable) one


Fig. 5. The BC approximation: the oblique hatched area is replaced by the horizontally hatched area $(a)$ in the original paper of BC (1974a) and in the present paper for the integral (A5), (b) in this paper for the integral (A1); (c) the regions where the BC approximation acts differently.
obtains

$$
\begin{align*}
J_{1}\left(x_{1}, x_{2}\right)= & \frac{l!m!n!i!}{(l+n+1)!(i+m+1)!} \\
& \times \int_{S\left(x_{1}, x_{2}\right)} \mathrm{d} u_{1} \mathrm{~d} u_{2} g\left(u_{2}\right) \\
& \times\left(x_{1}-u_{1}\right)^{k}\left(x_{2}-u_{2}\right)^{l+n+1} t_{1}^{i+m+1}\left(u_{1}, u_{2}\right) . \tag{A4}
\end{align*}
$$

In the second case we have to reduce the integral

$$
\begin{align*}
J_{2}= & \int_{S_{0}} \mathrm{~d} x_{1} \mathrm{~d} x_{2} f\left(x_{1}\right) t_{2}^{n}\left(x_{1}, x_{2}\right) \\
& \times \int_{S\left(x_{1}, x_{2}\right)} \mathrm{d} u_{1} \mathrm{~d} u_{2} g\left(u_{2}\right) \\
& \times t_{1}^{k}\left(u_{1}, u_{2}\right)\left(x_{1}-u_{1}\right)^{m}\left(x_{2}-u_{2}\right)^{l} \tag{A5}
\end{align*}
$$

Replacing $S\left(x_{1}, x_{2}\right)$ by $S_{a}\left(x_{1}, x_{2}\right)$ (Fig. $\left.5 a\right)$ and following exactly the same path as for $J_{1}$ one obtains

$$
\begin{align*}
J_{2}= & \frac{k!m!n!l!}{(m+k+1)!(n+l+1)!} \int_{S_{0}} \mathrm{~d} x_{1} \mathrm{~d} x_{2} f\left(x_{1}\right) g\left(x_{2}\right) \\
& \times t_{1}^{m+k+1}\left(x_{1}, x_{2}\right) t_{2}^{n+l+1}\left(x_{1}, x_{2}\right) \tag{A6}
\end{align*}
$$

## APPENDIX B

Kinematical cross sections for the perfect block of ellipsoidal shape
Let $r_{i}$ be the principal axes of the ellipsoid oriented along the unit vectors $c_{i}(i=1,2,3)$. If the ellipsoid is transformed into a sphere of unit radius (see Appendix $C$ ) and one denotes $\beta=\sum_{i} r_{i} h_{i} \mathbf{c}_{i}$, then the integral in (2) can be easily performed and for the most probable block one has

$$
\begin{align*}
\sigma\left(\Gamma_{1}, \Gamma_{2}, \Gamma_{3}\right) & =(4 \pi / 3) r_{1} r_{2} r_{3} n^{2}|F|^{2} \Phi_{3}(\beta)  \tag{B1a}\\
\Phi_{3}(\beta) & =9(\sin \beta-\beta \cos \beta)^{2} / \beta^{6} \tag{B1b}
\end{align*}
$$

If one takes account of (4), (38) and of the transformation (Fig. 1)

$$
\begin{array}{r}
\tau_{i}=\left(1-\delta_{i 3}\right)\left[(-1)^{i-1} \sin \theta \mathbf{n}_{1}+\cos \theta \mathbf{n}_{2}\right]+\delta_{i 3} \mathbf{n}_{3} \\
(i=1,2,3), \tag{B2}
\end{array}
$$

the modulus of $\beta$ may be expressed as

$$
\begin{equation*}
\beta=k_{0}\left(\Gamma^{*} \mathrm{~A} \Gamma^{*}\right)^{1 / 2} \tag{B3}
\end{equation*}
$$

where

$$
\begin{gather*}
\Gamma^{* \prime}=\left(-\Gamma_{1}, \Gamma_{2}, \Gamma_{3}\right), \quad \mathrm{A}=\mathrm{BPB}^{\prime}  \tag{B4a,b}\\
\mathrm{P}=\mathrm{E}^{(c)} \mathrm{R}^{2} \mathrm{E}^{(c)}, \quad \mathrm{R}_{i j}=r_{i} \delta_{i j} \tag{B5a,b}
\end{gather*}
$$

$\delta_{i j}$ being the Kronecker symbol and $B$ the matrix of the transformation ( $B 2$ ). Now ( $B 1$ ) is integrated over $\Gamma_{3}$ to obtain $\sigma\left(\Gamma_{1}, \Gamma_{2}\right)$. Although difficult, this integral can be performed exactly, but unfortunately the result cannot be convoluted analytically with $W\left(\varepsilon_{3}\right)$. In this
situation it is better to approximate $\Phi_{3}(\beta)$ in the beginning and then integrate over $\Gamma_{3}$. If $W\left(\varepsilon_{3}\right)$ is Gaussian, it is convenient to take

$$
\begin{equation*}
\Phi_{3}(\beta) \simeq \exp \left(-c_{0}^{2} \beta^{2} / 4 \pi\right) \tag{B6}
\end{equation*}
$$

where $c_{0}=1 \cdot 612 \simeq 5 / 3$ is determined from the normalization condition for $\sigma\left(\Gamma_{1}, \Gamma_{2}, \Gamma_{3}\right)$. Now, integrating over $\Gamma_{3}$ one has as a result a Gaussian with argument proportional to

$$
\begin{equation*}
\omega=\left[\sum_{i, j=1}^{2}\left(A_{i j} A_{33}-A_{i 3} A_{j 3}\right) \Gamma_{i}^{*} \Gamma_{j}^{*} / A_{33}\right]^{1 / 2} \tag{B7}
\end{equation*}
$$

By taking into account ( $B 4 b$ ) one has

$$
\begin{align*}
& A_{i j} A_{33}-A_{i 3} A_{j 3} \\
& =|\mathrm{P}|\left[\cos ^{2} \theta P_{11}^{-1}+(-1)^{1+\delta_{i j}} \sin ^{2} \theta P_{22}^{-1}\right. \\
& \left.\quad+(-1)^{i} \delta_{i j} \sin 2 \theta P_{12}^{-1}\right] . \tag{B8}
\end{align*}
$$

The inverse of P is trivially calculated from ( $B 5 a$ ); replacing it in (B8) and comparing with (40) and (41) we can write

$$
\begin{equation*}
A_{i j} A_{33}-A_{i 3} A_{j 3}=|\mathrm{P}| \cos \left[\left(1-\delta_{i j}\right) 2 \theta^{\prime}\right] /\left(\rho_{1} \rho_{2}\right) \tag{B9}
\end{equation*}
$$

On the other hand,

$$
\begin{align*}
A_{33}=P_{33} & =|\mathrm{P}|\left[P_{11}^{-1} P_{22}^{-1}-\left(P_{12}^{-1}\right)^{2}\right] \\
& =|\mathrm{P}| \sin ^{2} 2 \theta^{\prime} /\left(\rho_{1} \rho_{2} \sin 2 \theta\right)^{2} \tag{B10}
\end{align*}
$$

if (B8) and (B9) are considered. Finally, we find

$$
\begin{align*}
\omega= & \left(\rho_{2}^{2} \Gamma_{1}^{2}+\rho_{1}^{2} \Gamma_{2}^{2}-2 \cos 2 \theta^{\prime} \rho_{1} \rho_{2} \Gamma_{1} \Gamma_{2}\right)^{1 / 2} \\
& \times \sin 2 \theta / \sin 2 \theta^{\prime}  \tag{B11}\\
\sigma\left(\Gamma_{1}, \Gamma_{2}\right)= & \frac{4 \pi}{3 c_{0}} \frac{\sin 2 \theta}{\sin 2 \theta^{\prime}} n^{2}|F|^{2} \lambda \rho_{1} \rho_{2} \\
& \times \exp \left(-\pi c_{0}^{2} \omega^{2} / \lambda^{2}\right) \tag{B12}
\end{align*}
$$

The exact expression is (Popa, unpublished work)

$$
\begin{align*}
\sigma\left(\Gamma_{1}, \Gamma_{2}\right)= & \left(4 \pi \sin 2 \theta / 5 \sin 2 \theta^{\prime}\right) n^{2}|F|^{2} \lambda \\
& \times \rho_{1} \rho_{2} \Phi_{2}(2 \pi \omega / \lambda) \tag{B13a}
\end{align*}
$$

$$
\begin{align*}
\Phi_{2}(x)= & 15 \sum_{k=0}^{\infty}(-1)^{k} x^{2 k} \\
& \times[(2 k+3)(2 k+5) k!(k+1)!]^{-1} \tag{B13b}
\end{align*}
$$

Integrating ( $B 12$ ) over $\Gamma_{2}$, one gets

$$
\begin{align*}
\sigma\left(\Gamma_{1}\right)= & Q\left(c_{0} / \lambda\right) \rho_{2} \sin 2 \theta \\
& \times \exp \left[-\pi\left(c_{0}^{2} / \lambda^{2}\right) \rho_{2}^{2} \sin ^{2} 2 \theta \Gamma_{1}^{2}\right] \tag{B14}
\end{align*}
$$

the exact expression being (BC)
$\sigma\left(\Gamma_{1}\right)=Q(3 / 2 \lambda) \rho_{2} \sin 2 \theta \Phi_{1}\left[(2 \pi / \lambda) \rho_{2} \sin 2 \theta \Gamma_{1}\right]$

$$
\begin{equation*}
\Phi_{1}(x)=\left(x^{2}-x \sin 2 x+\sin ^{2} x\right) / x^{4} \tag{B15a}
\end{equation*}
$$

The Lorentzian approximation of ( $B 15 b$ ) needed in $\S 4$ is

$$
\begin{equation*}
\Phi_{1}(x) \simeq\left(1+9 x^{2} / 16\right)^{-1} \tag{B16}
\end{equation*}
$$

For $\sigma_{2}\left(\Gamma_{2}\right)$ the same expressions are obtained with $\rho_{1}$ in place of $\rho_{2}$.

## APPENDIX $C$ <br> Transformation of the ellipsoid into a sphere of unit radius

Remember here some mathematical relations used above (see also Becker \& Coppens, 1975). If the equation of an ellipsoidal surface in the system ( $\mathbf{c}_{\boldsymbol{i}}$ ) of its principal axes is $\sum_{i} z_{i}^{2} / r_{i}^{2}=1$ and the transformation $z_{i}=r_{i} z_{i}^{\prime}$ is performed then this equation become $\sum_{i} z_{i}^{\prime 2}=1$, which represents a sphere of unit radius. By this transformation any unit vector $\mathbf{u}=\sum_{i} u_{i} \mathbf{c}_{i}$ is transformed into the vector $\mathbf{U}^{\prime}=\sum_{i} U_{i}^{\prime} \mathbf{c}_{i}=\sum_{i} u_{i} \mathbf{c}_{i} / r_{i}$. If one denotes by $\rho_{u}$ the ellipsoid radius along the vector $\mathbf{u}$, then the vector $\rho_{u} \mathbf{u}$ is transformed into $\mathbf{u}^{\prime}=\rho_{u} \mathbf{U}^{\prime}$ of unit length. In consequence we can write

$$
\begin{equation*}
1 / \rho_{u}^{2}=\sum_{i} u_{i}^{2} / r_{i}^{2} \tag{C1}
\end{equation*}
$$

Now, if $\mathbf{u}$ and $\mathbf{v}$ are a pair of unit vectors whose mutual angle is $\varphi$, after transformation this angle becomes

$$
\begin{equation*}
\cos \varphi^{\prime}=\mathbf{u}^{\prime} \cdot \mathbf{v}^{\prime}=\rho_{u} \rho_{v} \sum_{i} u_{i} v_{i} / r_{i}^{2} \tag{C2}
\end{equation*}
$$

Finally, if $t$ is a segment in the $\mathbf{u}$ direction, the vector $t \mathbf{u}$ is transformed into $t \mathbf{U}^{\prime}=t\left|\mathbf{U}^{\prime}\right| \mathbf{u}^{\prime}=t^{\prime} \mathbf{u}^{\prime}$. Then the transformation $t^{\prime}$ of $t$ is

$$
\begin{equation*}
t^{\prime}=t / \rho_{u} \tag{C3}
\end{equation*}
$$

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# X-ray Diffraction by a Low-Angle Twist Boundary Perpendicular to Crystal Surface. I. Superstructure Factor of Screw Dislocation Superlattice 

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#### Abstract

The X-ray two-wave diffraction on a dislocation wall perpendicular to a crystal surface, consisting of periodically arranged dislocations (low-angle twist boundary), is considered in the case when the dislocation superlattice period is much less than the crystal extinction length. The formula obtained for the reflected intensity is of the same form as that for an ideal crystal with a modified crystal structure factor. The superstructure factor of a dislocation superlattice is


calculated. The recurrence relations are produced which enable a superstructure factor to be calculated for a satellite of any order and magnitude hb (h is the diffraction vector, $b$ is the Burgers vector).

## 1. Introduction

The grain boundary (GB) is a surface between two misorientated single crystals. The dislocation structure of a GB is well known (Hirth \& Lothe, 1968;


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